# Stochastic Riemannian Geometry 

Anton Thalmaier<br>Department of Mathematics<br>University of Luxembourg<br>Campus Belval - Maison du Nombre Maison<br>L-4364 Esch-Sur-AlZETte<br>LUXEMBOURG<br>Email address: anton.thalmaier@uni.lu<br>URL: http://math.uni.lu/thalmaier

[^0]
## Contents

Preface ..... v
Chapter 1. Stochastic Analysis on Manifolds ..... 1
1.1. Stochastic Flows ..... 11
1.2. Construction of Stochastic Flows ..... 21
1.3. Quadratic Variation and Integration of one-forms ..... 35
1.4. Linear Connections and Martingales on Manifolds ..... 43
1.5. Riemannian Metrics and Brownian Motions ..... 54
1.6. Parallel Transport and Stochastically Moving Frames ..... 63
1.7. Morphisms of Martingales and Brownian Motions ..... 87
1.8. Convergence and Confluence of Martingales ..... 100
1.9. Stochastic Differentials and Second Order Tangent Spaces ..... 111
Chapter 2. Geometry of Brownian Motion ..... 123
2.1. The Curvature Tensor and Jacobi Fields ..... 123
2.2. Brownian Motion and Curvature ..... 156
2.3. Heat Equation on Sections of a Vector Bundle ..... 176
2.4. Brownian Bridges and Brownian Loops ..... 186
Appendix A. Background on SDEs ..... 191
A.1. One-dimensional Stochastic Differential Equations ..... 191
A.2. Derivative Flows ..... 200
Bibliography ..... 203
Index ..... 207
Notations ..... 213

## Preface

The purpose of these notes is to develop fundamental tools of Stochastic Analysis on differentiable manifolds, and to provide a unified and comprehensive introduction to stochastic methods in Riemannian geometry. Right from the beginning, these objectives demand to carry over classical notions from Stochastic Analysis on Euclidean space to general manifolds and to develop the necessary concepts in a coordinate-free manner.

One of the immediate obstacles of Stochastic Analysis on manifolds is related to the fact that, in general, it is not feasible to transfer processes via charts from $\mathbb{R}^{n}$ to curved spaces, and to deal appropriately with certain classes of manifold-valued processes in terms of local coordinates. Itô's formula for $\mathbb{R}^{n}$-valued semimartingales shows that concepts like Brownian motions or local martingales are not invariant under coordinate transformations.

It is an elementary observation based on Itô's formula which leads to an intrinsic notion of manifold-valued semimartingales. However it turns out that martingale theory, traditionally based on the linear concept of conditional expectations, requires on manifolds an additional geometric structure such as a linear connection in the tangent bundle.

In the situation of Riemannian manifolds there is a canonical linear connection linked to the Riemannian geometry of the manifold, the so-called Levi-Civita connection, but for various purposes it is desirable to work also with more general linear connections. We develop martingale theory on general manifolds endowed with a linear connection.

Brownian motion on a Riemannian manifold is the special case of a martingale related to the Levi-Civita connection. Brownian motions are associated to the Riemannian metric via the Laplace-Beltrami operator and generalize the class of standard $\mathbb{R}^{n}$-valued Brownian motions. By definition, Brownian motions are local objects in the sense that for small times their behaviour is controlled by local geometry. However, their large-scale probabilistic behaviour reflects global aspects of topology and geometry of the manifold. Brownian motion picks up global invariants of the manifold, in their behaviour as time goes to infinity, and allows to interpolate between local and global geometry.

## CHAPTER 1

## Stochastic Analysis on Manifolds

In this chapter we deal with the theory of continuous manifold-valued semimartingales and develop fundamental tools about diffusions, martingales and Brownian motions on differentiable manifolds.

We start with a brief review of basic concepts from differential topology, mainly to fix the notions for further reference. For more details and additional information the reader is referred to standard textbooks (e.g., [15] or [45]).

DEFINITION 1.0.1 (Topological manifold). A Hausdorff topological space $M$ endowed with a countable basis for the topology is called $n$-dimensional topological manifold, if for every point $x \in M$ there is an open neighbourhood $U$ of $x$ in $M$ and a homeomorphism $h: U \rightarrow U^{\prime}$ onto an open subset $U^{\prime} \subset \mathbb{R}^{n}$.

DEFINITION 1.0.2 (Chart). A homeomorphism $h: U \rightarrow U^{\prime}$ from some open subset $U \subset M$ onto an open subset $U^{\prime} \subset \mathbb{R}^{n}$ is called a ( $n$-dimensional) chart for $M$. Charts are denoted by $(h, U)$. A chart for $M$ is said to be a chart about $x \in M$ if $x \in U$.

DEFINITION 1.0.3 (Transition map). Let $(h, U)$ and $(k, V)$ be charts for $M$. The homeomorphism

$$
k \circ h^{-1} \mid h(U \cap V): h(U \cap V) \rightarrow k(U \cap V)
$$ is called transition map from $(h, U)$ to $(k, V)$.



Figure 1.0.1. Transition map from $(h, U)$ to $(k, V)$
DEFINITION 1.0.4 (Atlas, differentiable structure). A family $\left(h_{i}, U_{i}\right)_{i \in I}$ of $n$-dimensional charts for $M$ is called atlas for $M$ if the $U_{i}$ cover $M$. An atlas is said to be differentiable if all its transition maps are differentiable (i.e., $C^{\infty}$ ). A maximal differentiable atlas for $M$ is called a $n$-dimensional differentiable structure for $M$.

If $\mathfrak{A}$ denotes a differentiable atlas for $M$ and if $(h, U),(k, V)$ are two additional $n$ dimensional charts with smooth transition maps to all charts of $\mathfrak{A}$, then they also change smoothly between each other. In particular,

$$
\mathscr{D}(\mathcal{A}):=\{(h, U) n \text {-dimensional chart for } M \mid(h, U) \text { changes smoothly with } \mathcal{A}\}
$$

defines an $n$-dimensional differentiable structure for $M$.
DEFINITION 1.0.5 (Differentiable manifold). A $n$-dimensional differentiable manifold is a pair $(M, \mathscr{D})$ where $M$ is a topological Hausdorff space with a countable basis for the topology and $\mathscr{D}$ an $n$-dimensional differentiable structure for $M$.

In the sequel we deal with differentiable manifolds only; the addition "differentiable" or "smooth" is mostly omitted. Furthermore, the differentiable structure $\mathscr{D}$ of a manifold $(M, \mathscr{D})$ is suppressed in the notation; one writes simply $M$ and refers to the charts of $\mathscr{D}$ also as charts of $M$. By convention, the empty topological space is assumed to be a manifold of arbitrary (also negative) dimension; the (well-defined) dimension of nonempty manifolds is denoted $\operatorname{dim} M$.

Example 1.0.6. The direct product $M \times N$ of two manifolds $M$ and $N$ is canonically a manifold of dimension $\operatorname{dim} M+\operatorname{dim} N$ (products of charts define a suitable atlas and the required differentiable structure).

DEFINITION 1.0.7 (Submanifold). Let $M$ be an $n$-dimensional manifold and $0 \leq$ $k \leq n$. A subspace $M_{0} \subset M$ is said to be a $k$-dimensional (or $(n-k)$-codimensional) submanifold of $M$, if about every point in $M_{0}$ there exists a chart $(h, U)$ for $M$ such that $h\left(U \cap M_{0}\right)=h(U) \cap\left(\mathbb{R}^{k} \times\{0\}\right)$. The subspace $M_{0}$ then itself is a $k$-dimensional manifold in the obvious way.


Figure 1.0.2. Submanifold of $M$
DEFINITION 1.0.8 (Differentiable map). Let $f: M \rightarrow N$ be a continuous map between manifolds and $x \in M$. The map $f$ is said to be differentiable at $x$, if for one (and then every) chart $(h, U)$ at $x$ and for one (and then every) chart $(k, V)$ at $f(x)$, the "pushed down" mapping $k \circ f \circ h^{-1}$ (defined on $h\left(f^{-1}(V) \cap U\right)$ ) is differentiable at $h(x)$.

Analogously, $f$ is said to be $k$-times differentiable at $x$ (where $k \in \mathbb{N}$ ), resp., infinitely often differentiable at $x$, if the same property holds true for the pushed down mapping
at $h(x)$. In the case when $f$ is $k$-times, resp. infinitely often differentiable at any point $x \in M$, we write $f \in C^{k}(M ; N)$ and $f \in C^{\infty}(M ; N)$, respectively. The expression " $f$ is differentiable" or "smooth" always means $f \in C^{\infty}(M ; N)$, i.e. "infinitely often differentiable".


Figure 1.0.3. Maps in local coordinates
Finally, $f$ is said to be a diffeomorphism, if $f$ is bijective and both $f$ as well as $f^{-1}$ are differentiable.

The space of real-valued differentiable functions on $M$ is denoted $C^{\infty}(M)$. Differentiable functions on $M$ of compact support are called test functions on $M$; the space of test functions on $M$ is denoted $C_{c}^{\infty}(M)$.

Note that the derivative of the pushed down map at $h(x)$, expressed in terms of the Jacobian $J_{h(x)}\left(k \circ f \circ h^{-1}\right)$, depends on the specific choice of charts, whereas the rank of the derivative at $h(x)$, denoted $\operatorname{rank}_{x} f$, is independent of coordinates.

A useful fact (inverse function theorem) is that a differentiable map $f: M \rightarrow N$ between manifolds of equal dimension $n$ is a local diffeomorphism at $x$ (i.e., a diffeomorphism of an open neighbourhood of $x$ onto some open neighbourhood of $f(x)$ ) if and only if $\operatorname{rank}_{x} f=n$.

For the construction of a chart independent version of the differential (as a canonical linearization of the differentiable map $f: M \rightarrow N$ locally at $x$ ) it is suitable to approximate the manifold $M$ at $x$ itself by a linear object, i.e. the tangent space $T_{x} M$.

DEFINITION 1.0.9 (Tangent space; geometric definition). Let $M$ be a manifold and $x \in M$. Let

$$
K_{x} M:=\{\alpha:]-\varepsilon, \varepsilon[\rightarrow M \text { differentiable } \mid \varepsilon>0, \alpha(0)=x\}
$$

denote the set of differentiable curves $\alpha$ through $x$. Two curves $\alpha, \beta \in K_{x} M$ are called tangentially equivalent, written $\alpha \sim \beta$, if $(h \circ \alpha)^{\cdot}(0)=(h \circ \beta)^{\cdot}(0)$ for one (and then any) chart $(h, U)$ at $x$. The quotient $\left(T_{x} M\right)_{\text {geom }}:=K_{x} M / \sim$ is called the (geometric) tangent space of $M$ in $x$, and the classes $[\alpha] \in\left(T_{x} M\right)_{\text {geom }}$ are called the (geometric) tangent vectors of $M$ in the point $x$.

Note that by definition $\alpha \sim \alpha_{h}$ for $\alpha \in K_{x} M$ and $(h, U)$ a chart about $x$, where $\alpha_{h}(t):=h^{-1}\left(h(x)+t(h \circ \alpha)^{\cdot}(0)\right)$ for $t$ sufficiently small.

The fact that $\left(T_{x} M\right)_{\text {geom }}$ is a finite-dimensional real vector space is not obvious from the given definition. It becomes however evident by adopting a slightly different point of view. First of all, two real differentiable functions defined locally about $x$ are called equivalent, if they coincide on some neighbourhood of $x$. The resulting equivalence classes are called germs of differentiable functions at $x$. The set $\mathcal{E}_{x} M$ of these germs inherits the structure of a real algebra in a natural way. In the notation it is usually not distinguished between a germ $\varphi \in \mathcal{E}_{x} M$ and its representative (a differentiable function defined locally about $x$ ).

The scalar multiplication $\varphi a=\varphi(x) a$ for $\varphi \in \mathcal{E}_{x} M, a \in \mathbb{R}$, gives $\mathbb{R}$ the structure of an $\mathcal{E}_{x} M$-module. An $\mathbb{R}$-derivation of $\mathcal{E}_{x} M$ in $\mathbb{R}$ is an $\mathbb{R}$-linear map $v: \mathcal{E}_{x} M \rightarrow \mathbb{R}$ satisfying the product rule

$$
v(\varphi \psi)=\varphi v(\psi)+\psi v(\varphi) \quad \text { for } \varphi, \psi \in \mathcal{E}_{x} M
$$

The set $\operatorname{Der}_{\mathbb{R}}\left(\mathcal{E}_{x} M, \mathbb{R}\right)$ of $\mathbb{R}$-derivations of $\mathcal{E}_{x} M$ in $\mathbb{R}$ forms naturally an $\mathcal{E}_{x} M$-module, and in particular a real vector space.

DEFINITION 1.0.10 (Tangent space; algebraic definition). Let $M$ be a manifold and $x \in M$. The real vector space

$$
\left(T_{x} M\right)_{\mathrm{alg}}:=\operatorname{Der}_{\mathbb{R}}\left(\mathcal{E}_{x} M, \mathbb{R}\right)
$$

is called the (algebraic) tangent space of $M$ at $x$, and $\mathbb{R}$-derivations $v \in\left(T_{x} M\right)_{\text {alg }}$ are called (algebraic) tangential vectors of $M$ at the point $x$.

REMARK 1.0.11. Let $M$ be manifold. For any $x \in M$ the spaces $\left(T_{x} M\right)_{\text {geom }}$ and $\left(T_{x} M\right)_{\text {alg }}$ are canonically identified; more precisely the following maps are inverse to each other:

$$
\begin{array}{ll}
\left(T_{x} M\right)_{\text {geom }} \rightarrow\left(T_{x} M\right)_{\mathrm{alg}}, & {[\alpha] \mapsto\left(\mathcal{E}_{x} M \rightarrow \mathbb{R}, \varphi \mapsto(\varphi \circ \alpha)^{\cdot}(0)\right)} \\
\left(T_{x} M\right)_{\mathrm{alg}} \rightarrow\left(T_{x} M\right)_{\text {geom }}, & v \mapsto[]-\varepsilon, \varepsilon\left[\rightarrow M, t \mapsto h^{-1}(h(x)+t v(h))\right]
\end{array}
$$

where $(h, U)$ is a chart for $M$ at $x$ and $v(h):=\left(v\left(h^{1}\right), \ldots, v\left(h^{n}\right)\right) \in \mathbb{R}^{n}$.
DEFINITION 1.0.12 (Tangent space). Let $M$ be a manifold and $x \in M$. The real vector space $T_{x} M:=\left(T_{x} M\right)_{\text {alg }} \equiv\left(T_{x} M\right)_{\text {geom }}$ is called the tangent space of $M$ at $x$, its elements (considered either as derivations or represented by curves) are the tangent vectors of $M$ at the point $x$.

EXAMPLE 1.0 .13 . Any $n$-dimensional real vector space $V$ is a $n$-dimensional manifold in a canonical way. Furthermore, for $x \in V$, we have $T_{x} V \cong V$ canonically (as real vector spaces). Indeed, if $h: V \xrightarrow{\sim} \mathbb{R}^{n}$ is an isomorphism of vector spaces (and hence a global chart), then the homomorphisms

$$
\begin{array}{ll}
\left(T_{x} V\right)_{\mathrm{alg}} \rightarrow V, & v \mapsto h^{-1}\left(v\left(h^{1}\right), \ldots, v\left(h^{n}\right)\right) \\
V \rightarrow\left(T_{x} V\right)_{\mathrm{alg}}, & v \mapsto\left(\mathcal{E}_{x} V \rightarrow \mathbb{R},\left.\varphi \mapsto \frac{d}{d t} \varphi(x+t v)\right|_{t=0}\right)
\end{array}
$$

are inverse to each other and independent of the particular choice of $h$.
Definition 1.0.14 (Differential). Let $f: M \rightarrow N$ be a differentiable map between manifolds and $x \in M$. The differential of $f$ at $x$

$$
d f_{x} \equiv f_{* x}: T_{x} M \rightarrow T_{f(x)} N
$$

is respectively geometrically or algebraically explained as

$$
\left(d f_{x}\right)_{\text {geom }}:\left(T_{x} M\right)_{\text {geom }} \rightarrow\left(T_{f(x)} N\right)_{\text {geom }}, \quad[\alpha] \mapsto[f \circ \alpha]
$$

$$
\left(d f_{x}\right)_{\mathrm{alg}}:\left(T_{x} M\right)_{\mathrm{alg}} \rightarrow\left(T_{f(x)} N\right)_{\mathrm{alg}}, \quad v \mapsto\left(\mathcal{E}_{f(x)} N \rightarrow \mathbb{R}, \varphi \mapsto v(\varphi \circ f)\right)
$$

Both mappings are well-defined and consistent with respect to the canonical identification of geometric and algebraic tangent space.

REMARK 1.0.15 (Functorality of the differential). We have $d\left(\mathrm{id}_{M}\right)_{x}=\mathrm{id}_{T_{x} M}$ for $x \in M$. Further, for any differentiable maps $f: M \rightarrow N$ and $g: N \rightarrow L$ between manifolds, the chain rule $d(g \circ f)_{x}=d g_{f(x)} \circ d f_{x}$ holds. In particular, if $f$ is a local diffeomorphism at $x$, then $d f_{x}: T_{x} M \rightarrow T_{f(x)} N$ is a linear isomorphism.

The definitions of tangent spaces and differentials are obviously of local nature; for instance, let $U \subset M$ be open and $x \in U$, then $T_{x} M \cong T_{x} U$ in a canonical (and trivial) way, namely via $(d \iota)_{x}$ where $\iota: U \hookrightarrow M$ denotes the inclusion, and one identifies the tangent spaces $T_{x} M$ and $T_{x} U$.

Example 1.0.16 (Basis for $\left.T_{x} M\right)$. Let $M$ be an $n$-dimensional manifold and $(h, U)$ be a chart about $x \in M$. Then

$$
d h_{x}: T_{x} M \xrightarrow{\sim} T_{h(x)} \mathbb{R}^{n} \cong \mathbb{R}^{n}, \quad v \mapsto\left(v\left(h^{1}\right), \ldots, v\left(h^{n}\right)\right),
$$

is an isomorphism of real vector spaces, in particular, $\operatorname{dim}_{\mathbb{R}} T_{x} M=n$. Thus, by means of

$$
\left(\frac{\partial}{\partial h^{i}}\right)_{x}:=\left(d h_{x}\right)^{-1}\left(e_{i}\right)=d\left(h^{-1}\right)_{h(x)}\left(e_{i}\right), \quad i=1, \ldots, n,
$$

an $\mathbb{R}$-basis for $T_{x} M$ is given; note that here $\partial_{i, x} \equiv\left(\frac{\partial}{\partial h^{i}}\right)_{x} \in T_{x} M$ represents the derivation $\varphi \mapsto \frac{\partial}{\partial x^{i}}\left(\varphi \circ h^{-1}\right)(h(x))$.

THEOREM 1.0.17 (Differentials in coordinates). Let $M$ be an n-dimensional manifold, $N$ an n-dimensional manifold, $f: M \rightarrow N$ a differentiable map and $x \in M$. Choosing charts $(h, U)$ for $M$ about $x$ and $(k, V)$ for $N$ about $f(x)$, the following diagram commutes:

where $J_{h(x)}\left(k \circ f \circ h^{-1}\right) \in \mathbf{M}(n \times n ; \mathbb{R})$ is the Jacobian of $k \circ f \circ h^{-1}$ at $h(x)$.
Proof. Any $v \in T_{x} M$ can be written as $v=\sum_{i} v^{i}\left(\frac{\partial}{\partial h^{i}}\right)_{x}$ where $v^{i}=v\left(h^{i}\right)$. Upon Definition 1.0.14, the differential $(d f)_{x} v \in T_{f(x)} N$ is represented by the derivation $\varphi \mapsto$ $v(\varphi \circ f)$, so that $(d f)_{x} v=\sum_{j} v\left(k^{j} \circ f\right)\left(\frac{\partial}{\partial k^{j}}\right)_{f(x)}$. Thus, if $v=\sum_{i} v^{i}\left(\frac{\partial}{\partial h^{i}}\right)_{x}$, then

$$
v\left(k^{j} \circ f\right)=\sum_{i} v^{i}\left(\frac{\partial}{\partial h^{i}}\right)_{x}\left(k^{j} \circ f\right)=\sum_{i} \frac{\partial\left(k^{j} \circ f \circ h^{-1}\right)}{\partial x^{i}}(h(x)) v^{i},
$$

which shows the claim.
The examples above show that $\operatorname{rank}_{x} f=\operatorname{rank}\left(d f_{x}\right)$ for a differentiable map $f: M \rightarrow$ $N$ between manifolds and $x \in M$. In particular, if $d f_{x}$ is an isomorphism, then necessarily $\operatorname{dim} M=\operatorname{dim} N$ and $f$ is a local diffeomorphism at $x$ by the local inverse theorem.

DEFINITION 1.0.18 (Immersion, embedding). A map $f: M \rightarrow N$ between manifolds is called an immersion, if $f$ is differentiable and the linear map $d f_{x}: T_{x} M \rightarrow T_{f(x)} N$ is injective for any $x \in M$. A map $f: M \rightarrow N$ is called an embedding, if $f(M) \subset N$ is a submanifold and $f: M \rightarrow f(M)$ a diffeomorphism.

Any embedding is obviously an immersion. Immersions however are not injective in general; even an injective immersion is not necessarily an embedding.

DEFINITION 1.0.19 (Velocity of a curve). Let $\alpha: I \rightarrow M$ be a curve in $M$, defined on an open real interval $I=] a, b\left[\right.$, and let $d \alpha_{t}: \mathbb{R} \cong T_{t} I \rightarrow T_{\alpha(t)} M$ be the differential of $\alpha$ at $t \in I$. The vector $\dot{\alpha}(t):=d \alpha_{t}(1) \in T_{\alpha(t)} M$ is called velocity of $\alpha$ at $t$; algebraically it is the derivation $\varphi \mapsto(\varphi \circ \alpha)^{\circ}(t)$, geometrically $\dot{\alpha}(t)$ is represented by $s \mapsto \alpha(t+s)$. Obviously any equivalence class $[\alpha] \in\left(T_{x} M\right)_{\text {geom }}$ can be written as $\dot{\alpha}(0)$.

DEFINITION 1.0.20 (Locally trivial fibration, fiber bundle). Let $E, M$ and $F$ be manifolds. A differentiable map $\pi: E \rightarrow M$ is called a locally trivial fibration with typical fiber $F$ (or a fiber bundle), if about any point of $M$ there exists an open neighborhood $U$ and a diffeomorphism $\varphi: \pi^{-1}(U) \xrightarrow{\sim} U \times F$ above $U$, i.e. such that the following diagram commutes:


The pair $(\varphi, U)$ is said to be a bundle chart (or local trivialization) of the fibration, a family $\left(\varphi_{i}, U_{i}\right)_{i \in I}$ of bundle charts with $M=\bigcup_{i \in I} U_{i}$ is said to be a bundle atlas for $E$. In this situation, $M$ is called the basis, $E$ the total space, $\pi$ the projection, $F$ the typical fiber and $E_{x} \equiv \pi^{-1}(\{x\})$ the fiber at $x \in M$.

One commonly writes $E / M$ or just $E$ instead of $\pi: E \rightarrow M$ or $E \rightarrow M$. Furthermore, if $M_{0} \subset M$, we use occasionally the notation $E / M_{0}:=\pi^{-1}\left(M_{0}\right)$.

It is an immediate consequence of the definition that each fiber $E_{x}$ is a submanifold of $E$ diffeomorphic to the typical fiber $F$. For any two bundle charts $\left(\varphi_{i}, U_{i}\right),\left(\varphi_{j}, U_{j}\right)$ the composition $\varphi_{j} \circ \varphi_{i}^{-1}$ defines a diffeomorphism on $\left(U_{i} \cap U_{j}\right) \times F$; the corresponding maps $\phi_{i j}: U_{i} \cap U_{j} \rightarrow \operatorname{Diff}(F)$ into the group of diffeomorphisms of the typical fiber $F$ are called transition functions.

A bundle atlas provides automatically a differentiable atlas for the manifold $E$ and determines in that way the differentiable structure of $E$. Moreover in a canonical way, for $U \subset M$ open, $\pi \mid U: E / U \rightarrow U$ is a fiber bundle as well.

Definition 1.0.21 (Trivial fiber bundle). A fiber bundle $\pi: E \rightarrow M$ is said to be trivial, if there exists a trivialization, i.e. a global bundle chart $(\varphi, M)$.

DEFINITION 1.0 .22 (Vector bundle). A locally trivial fibration $\pi: E \rightarrow M$ with an $m$-dimensional real vector space $F$ as typical fiber is said to be an $m$-dimensional vector bundle over $M$, if there exists a bundle atlas $\left(\varphi_{i}, U_{i}\right)_{i \in I}$ for $E$ such that the diffeomorphisms

$$
\varphi_{j} \circ \varphi_{i}^{-1}:\{x\} \times F \rightarrow\{x\} \times F, \quad x \in U_{i} \cap U_{j}
$$

are linear isomorphisms of $F$.
Each fiber $E_{x}$ then carries the structure of a real vector space such that the bundle charts above are fiberwise linear: for $x \in U$ the restriction $\varphi \mid E_{x}$ maps the fiber $E_{x}$ linearly to $\{x\} \times F$. Without restriction of generality one can take $F=\mathbb{R}^{n}$ as the typical fiber.

DEFINITION 1.0.23 (Subbundle). Let $\pi: E \rightarrow M$ an $m$-dimensional vector bundle and $k \leq m$. A subset $E_{0} \subset E$ is said to be an $k$-dimensional subbundle of $E$, if about each point $x \in M$ there exists a fiberwise linear bundle chart $\varphi: E / U \rightarrow U \times \mathbb{R}^{m}$ for $E$ such that $\varphi\left(E_{0} / U\right)=U \times\left(\mathbb{R}^{k} \times\{0\}\right)$. Then $\pi \mid E_{0}: E_{0} \rightarrow M$ itself is a $k$-dimensional vector bundle over $M$.

DEFINITION 1.0.24 (Bundle homomorphism). Let $E$ and $E^{\prime}$ be vector bundles over the same base manifold $M$. A differentiable map $\phi: E \rightarrow E^{\prime}$ is called bundle homomorphism or homomorphism of vector bundles, if $\phi$ is a map over $M$ and linear in each fiber, i.e., if $\phi$ maps each $E_{x}$ to $E_{x}^{\prime}$ and each $\phi_{x}=\phi \mid E_{x}: E_{x} \rightarrow E_{x}^{\prime}$ is given by a linear map. This constitutes the category $\mathscr{V}_{M}$ of vector bundles over $M$.

DEFInition 1.0.25 (Section). Let $\pi: E \rightarrow M$ be a vector bundle. A section of $E$ is a differentiable map $A: M \rightarrow E$ such that $\pi \circ A=\operatorname{id}_{M}$ (i.e. a right inverse to $\pi$ ). The set $\Gamma(E)$ of sections of $E$ constitutes a $C^{\infty}(M)$-module in a natural way via $(\varphi A)(x)=$ $\varphi(x) A(x), \varphi \in C^{\infty}(M)$. The value of a section $A$ at $x \in M$ is also denoted $A_{x}$ instead of $A(x)$.

REMARK 1.0.26 (Local frame). Let $E \rightarrow M$ be an $m$-dimensional vector bundle and $x_{0} \in M$. A local frame for $E$ at $x_{0}$ consists of an open neighbourhood $U$ of $x_{0}$, together with sections $e_{1}, \ldots, e_{m} \in \Gamma(E / U)$ such that for any $x \in U$ the family $\left(e_{1}(x), \ldots, e_{m}(x)\right)$ provides an $\mathbb{R}$-basis of $E_{x}$. By means of appropriate bundle charts it is possible to construct local frames for $E$ at any $x_{0} \in M$. Then to each section $A \in \Gamma(E)$ there exist uniquely determined functions $a^{i} \in C^{\infty}(U)$ such that $A \mid U=\sum_{i=1}^{m} a^{i} e_{i}$.

When constructing fibrations one often starts with a basis $M$, a typical fiber $F$ and a family $\left(E_{x}\right)_{x \in M}$ of manifolds $E_{x}$ diffeomorphic to $F$. Then $E:=\cup_{x \in M} E_{x}\left(\equiv \bigcup_{x}\{x\} \times\right.$ $E_{x}$ ) and $\pi: E \rightarrow M, E_{x} \ni e \mapsto x$ gives the total space $E$, at first just as a set with the corresponding projection. The still missing topology and differentiable structure on $E$, as well as appropriate bundle charts, are then typically provided by canonical pre-bundle charts: A pre-bundle chart of $E$ is a pair $(\varphi, U)$, consisting of an open subset $U \subset M$ and a fiberwise diffeomorphic bijection $\varphi: E / U=\bigcup_{x \in U} E_{x} \rightarrow U \times F$ over $U$. A family $\left(\varphi_{i}, U_{i}\right)_{i \in I}$ of pre-bundle charts such that $\bigcup_{i \in I} U_{i}=M$ is called a pre-bundle atlas for $E$, if all transition maps

are differentiable, and thus diffeomorphisms.
LEMMA 1.0.27. To each pre-bundle atlas for $E$ there exists precisely one topology and differentiable structure on $E$ which make $\pi: E \rightarrow M$ a locally trivial fibration with typical fiber $F$ and the pre-bundle atlas to a bundle atlas.

Proof. Let $e \in E$ and $x:=\pi(e) \in M$. Via a pre-bundle chart $(\varphi, U)$ with $x \in U$ we have $\varphi: E / U \xrightarrow{\sim} U \times F$ where $e$ is mapped to some point $(x, v) \in U \times F$. A basis of neighbourhoods at $e \in E$ for the wanted topology on $E$ is found by pulling back via $\varphi$ a basis of open sets at $(x, v)$ in $U \times F$. The remaining claimed properties are then easily checked.

Example 1.0.28 (Tangent bundle). Let $M$ be an $n$-dimensional manifold. The tangent spaces $T_{p} M, p \in M$ are isomorphic to $\mathbb{R}^{n}$ as vector spaces (and hence as manifolds) and thus as described above they form a locally trivial fibration $T M:=\bigcup_{x \in M} T_{x} M \rightarrow$ $M$ : Each chart $(h, U)$ for $M$ induces a pre-bundle chart for $T M$ via

$$
\varphi_{(h, U)}: T M / U \rightarrow U \times \mathbb{R}^{n}, \quad v \mapsto\left(\pi(v), v\left(h^{1}\right), \ldots, v\left(h^{n}\right)\right)
$$

For any further chart $(k, V)$ for $M$ the transition between the pre-bundle charts is given by

$$
(U \cap V) \times \mathbb{R}^{n} \rightarrow(U \cap V) \times \mathbb{R}^{n}, \quad(x, w) \mapsto\left(x, J_{h(x)}\left(k \circ h^{-1}\right) w\right)
$$

and hence is differentiable. Thus $T M \rightarrow M$ constitutes a fiber bundle. Moreover, since the bundle charts $\varphi_{(h, U)}$ are linear in each fiber, $T M \rightarrow M$ defines an $n$-dimensional vector bundle, the tangent bundle of $M$.

DEFINITION 1.0.29 (Induced fibration). Let $f: M \rightarrow N$ be a differentiable map between manifolds and $\pi: E \rightarrow N$ a locally trivial fibration with typical fiber $F$. Then also $f^{*} E:=\bigcup_{x \in M} E_{f(x)} \rightarrow M$ with the canonical projection is a locally trivial fibration with typical fiber $F$. To this end bundle charts $(\varphi, U)$ for $E$ provide fiberwise "induced" pre-bundle charts $\left(f^{*} \varphi, f^{-1}(U)\right)$ for $f^{*} E \equiv\{(x, e) \in M \times E: f(x)=\pi(e)\}$ via

$$
f^{*} \varphi: f^{*} E / f^{-1}(U) \rightarrow f^{-1}(U) \times F, \quad f^{*} \varphi\left|\left(f^{*} E\right)_{x} \equiv \varphi\right| E_{f(x)} \text { for } x \in f^{-1}(U)
$$

These induced charts change in a differentiable way, and by Lemma $1.0 .27, f^{*} E$ is a locally trivial fibration with base $M$, called the fibration induced from $E$ by $f$ or the pullback fibration under $f$.

Example 1.0.30 (Induced vector bundle). Let $f: M \rightarrow N$ be a differentiable map between manifolds, and $E \rightarrow N$ be a vector bundle. Then $f^{*} E \rightarrow M$ is a vector bundle as well, the so-called pullback of $E$ under $f$. For a bundle homomorphism $\phi: E \rightarrow E^{\prime}$ over $N$ there is again fiberwise a bundle homomorphism $f^{*} \phi: f^{*} E \rightarrow f^{*} E^{\prime}$ over $M$, defined via $f^{*} \phi\left|\left(f^{*} E\right)_{x} \equiv \phi\right| E_{f(x)}$. This constitutes a covariant functor $f^{*}: \mathscr{V}_{N} \rightarrow \mathscr{V}_{M}$.

Example 1.0.31. Let $f: M \rightarrow N$ be a differentiable map between manifolds. There is a canonical bundle homomorphism $d f: T M \rightarrow f^{*} T N$ over $M$ fiberwise explained by the differential $d f_{x}: T_{x} M \rightarrow T_{f(x)} N$.

DEFINITION 1.0.32 (Section, vector field along a map). Let $f: M \rightarrow N$ be a differentiable map between manifolds, and $E$ a vector bundle over $N$. The elements of the $C^{\infty}(M)$-module

$$
\Gamma\left(f^{*} E\right) \equiv\{A: M \rightarrow E \mid A \text { differentiable with } \pi \circ A=f\}
$$

are called the sections along $f$, or in the special case of the $C^{\infty}(M)$-module $\Gamma\left(f^{*} T N\right)$, the vector fields along $f$. In particular, if $I \subset \mathbb{R}$ is an interval and $\alpha: I \rightarrow N$ a differentiable curve, then

$$
\Gamma\left(\alpha^{*} E\right) \equiv\left\{\sigma: I \rightarrow E \mid \sigma \text { differentiable with } \sigma(t) \in E_{\alpha(t)} \text { for each } t \in I\right\}
$$

and the vector field along $\alpha$ given by

$$
\dot{\alpha} \in \Gamma\left(\alpha^{*} T N\right), \quad \dot{\alpha}_{t}:=\dot{\alpha}(t)
$$

is called the tangential vector field along $\alpha$.


Figure 1.0.4. Vector field $\sigma$ along the curve $\alpha$

THEOREM 1.0.33 (Linear algebra for vector bundles). Let $\mathscr{V}$ be the category of finitedimensional real vector spaces and $\mathscr{V}_{M}$ the category of vector bundles over a manifold $M$. Further let

$$
\mathcal{F}: \mathscr{V} \times r \times \mathscr{V} \times s \rightarrow \mathscr{V}
$$

be an r-times covariant and s-times contravariant functor which is differentiable in the sense that the maps induced by $\mathcal{F}$

$$
\begin{aligned}
\operatorname{Hom}\left(V_{1}, V_{1}^{\prime}\right) & \times \cdots \times \operatorname{Hom}\left(V_{r}, V_{r}^{\prime}\right) \times \operatorname{Hom}\left(W_{1}^{\prime}, W_{1}\right) \times \cdots \times \operatorname{Hom}\left(W_{s}^{\prime}, W_{s}\right) \\
& \rightarrow \operatorname{Hom}\left(\mathcal{F}\left(V_{1}, \ldots, V_{r}, W_{1}, \ldots, W_{s}\right), \mathcal{F}\left(V_{1}^{\prime}, \ldots, V_{r}^{\prime}, W_{1}^{\prime}, \ldots, W_{s}^{\prime}\right)\right)
\end{aligned}
$$

are differentiable. Then, by fiberwise application, $\mathcal{F}$ induces canonically an $r$-times covariant and $s$-times contravariant functor

$$
\mathcal{F}_{M}: \mathscr{V}_{M}^{\times r} \times \mathscr{V}_{M}^{\times s} \rightarrow \mathscr{V}_{M}
$$

In a sloppy form Theorem 1.0.33 means the following: One decomposes vector bundle, bundle charts, resp. bundle homomorphisms, into its fiber parts, applies fiberwise the construction rule $\mathcal{F}$ in $\mathscr{V}$, and glues the result again together to new bundles, bundle charts and morphisms. The differentiability condition on $\mathcal{F}$ guarantees automatically the conditions of Lemma 1.0.27, necessary to give the still missing differentiable structures. Canonical examples for suitable functors are:

| $\mathcal{F}\left(V_{1}, \ldots, V_{r}, W_{1}, \ldots, W_{s}\right)$ | $r$ | $s$ |
| :--- | :--- | :--- |
| $V_{1} \oplus \cdots \oplus V_{r}$ | $r$ | 0 |
| $V_{1} \otimes \cdots \otimes V_{r}$ | $r$ | 0 |
| $W^{*}$ | 0 | 1 |
| $\operatorname{Hom}(W, V)$ | 1 | 1 |
| $\operatorname{Mult}\left(W_{1}, \ldots, W_{s} ; V\right)$ | 1 | $s$ |
| $\operatorname{Bil}\left(W_{1}, W_{2} ; \mathbb{R}\right)$ | 0 | 2 |
| $\operatorname{Alt}^{k}(W, V)$ | 1 | 1 |

In the case $W_{1}=\cdots=W_{s}=W$ one writes $\operatorname{Mult}\left(W^{s} ; V\right)$ for $\operatorname{Mult}\left(W_{1}, \ldots, W_{s} ; V\right)$. Usually one also writes furthermore $\mathcal{F}$ instead of $\mathcal{F}_{M}$, e.g. $E_{1} \oplus E_{2}$ instead of $E_{1} \oplus_{M} E_{2}$ for vector bundle $E_{1}, E_{2}$ over $M$.

Canonical isomorphisms in $\mathscr{V}$ carry over to canonical isomorphisms in $\mathscr{V}_{M}$. Typical examples are among others:

$$
\begin{array}{ll}
\mathcal{F}\left(V_{1}, \ldots, V_{r}, W_{1}, \ldots, W_{s}\right) \cong & \mathcal{F}^{\prime}\left(V_{1}, \ldots, V_{r}, W_{1}, \ldots, W_{s}\right) \\
\hline \operatorname{Hom}(W, V) & W^{*} \otimes V \\
W_{1}^{*} \otimes W_{2}^{*} & \left(W_{1} \otimes W_{2}\right)^{*} \\
\operatorname{Bil}(W, W ; \mathbb{R}) & W^{*} \otimes W^{*} \\
\operatorname{Mult}\left(W_{1}, \ldots, W_{s} ; V\right) & W_{1}^{*} \otimes \cdots \otimes W_{s}^{*} \otimes V
\end{array}
$$

DEFINITION 1.0.34 (Vector field). Let $M$ be a manifold and $\pi: T M \rightarrow M$ the tangent bundle of $M$. The elements of the $C^{\infty}(M)$-module $\Gamma(T M)$ are called vector fields on $M$.

Vector fields can be read as derivations by means of the canonical $C^{\infty}(M)$-isomorphism

$$
\Gamma(T M) \rightarrow \operatorname{Der}_{\mathbb{R}} C^{\infty}(M), \quad A \mapsto(f \mapsto A f) ;
$$

here for $f \in C^{\infty}(M)$ the function $A f: M \rightarrow \mathbb{R}$ is explained by $A f(x):=A_{x}(f)$. This gives the product rule $A(f g)=f A g+g A f$ for $f, g \in C^{\infty}(M)$. For an arbitrary
map $A: M \rightarrow T M$ with $\pi \circ A=\operatorname{id}_{M}$ one verifies that $A \in \Gamma(T M)$ if and only if $A f \in C^{\infty}(M)$ for each function $f \in C^{\infty}(M)$.

Example 1.0.35 (Vector fields in coordinates). Let $A \in \Gamma(T M)$ be a vector field on $M$ and $(h, U)$ be a chart for $M$. There exist uniquely determined functions $a_{i} \in C^{\infty}(U)$ such that $A \mid U=\sum a_{i} \partial_{i}$; here $\partial_{i}=\frac{\partial}{\partial h^{i}}$ denotes for $i=1, \ldots, n$ the derivation defined by

$$
\left(\frac{\partial}{\partial h^{i}}\right)_{x}(f)=\frac{\partial}{\partial x^{i}}\left(f \circ h^{-1}\right)(h(x)), \quad x \in U
$$

(see Example 1.0.16). In the special case $M=U \subset \mathbb{R}^{n}$, according to the canonical trivialization $T U \cong U \times \mathbb{R}^{n}$ (via the global chart id ${ }_{U}$ ), each vector field $A \in \Gamma(T U)$ is of the form $A=\left(\mathrm{id}_{U}, a\right)$ where $a \in C^{\infty}\left(U ; \mathbb{R}^{n}\right)$, and the map

$$
C^{\infty}\left(U ; \mathbb{R}^{n}\right) \xrightarrow{\sim} \Gamma(T U), \quad a \mapsto\left(\mathrm{id}_{U}, a\right)
$$

is a $C^{\infty}(U)$-isomorphism. In this situation the canonical vector fields to the constant maps $\left(x \mapsto e_{i}\right) \in C^{\infty}\left(U ; \mathbb{R}^{n}\right)$ are denoted by $D_{i}$ (or $D$ if $n=1$ ); as derivations the $D_{i}$ operate via $D_{i} f=\frac{\partial}{\partial x^{i}} f$ for $f \in C^{\infty}(U)$ (and for $n=1$ again by $D f=\frac{d}{d x} f$ ).

DEfinition 1.0.36 (Cotangent bundle, differential form). Let $M$ be a manifold. The vector bundle $T^{*} M \equiv(T M)^{*}$ over $M$ is called the cotangent bundle of $M$; the elements of $A^{1}(M):=\Gamma\left(T^{*} M\right)$ are denoted differential forms on $M$.

For $f \in C^{\infty}(M)$ let $d f \in A^{1}(M)$ be the differential form defined by

$$
(d f)_{x} \equiv T_{x} f \in T_{x}^{*} M
$$

Given $\alpha \in A^{1}(M)$ and $(h, U)$ a chart for $M$, there are unique functions $\alpha_{i} \in C^{\infty}(U)$ such that $\alpha \mid U=\sum \alpha_{i} d h^{i}$. Note that $d h^{i}\left(\frac{\partial}{\partial h^{j}}\right)=\delta_{i j}$ for $i, j=1, \ldots, n$.

REMARK 1.0.37 (Integral curve). Vector fields can be integrated to integral curves. Let $A \in \Gamma(T M)$ be a vector field on $M$ and $x \in M$. A differentiable curve $\varphi: I \rightarrow M$ (where $I \subset \mathbb{R}$ is an open interval about 0 ) is said to be an integral curve to the vector field $A$ with starting point $x$, if

$$
\varphi(0)=x \quad \text { and } \quad \dot{\varphi}(t)=A(\varphi(t)) \quad \text { for } t \in I
$$

DEFINITION 1.0.38 (Local flow). A local flow on a manifold $M$ is a differentiable $\operatorname{map} \phi: D \rightarrow M$, where $D \subset \mathbb{R} \times M$ is an open neighbourhood of $\{0\} \times M$ and each $I_{x}:=\{t \in \mathbb{R}:(t, x) \in D\}$ an interval, such that the following two conditions are satisfied:
(i) $\phi(0, x)=x$
(ii) $\phi(s, \phi(t, x))=\phi(s+t, x)$ whenever the left-hand side is explained.

For any $x \in M$ the curve $\varphi_{x}: I_{x} \rightarrow M, t \mapsto \phi(t, x)$ is called flow line with starting point $x$. (As a consequence of condition (ii) along with the fact that $D$ is open, flow lines are automatically maximal).

REMARK 1.0.39. Via reduction to the Existence and Uniqueness Theorem for solutions of first order ordinary differential equations, we conclude that to any vector field $A$ on a manifold $M$ there exists a local flow $\phi$ on $M$ whose flow lines coincide with the maximal integral curves to $A$, i.e. such that for $\varphi_{x}(t)=\phi(t, x)$ the following flow equation holds:

$$
\begin{equation*}
\dot{\varphi}_{x}(t)=A\left(\varphi_{x}(t)\right), \quad \varphi_{x}(0)=x \tag{1.0.1}
\end{equation*}
$$

### 1.1. Stochastic Flows

In the same way as a vector field on a differentiable manifold induces a flow, second order differential operators induce stochastic flows with similar properties. In this sense, Brownian motion on a Riemannian manifold $M$ appears as the stochastic flow associated to the canonical Laplacian on $M$, the so-called Laplace-Beltrami operator. The new feature of stochastic flows is that the flow curves depend on a random parameter and behave irregularly as functions of time [29]. This irregularity reveals an irreversibility of time which is inherent to stochastic phenomena.

Let $M$ be a differentiable manifold of dimension $n$ and denote by

$$
T M \xrightarrow{\pi} M
$$

its tangent bundle. In particular, from a set-theoretical point of view, we have

$$
T M=\dot{\cup}_{x \in M} T_{x} M, \quad \pi \mid T_{x} M=x
$$

The space of smooth sections of $T M$ is denoted by

$$
\begin{aligned}
\Gamma(T M) & =\left\{A: M \rightarrow T M \text { smooth } \mid \pi \circ A=\mathrm{id}_{M}\right\} \\
& =\left\{A: M \rightarrow T M \text { smooth } \mid A(x) \in T_{x} M \text { for all } x \in M\right\}
\end{aligned}
$$

and constitutes the vector fields on $M$. As usual, we identify vector fields on $M$ and $\mathbb{R}$-derivations on $C^{\infty}(M)$ as follows:
$\Gamma(T M) \widehat{=}\left\{A: C^{\infty}(M) \rightarrow C^{\infty}(M) \mathbb{R}\right.$-linear $\left.\mid A(f g)=f A(g)+g A(f) \forall f, g \in C^{\infty}(M)\right\}$
where a vector field $A \in \Gamma(T M)$ is considered as $\mathbb{R}$-derivation via

$$
\begin{equation*}
A(f)(x):=d f_{x} A(x) \in \mathbb{R}, \quad x \in M \tag{1.1.1}
\end{equation*}
$$

using the differential $d f_{x}: T_{x} M \rightarrow \mathbb{R}$ of $f$ at $x$.
There is a dynamical point of view to vector fields on manifolds: it associates to each vector field a dynamical system given by the flow of the vector field.
1.1.1. Flow of a vector field. Given a vector field $A \in \Gamma(T M)$. For each $x \in M$ we consider the smooth curve $t \mapsto x(t)$ in $M$ with the properties

$$
x(0)=x \text { and } \dot{x}(t)=A(x(t))
$$

We write $\phi_{t}(x):=x(t)$. In this way, for $A \in \Gamma(T M)$, the flow to $A$ is given by

$$
\left\{\begin{array}{l}
\frac{d}{d t} \phi_{t}=A\left(\phi_{t}\right)  \tag{1.1.2}\\
\phi_{0}=\mathrm{id}_{M}
\end{array}\right.
$$

System (1.1.2) is understood in the sense that for any $f \in C_{c}^{\infty}(M)$ (space of compactly supported smooth functions on $M$ ) the following conditions hold:

$$
\left\{\begin{array}{l}
\frac{d}{d t}\left(f \circ \phi_{t}\right)=A(f) \circ \phi_{t}  \tag{1.1.3}\\
f \circ \phi_{0}=f
\end{array}\right.
$$

Indeed, by the chain rule along with definition (1.1.1), we have for each $f \in C_{c}^{\infty}(M)$,

$$
\frac{d}{d t}\left(f \circ \phi_{t}\right)=(d f)_{\phi_{t}} \frac{d}{d t} \phi_{t}=(d f)_{\phi_{t}} A\left(\phi_{t}\right)=A(f)\left(\phi_{t}\right)
$$

In integrated form, for each $f \in C_{c}^{\infty}(M)$, the conditions (1.1.3) write as:

$$
\begin{equation*}
f \circ \phi_{t}(x)-f(x)-\int_{0}^{t} A(f)\left(\phi_{s}(x)\right) d s=0, \quad t \geq 0, x \in M \tag{1.1.4}
\end{equation*}
$$

As usual, the curve

$$
\phi .(x): t \mapsto \phi_{t}(x)
$$

is called flow curve (or integral curve) to $A$ starting at $x$.
Remark 1.1.1. Defining $P_{t} f:=f \circ \phi_{t}$, we observe that $\frac{d}{d t} P_{t} f=P_{t}(A(f))$, in particular

$$
\begin{equation*}
\left.\frac{d}{d t}\right|_{t=0} P_{t} f=A(f) \tag{1.1.5}
\end{equation*}
$$

In other words, from the knowledge of the flow $\phi_{t}$, the underlying vector field $A$ can be recovered by taking the derivative at zero as in Eq. (1.1.5).
1.1.2. Flow to a second order differential operator. Now let $L$ be a second order partial differential operator (PDO) on $M$, e.g. of the form

$$
\begin{equation*}
L=A_{0}+\sum_{i=1}^{r} A_{i}^{2} \tag{1.1.6}
\end{equation*}
$$

where $A_{0}, A_{1}, \ldots, A_{r} \in \Gamma(T M)$ for some $r \in \mathbb{N}$. Note that $A_{i}^{2}=A_{i} \circ A_{i}$ is understood as composition of derivations, i.e.

$$
A_{i}^{2}(f)=A_{i}\left(A_{i}(f)\right), \quad f \in C^{\infty}(M)
$$

Example 1.1.2. Let $M=\mathbb{R}^{n}$ and consider

$$
A_{0}=0 \text { and } A_{i}=\frac{\partial}{\partial x_{i}} \text { for } i=1, \ldots, n
$$

Then $L=\Delta$ is the classical Laplace operator on $\mathbb{R}^{n}$.
Alternatively, we may consider partial differentiable operators $L$ on $M$ which locally in a chart $(h, U)$ can be written as

$$
\begin{equation*}
L \mid U=\sum_{i=1}^{n} b_{i} \partial_{i}+\sum_{i, j=1}^{n} a_{i j} \partial_{i} \partial_{j} \tag{1.1.7}
\end{equation*}
$$

where $b \in C^{\infty}\left(U, \mathbb{R}^{n}\right)$ and $a \in C^{\infty}\left(U, \mathbb{R}^{n} \otimes \mathbb{R}^{n}\right)$ such that $a_{i j}=a_{j i}$ for all $i, j$ ( $a$ symmetric). Here we use the notation $\partial_{i}=\frac{\partial}{\partial h_{i}}$.

Motivated by the example of a flow to a vector field (vector fields can be seen as first order differential operators) we want to investigate the question whether an analogous concept of flow exists for second order PDOs.

Question. Is there a notion of a flow to $L$ if $L$ is a second order PDO given by (1.1.6) or (1.1.7)?

DEFINITION 1.1.3. Let $\left(\Omega, \mathscr{F}, \mathbb{P} ;\left(\mathscr{F}_{t}\right)_{t \geq 0}\right)$ be a filtered probability space, i.e. a probability space equipped with increasing sequence of sub- $\sigma$-algebras $\mathscr{F}_{t}$ of $\mathscr{F}$. An adapted continuous process

$$
X .(x) \widehat{=}\left(X_{t}(x)\right)_{t \geq 0}
$$

on $\left(\Omega, \mathscr{F}, \mathbb{P} ;\left(\mathscr{F}_{t}\right)_{t \geq 0}\right)$ taking values in $M$, is called flow process to $L$ (or L-diffusion) with starting point $x$ if $\bar{X}_{0}(x)=x$ and if, for all test functions $f \in C_{c}^{\infty}(M)$, the process

$$
\begin{equation*}
N_{t}^{f}(x):=f\left(X_{t}(x)\right)-f(x)-\int_{0}^{t}(L f)\left(X_{s}(x)\right) d s, \quad t \geq 0 \tag{1.1.8}
\end{equation*}
$$

is a martingale, i.e.

$$
\mathbb{E}^{\mathscr{F}_{s}} \underbrace{\left[f\left(X_{t}(x)\right)-f\left(X_{s}(x)\right)-\int_{s}^{t}(L f)\left(X_{r}(x)\right) d r\right]}_{=N_{t}^{f}(x)-N_{s}^{f}(x)}=0, \quad \text { for all } s \leq t
$$

Note that, by definition, flow processes to a second order PDO depend on an additional random parameter $\omega \in \Omega$. For each $t \geq 0, X_{t}(x) \equiv\left(X_{t}(x, \omega)\right)_{\omega \in \Omega}$ is an $\mathscr{F}_{t}$-measurable random variable. The defining equation (1.1.4) for flow curves translates to the martingale property of (1.1.8), i.e. the flow curve condition (1.1.4) only holds under conditional expectations. The theory of martingales gives a rigorous meaning to the idea of a process without systematic drift [46].

Flow processes will be constructed as solutions to certain stochastic differential equations on $M$, which degenerate to the flow equation (1.0.1) in the particular case of vector fields. The second order part of the differential operator causes the "flow lines" to depend now on random in an intriguing way. The paths of flow processes are still continuous, but are in general nowhere differentiable anymore.

REMARK 1.1.4. Since $N_{0}^{f}(x)=0$, we get from the martingale property of $N^{f}(x)$ that

$$
\mathbb{E}\left[N_{t}^{f}(x)\right]=\mathbb{E}\left[N_{0}^{f}(x)\right]=0
$$

Hence, defining $P_{t} f(x):=\mathbb{E}\left[f\left(X_{t}(x)\right)\right]$, we observe that

$$
P_{t} f(x)=f(x)+\int_{0}^{t} \mathbb{E}\left[(L f)\left(X_{s}(x)\right)\right] d s
$$

and thus

$$
\frac{d}{d t} P_{t} f(x)=\mathbb{E}\left[(L f)\left(X_{t}(x)\right)\right]=P_{t}(L f)(x)
$$

in particular

$$
\left.\left.\frac{d}{d t}\right|_{t=0} \mathbb{E}\left[f\left(X_{t}(x)\right)\right] \equiv \frac{d}{d t}\right|_{t=0} P_{t} f(x)=L f(x)
$$

The last formula shows that as for deterministic flows we can recover the operator $L$ from its stochastic flow process. To this end however, we have to average over all possible trajectories starting from $x$.

For background on stochastic flows we refer to the monograph of Kunita [29].
EXAMPLE 1.1.5 (Brownian motion). Let $M=\mathbb{R}^{n}$ and $L=\frac{1}{2} \Delta$ where $\Delta$ is the Laplacian on $\mathbb{R}^{n}$. Let $X \equiv\left(X_{t}\right)$ be a Brownian motion on $\mathbb{R}^{n}$ starting at the origin. By Itô's formula [38], for $f \in C^{\infty}\left(\mathbb{R}^{n}\right)$, we have

$$
\begin{aligned}
d\left(f \circ X_{t}\right) & =\sum_{i=1}^{n} \partial_{i} f\left(X_{t}\right) d X_{t}^{i}+\frac{1}{2} \sum_{i, j=1}^{n} \partial_{i} \partial_{j} f\left(X_{t}\right) d X_{t}^{i} d X_{t}^{j} \\
& =\left\langle(\nabla f)\left(X_{t}\right), d X_{t}\right\rangle+\frac{1}{2}(\Delta f)\left(X_{t}\right) d t
\end{aligned}
$$

Thus, for each $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
f\left(X_{t}\right)-f\left(X_{0}\right)-\int_{0}^{t} \frac{1}{2}(\Delta f)\left(X_{s}\right) d s, \quad t \geq 0
$$

is a martingale. This means that the process

$$
X_{t}(x):=x+X_{t}
$$

is an $L$-diffusion to $\frac{1}{2} \Delta$ in the sense of Definition 1.1.3.
REMARKS 1.1.6. As for deterministic flows, we have to deal with the problem that stochastic flows may explode in finite times.

1. We allow $X .(x)$ to be defined only up to some stopping time $\zeta(x)$, i.e.

$$
X .(x) \mid[0, \zeta(x)[
$$

where

$$
\begin{equation*}
\{\zeta(x)<\infty\} \subset\left\{\lim _{t \uparrow \zeta(x)} X_{t}(\omega)=\infty \text { in } \hat{M}:=M \dot{\cup}\{\infty\}\right\} \quad \mathbb{P} \text {-a.s. } \tag{1.1.9}
\end{equation*}
$$

Here $\hat{M}$ denotes the one-point compactification of $M$. A stopping time $\zeta(x)$ with property (1.1.9) is called (maximal) lifetime for the process $X .(x)$ starting at $x$. In equivalent terms, let $U_{n} \subset M$ be open, relatively compact subsets exhausting $M$ in the sense that

$$
U_{n} \subset \bar{U}_{n} \subset U_{n+1} \subset \ldots, \quad \bar{U}_{n} \text { compact, and } \cup_{n} U_{n}=M
$$

Then we have $\zeta(x)=\sup _{n} \tau_{n}(x)$ for the maximal lifetime of $X .(x)$ where $\tau_{n}(x)$ is the family of stopping times (first exit times of $U_{n}$ ) defined by

$$
\tau_{n}(x):=\inf \left\{t \geq 0: X_{t}(x) \notin U_{n}\right\}
$$

2. For $f \in C^{\infty}(M)$ (not necessarily compactly supported), the process $N^{f}(x)$ will in general only be a local martingale [38], i.e. there exist stopping times $\tau_{n} \uparrow \zeta(x)$ such that

$$
\forall n \in \mathbb{N}, \quad\left(N_{t \wedge \tau_{n}}^{f}(x)\right)_{t \geq 0} \text { is a (true) martingale. }
$$

3. The following two statements are equivalent (the proof will be given later):
(a) The process

$$
f(X .(x))=\left(f\left(X_{t}(x)\right)\right)_{t \geq 0}
$$

is of locally bounded variation for all $f \in C_{c}^{\infty}(M)$.
(b) The operator $L$ is of first order, i.e. $L$ is a vector field (in which case the flow is deterministic).
In other words, flow processes have "nice paths" (for instance, paths of bounded variation) if and only if the corresponding operator is first order (i.e. a vector field).
1.1.3. What are $L$-diffusions good for? Before discussing the problem of how to construct $L$-diffusions, we want to study some implications to indicate the usefulness and power of this concept. In the following two examples we only assume existence of an $L$-diffusion to a given operator $L$.
A. (Dirichlet problem) Let $\varnothing \neq D \subsetneq M$ be an open, connected, relatively compact domain, $\varphi \in C(\partial D)$ and let $L$ be a second order PDO on $M$. The Dirichlet problem (DP) is the problem to find a function $u \in C(\bar{D}) \cap C^{2}(D)$ such that

$$
\left\{\begin{array}{l}
L u=0 \text { on } D  \tag{DP}\\
\left.u\right|_{\partial D}=\varphi
\end{array}\right.
$$

Suppose that there is an $L$-diffusion $\left(X_{t}(x)\right)_{t \geq 0}$. We choose a sequence of open domains $D_{n} \uparrow D$ such that $\bar{D}_{n} \subset D$, and for each $n$, we consider the first exit time of $D_{n}$,

$$
\tau_{n}(x)=\inf \left\{t \geq 0, X_{t}(x) \notin D_{n}\right\}
$$

Then $\tau_{n}(x) \uparrow \tau(x)$ where

$$
\tau(x)=\sup _{n} \tau_{n}(x)=\inf \left\{t \geq 0, X_{t}(x) \notin D\right\}
$$

Now assume that $u$ is a solution to (DP). We may choose test functions $u_{n} \in C_{c}^{\infty}(M)$ such that $u_{n}\left|D_{n}=u\right| D_{n}$ and $\operatorname{supp} u_{n} \subset D$. Then, by the property of an $L$-diffusion,

$$
N_{t}(x):=u_{n}\left(X_{t}(x)\right)-u_{n}(x)-\int_{0}^{t}\left(L u_{n}\right)\left(X_{r}(x)\right) d r
$$

is a martingale. We suppose that $x \in D_{n}$. Then

$$
\begin{align*}
N_{t \wedge \tau_{n}(x)}(x) & =u_{n}\left(X_{t \wedge \tau_{n}(x)}(x)\right)-u_{n}(x)-\int_{0}^{t \wedge \tau_{n}(x)} \underbrace{\left(L u_{n}\right)\left(X_{r}(x)\right)}_{=0} d r  \tag{1.1.10}\\
& =u\left(X_{t \wedge \tau_{n}(x)}(x)\right)-u(x)
\end{align*}
$$

is also a martingale (here we used that the integral in (1.1.10) is zero since $L u_{n}=L u=0$ on $D_{n}$ ). Thus we get

$$
\mathbb{E}\left[N_{t \wedge \tau_{n}(x)}(x)\right]=\mathbb{E}\left[N_{0}(x)\right]=0
$$

which shows that for each $n \in \mathbb{N}$,

$$
\begin{equation*}
u(x)=\mathbb{E}\left[u\left(X_{t \wedge \tau_{n}(x)}(x)\right)\right] \tag{1.1.11}
\end{equation*}
$$

From Eq. (1.1.11) we may conclude by dominated convergence and since $\tau_{n}(x) \uparrow \tau$ that

$$
u(x)=\lim _{n \rightarrow \infty} \mathbb{E}\left[u\left(X_{t \wedge \tau_{n}(x)}(x)\right)\right]=\mathbb{E}\left[\lim _{n \rightarrow \infty} u\left(X_{t \wedge \tau_{n}(x)}(x)\right)\right]=\mathbb{E}\left[u\left(X_{t \wedge \tau(x)}(x)\right)\right]
$$

We now make the hypothesis that $\tau(x)<\infty$ a.s. (the process exits the domain $D$ in finite time). Then

$$
\begin{aligned}
u(x) & =\lim _{t \rightarrow \infty} \mathbb{E}\left[u\left(X_{t \wedge \tau(x)}(x)\right)\right]=\mathbb{E}\left[\lim _{t \rightarrow \infty} u\left(X_{t \wedge \tau(x)}(x)\right)\right] \\
& =\mathbb{E}\left[u\left(X_{\tau(x)}(x)\right)\right]=\mathbb{E}\left[\varphi\left(X_{\tau(x)}(x)\right)\right]
\end{aligned}
$$

where for the last equality we used the boundary condition $u \mid \partial D=\varphi$. Note that by passing to the image measure $\mu_{x}:=\mathbb{P} \circ X_{\tau(x)}(x)^{-1}$ on the boundary we get

$$
\mathbb{E}\left[\varphi\left(X_{\tau(x)}(x)\right)\right]=\int_{\partial D} \varphi(z) \mu_{x}(d z)
$$

Notation 1.1.7. The measure $\mu_{x}$, defined on Borel sets $A \subset \partial D$,

$$
\mu_{x}(A)=\mathbb{P}\left\{X_{\tau(x)}(x) \in A\right\}
$$

is called exit measure from the domain $D$ of the diffusion $X_{t}(x)$. It represents the probability that the process $X_{t}$, when started at $x$ in $D$, exits the domain $D$ through the boundary set $A$.

Conclusions. From the discussion of the Dirichlet problem above we can make the following two observations.
(a) (Uniqueness) Under the hypothesis

$$
\tau(x)<\infty \text { a.s. for all } x \in D
$$

we have uniqueness of the solutions to the Dirichlet problem (DP). It will be shown later that this hypothesis concerns non-degeneracy of the operator $L$.
(b) (Existence) Under the hypothesis

$$
\tau(x) \rightarrow 0 \text { if } D \ni x \rightarrow a \in \partial D
$$

we have

$$
\mathbb{E}\left[\varphi\left(X_{\tau(x)}(x)\right)\right] \rightarrow \varphi(a), \quad \text { if } D \ni x \rightarrow a \in \partial D
$$

Thus one may define $u(x):=\mathbb{E}\left[\varphi\left(X_{\tau(x)}(x)\right)\right]$. It can be shown then that $u$ is $L$ harmonic on $D$ if it is twice differentiable; thus under the hypothesis in (b), $u$ will then satisfy the boundary condition and hence solve (DP). The hypothesis in (b) is obviously a regularity condition on the boundary $\partial D$.

Note that in the arguments above we nowhere used the explicit form of the operator $L$ nor of the domain $D$. We only used the general properties of a stochastic flow process associated to the given operator $L$. For a more complete discussion of the Dirichlet problem see $[43,2]$.

EXAMPLES 1.1.8.
(1) Let $M=\mathbb{R}^{2} \backslash\{0\}$ and $D=\left\{x \in \mathbb{R}^{2}: r_{1}<|x|<r_{2}\right\}$ with $0<r_{1}<r_{2}$. Consider the operator

$$
L=\frac{1}{2} \frac{\partial^{2}}{\partial \vartheta^{2}}
$$

where $\vartheta$ denotes the angle when passing to polar coordinates on $M$. If $u$ is a solution of (DP), then $u+v(r)$ is a solution of (DP) as well, for any radial function $v(r)$ satisfying $v\left(r_{1}\right)=v\left(r_{2}\right)=0$. Hence, uniqueness of solutions fails.


Note: For $x \in D$ with $|x|=r$, let $S_{r}=\left\{x \in \mathbb{R}^{2}:|x|=r\right\}$. Then, the flow process $X .(x)$ to $L$ is easily seen to be a (one-dimensional) Brownian motion on $S_{r}$. In particular,

$$
\tau(x)=+\infty \text { a.s. }
$$

(2) Let $M=\mathbb{R}^{2}$ and consider the operator

$$
L=\frac{1}{2} \frac{\partial^{2}}{\partial x_{1}^{2}}
$$

on a domain $D$ in $\mathbb{R}^{2}$ of the following shape:


Then, for $x=\left(x_{1}, x_{2}\right) \in D$, the flow process $X$. $(x)$ starting at $x$ is a (one-dimensional) Brownian motion on $\mathbb{R} \times\left\{x_{2}\right\}$. In other words, flow processes move on horizontal lines. In particular, when started at $x \in D$, the process can only exit at two points (e.g. $x_{\ell}$ and $x_{r}$ in the picture). Letting $x$ vertically approach $a$, by symmetry of the one-dimensional Brownian motion, we see that there exists a solution of (DP) if and only if

$$
\varphi(a)=\frac{\varphi(b)+\varphi(c)}{2}
$$

B. (Heat equation) Let $L$ be a second order PDO on $M$ and fix $f \in C(M)$. The heat equation on $M$ with initial condition $f$ concerns the problem of finding a real-valued function $u=u(t, x)$ defined on $\mathbb{R}_{+} \times M$ such that

$$
\left\{\begin{array}{l}
\left.\frac{\partial u}{\partial t}=L u \quad \text { on }\right] 0, \infty[\times M  \tag{HE}\\
\left.u\right|_{t=0}=f
\end{array}\right.
$$

Suppose now that there is an $L$-diffusion $X .(x)$. It is straightforward to see that the "timespace process" $\left(t, X_{t}(x)\right)$ will then be a $\hat{L}$-diffusion for the parabolic operator

$$
\hat{L}=\frac{\partial}{\partial t}+L
$$

with starting point $(0, x)$. By definition, this means that for all $\varphi \in C^{2}\left(\mathbb{R}_{+} \times M\right)$,

$$
d \varphi\left(t, X_{t}(x)\right)-(\hat{L} \varphi)\left(t, X_{t}(x)\right) d t \stackrel{\mathrm{~m}}{=} 0
$$

where $\underline{\underline{m}}$ denotes equality modulo differentials of local martingales.
From now on we assume non-explosion of the $L$-diffusion. In other words, we adopt the hypothesis that $\zeta(x)=+\infty$ a.s. for all $x \in M$, i.e.

$$
\mathbb{P}\left\{X_{t}(x) \in M, \forall t \geq 0\right\}=1, \quad \forall x \in M
$$

Suppose now that $u$ is a bounded solution of (HE). We fix $t \geq 0$ and consider the restriction $u \mid[0, t] \times M$. Then

$$
u\left(t-s, X_{s}(x)\right)-u(t, x)-\int_{0}^{s}\left[\left(\frac{\partial}{\partial r}+L\right) u(t-r, \cdot)\right]\left(X_{r}(x)\right) d r, \quad 0 \leq s<t
$$

is a local martingale. In other words, fixing $t>0$, we have for $0 \leq s<t$,

$$
\begin{align*}
u\left(t-s, X_{s}(x)\right)=u(t, x) & +\int_{0}^{s} \underbrace{\left(\frac{\partial}{\partial r}+L\right) u(t-r, \cdot)}_{=0, \text { since } u \text { solves (HE) }}\left(X_{r}(x)\right) d r  \tag{1.1.12}\\
& +(\text { local martingale })_{s}
\end{align*}
$$

Since the integral in (1.1.12) vanishes, we see that the local martingale term in (1.1.12) is actually a bounded local martingale (since $u\left(t-s, X_{s}(x)\right)-u(t, x)$ is bounded) and hence a true martingale (equal to zero at time 0 ). Using the martingale property we first take expectations and then pass to the limit as $s \uparrow t$ to obtain
(1.1.13) $u(t, x)=\mathbb{E}\left[u\left(t-s, X_{s}(x)\right)\right] \rightarrow \mathbb{E}\left[u\left(0, X_{t}(x)\right)\right]=\mathbb{E}\left[f\left(X_{t}(x)\right)\right], \quad$ as $s \uparrow t$, where for the limit in (1.1.13) we used dominated convergence (recall that $u$ is bounded).

Conclusion. Under the hypothesis $\zeta(x)=+\infty$ for all $x \in M$, we have uniqueness of bounded solutions to the heat equation (HE). Solutions are necessarily of the form

$$
u(t, x)=\mathbb{E}\left[f\left(X_{t}(x)\right)\right]
$$

Interpretation. The solution $u(t, x)$ at time $t$ and at point $x$ can be constructed as follows: run an $L$-diffusion process starting from $x$ up time $t$, apply the initial condition $f$ to the obtained random position $X_{t}(x)$ at time $t$ and average over all possible paths.

REMARK 1.1.9. If we drop the hypothesis of infinite lifetime $\zeta(x)=+\infty$ for all $x \in M$, then uniqueness of bounded solutions to the heat equation can no longer be expected. There exists always a minimal solution $u$ to the heat equation (HE) in the sense that $u(t, x) \rightarrow 0$ as $x \rightarrow \infty$ in the one-point compactification $\hat{M}=M \cup\{\infty\}$ of $M$. Let $\sigma_{n} \uparrow \zeta(x)$ be an increasing sequence of stopping times. Then the argument above shows

$$
\begin{aligned}
u(t, x) & =\mathbb{E}\left[u\left(t-t \wedge \sigma_{n}, X_{t \wedge \sigma_{n}}(x)\right)\right] \\
& =\mathbb{E}\left[\lim _{n \rightarrow \infty} u\left(t-t \wedge \sigma_{n}, X_{t \wedge \sigma_{n}}(x)\right)\right] \\
& =\mathbb{E}\left[\mathbb{1}_{\{t<\zeta(x)\}} u\left(0, X_{t}(x)\right)\right] \\
& =\mathbb{E}\left[\mathbb{1}_{\{t<\zeta(x)\}} f\left(X_{t}(x)\right)\right] .
\end{aligned}
$$

This gives for the minimal solution the representation

$$
u(t, x)=\mathbb{E}\left[\mathbb{1}_{\{t<\zeta(x)\}} f\left(X_{t}(x)\right)\right]
$$

### 1.1.4. $\Gamma$-operators and quadratic variation.

DEFINITION 1.1.10. Let $L: C^{\infty}(M) \rightarrow C^{\infty}(M)$ be a linear mapping (for instance a second order PDO). The $\Gamma$-operator associated to $L$ ("l'operateur carré du champ") is the bilinear map

$$
\begin{gathered}
\Gamma: C^{\infty}(M) \times C^{\infty}(M) \rightarrow C^{\infty}(M) \text { given as } \\
\Gamma(f, g):=\frac{1}{2}(L(f g)-f L(g)-g L(f))
\end{gathered}
$$

Example 1.1.11. Let $L$ be a second order PDO on $M$ without constant term (i.e. $L 1=0)$. Suppose that in a local chart $(h, U)$ for $M$ the operator $L$ writes as

$$
L \mid C_{U}^{\infty}(M)=\sum_{i, j=1}^{n} a_{i j} \partial_{i} \partial_{j}+\sum_{i=1}^{n} b_{i} \partial_{i}
$$

where $C_{U}^{\infty}(M)=\left\{f \in C^{\infty}(M): \operatorname{supp} f \subset U\right\}$ and $\partial_{i}=\frac{\partial}{\partial h_{i}}$. Then

$$
\Gamma(f, g)=\sum_{i, j=1}^{n} a_{i j}\left(\partial_{i} f\right)\left(\partial_{j} g\right), \quad \forall f, g \in C_{U}^{\infty}(M)
$$

For instance, in the special case that $M=\mathbb{R}^{n}$ and $L=\Delta$, we find

$$
\Gamma(f, f)=|\nabla f|^{2}
$$

REMARK 1.1.12. Let $L$ be a second order PDO. Then the following equivalence holds:
$\Gamma(f, g)=0 \forall f, g \in C^{\infty}(M)$ if and only if $L$ is of first order, i.e. $L \in \Gamma(T M)$.
For instance, if $L=A_{0}+\sum_{i=1}^{r} A_{i}^{2}$, then

$$
\Gamma(f, g)=\sum_{i=1}^{r} A_{i}(f) A_{i}(g)
$$

and in particular

$$
\Gamma \equiv 0 \quad \text { if and only if } \quad A_{1}=A_{2}=\ldots=A_{r}=0
$$

REMARK 1.1.13. A continuous real-valued stochastic process $\left(X_{t}\right)_{t \geq 0}$ is called a semimartingale if it can be decomposed as

$$
\begin{equation*}
X_{t}=X_{0}+M_{t}+A_{t} \tag{1.1.14}
\end{equation*}
$$

where $M$ is a local martingale and $A$ an adapted process of locally bounded variation (with $M_{0}=A_{0}=0$ ). The representation of a semimartingale $X$ as in (1.1.14) (Doob-Meyer decomposition) is unique: if $\mathscr{M}_{0}$ denotes the class of local martingales starting from 0 and $\mathscr{A}_{0}$ is the class of adapted process with paths of locally bounded variation starting from 0 , then $\mathscr{M}_{0} \cap \mathscr{A}_{0}=0$.

DEFINITION 1.1.14. Let $X$ be a continuous adapted process taking values in a manifold $M$. Then $X$ is called semimartingale on $M$ if

$$
f(X) \equiv\left(f\left(X_{t}\right)\right)_{t \geq 0}
$$

is a real-valued semimartingale for all $f \in C^{\infty}(M)$.
REMARK 1.1.15 (Semimartingale with lifetime). As already noted, semimartingales are often defined only up to some predictable stopping time $\xi>0$. By a transformation of time, if required, infinite lifetime can always be achieved. For instance, let $X$ be semimartingale defined on $\left[0, \xi\left[\right.\right.$ and let $\left(\tau_{n}\right)_{n \in \mathbb{N}}$ be an increasing sequence of finite stopping times such that $\tau_{0}=0, \tau_{n}<\xi$ and $\tau_{n} \uparrow \xi$, then

$$
\tau_{n+r}:=\left(\tau_{n}+\frac{r}{1-r}\right) \wedge \tau_{n+1}, \quad 0 \leq r<1,
$$

defines a continuous time-change $\left(\tau_{t}\right)_{t \geq 0}$ with $\tau_{0}=0$ and $\tau_{\infty}=\xi$, and the time-changed process $\hat{X}: \hat{X}_{t}:=X_{\tau_{t}}$ is a semimartingale (with respect to the time-changed filtration) of infinite lifetime.

Obviously the semimartingale property is a local property.
REMARK 1.1.16. Let $\xi$ be a predictable stopping time and $X$ be an $M$-valued process defined on $\left[0, \xi\left[\right.\right.$. Let $\left(\tau_{n}\right)_{n \in \mathbb{N}}$ be a sequence of finite stopping times such that $\tau_{0}=0$, $\tau_{n} \leq \tau_{n+1}$ for $n \in \mathbb{N}$ and $\sup _{n} \tau_{n}=\xi$. The following conditions are equivalent:
(i) $X$ is an $M$-valued semimartingale.
(ii) For any $n \in \mathbb{N}$ the stopped process $X^{\tau_{n}}$ is a semimartingale.
(iii) For any $n \in \mathbb{N}$ the restriction $X \mid\left[\tau_{n}, \tau_{n+1}[\right.$ is a semimartingale, i.e., the process $\left(Y_{t}^{n}\right)_{t \in \mathbb{R}_{+}}$with $Y_{t}^{n}:=X_{\left(\tau_{n}+t\right) \wedge \tau_{n+1}}$ is a semimartingale with respect to the filtration $\left(\mathscr{F}_{t}^{n}\right)_{t \in \mathbb{R}_{+}}$shifted by $\tau_{n}$, i.e. $\mathscr{F}_{t}^{n}:=\mathscr{F}_{\tau_{n}+t}$.

REmark 1.1.17. If $X$ has maximal lifetime $\zeta$, i.e.,

$$
\{\zeta<\infty\} \subset\left\{\lim _{t \uparrow \zeta} X_{t}=\infty \text { in } \hat{M}=M \dot{\cup}\{\infty\}\right\} \text { a.s. }
$$

then $f(X)$ is well-defined as a process globally on $\mathbb{R}_{+}$for all $f \in C_{c}^{\infty}(M)$ (with the convention $f(\infty)=0)$. For $f \in C^{\infty}(M)$, in general,

$$
f(X) \equiv\left(f\left(X_{t}\right)\right)_{t<\zeta}
$$

is only a semimartingale with lifetime $\zeta$.
Proposition 1.1.18. Let $L: C^{\infty}(M) \rightarrow C^{\infty}(M)$ be an $\mathbb{R}$-linear map and $X$ be a semimartingale on $M$ such that for all $f \in C^{\infty}(M)$,

$$
N_{t}^{f}:=f\left(X_{t}\right)-f\left(X_{0}\right)-\int_{0}^{t} L f\left(X_{r}\right) d r
$$

is a continuous local martingale (of same lifetime as $X$ ) (i.e. $d(f(X))-L f(X) d t \stackrel{m}{=} 0$ where $\stackrel{m}{=}$ denotes equality modulo differentials of local martingales). Then, for all $f, g \in$ $C^{\infty}(M)$, the quadratic variation $[f(X), g(X)]$ of $f(X)$ and $g(X)$ is given by

$$
d[f(X), g(X)] \equiv d\left[N^{f}, N^{g}\right]=2 \Gamma(f, g)(X) d t
$$

In particular, $\Gamma(f, f)(X) \geq 0$ a.s.
Proof. Let $f \in C^{\infty}\left(M, \mathbb{R}^{r}\right)$ and $\phi \in C^{\infty}\left(\mathbb{R}^{r}\right)$. Writing as above $\stackrel{\mathrm{m}}{=}$ for equality modulo differentials of local martingales, we have

$$
\begin{equation*}
d(\phi \circ f)(X) \stackrel{m}{=} L(\phi \circ f)(X) d t \tag{1.1.15}
\end{equation*}
$$

Developing the left-hand side in Eq. (1.1.15) by Itô's formula, the function $\phi$ being applied to the semimartingale $f(X)$, we get

$$
\begin{aligned}
& d(\phi(f(X))) \\
& \quad=\sum_{i=1}^{r}\left(D_{i} \phi\right)(f(X)) d\left(f^{i}(X)\right)+\frac{1}{2} \sum_{i, j=1}^{r}\left(D_{i} D_{j} \phi\right)(f(X)) d\left[f^{i}(X), f^{j}(X)\right] \\
& \quad \stackrel{\mathrm{m}}{=} \sum_{i=1}^{r}\left(D_{i} \phi\right)(f(X))\left(L f^{i}\right)(X) d t+\frac{1}{2} \sum_{i, j=1}^{r}\left(D_{i} D_{j} \phi\right)(f(X)) d\left[f^{i}(X), f^{j}(X)\right]
\end{aligned}
$$

where $D_{i}=\partial / \partial x_{i}$. By equating the drift parts we find
$\left(L(\phi \circ f)-\sum_{i=1}^{r}\left(\left(D_{i} \phi\right) \circ f\right)\left(L f^{i}\right)\right)(X) d t=\frac{1}{2} \sum_{i, j=1}^{r}\left(D_{i} D_{j} \phi\right)(f(X)) d\left[f^{i}(X), f^{j}(X)\right]$.
Taking now $r=2$ and considering the special case $\phi(x, y)=x y$, we get with $f=$ $\left(f^{1}, f^{2}\right)$,

$$
\left(L\left(f^{1} f^{2}\right)-f^{1} L\left(f^{2}\right)-f^{2} L\left(f^{1}\right)\right)(X) d t=d\left[f^{1}(X), f^{2}(X)\right] .
$$

This completes the proof since $\left(L\left(f^{1} f^{2}\right)-f^{1} L\left(f^{2}\right)-f^{2} L\left(f^{1}\right)\right)(X)=2 \Gamma\left(f^{1}, f^{2}\right)(X)$.

LEmma 1.1.19. For an $\mathbb{R}$-linear map $L: C^{\infty}(M) \rightarrow C^{\infty}(M)$ the following statements are equivalent:
(i) L is a second order PDO (without constant term)
(ii) $L$ satisfies the second order chain rule, i.e. for all $f \in C^{\infty}\left(M, \mathbb{R}^{r}\right)$ and $\phi \in$ $C^{\infty}\left(\mathbb{R}^{r}\right)$,

$$
L(\phi \circ f)=\sum_{i=1}^{r}\left(D_{i} \phi \circ f\right)\left(L f^{i}\right)+\sum_{i, j=1}^{r}\left(D_{i} D_{j} \phi \circ f\right) \Gamma\left(f^{i}, f^{j}\right) .
$$

Proof. (i) $\Rightarrow$ (ii): Write $L$ in local coordinates as

$$
L \mid C_{U}^{\infty}(M)=\sum_{i, j=1}^{n} a_{i j} \partial_{i} \partial_{j}+\sum_{i=1}^{n} b_{i} \partial_{i}
$$

and use that $\Gamma(f, g)=\sum_{i, j=1}^{n} a_{i j} \partial_{i} f \partial_{j} g$.
(ii) $\Rightarrow$ (i): Determine the action of $L$ on functions $\varphi$ written in local coordinates $(h, U)$ via

$$
L(\varphi) \mid U=L\left(\varphi \circ h^{-1} \circ h\right) \equiv L(\phi \circ f)
$$

where $\phi=\varphi \circ h^{-1}$ and $f=h$. Details are left as an exercise to the reader.
COROLLARY 1.1.20. Let $L: C^{\infty}(M) \rightarrow C^{\infty}(M)$ be an $\mathbb{R}$-linear mapping. Suppose that for each $x \in M$ there exists a semimartingale $X$ on $M$ such that $X_{0}=x$ and such that for each $f \in C^{\infty}(M)$,

$$
f\left(X_{t}\right)-f(x)-\int_{0}^{t} L f(X) d r
$$

is a local martingale. Then $L$ is necessarily a PDO of order at most 2.
In addition, $X$ has "nice" trajectories (e.g. in the sense that $[f(X), f(X)]=0$ for all $f \in C^{\infty}(M)$ ) if and only if $L$ is first order.

Proof. As in the proof of Proposition 1.1.18, for all $f \in C^{\infty}\left(M, \mathbb{R}^{r}\right)$ and $\phi \in$ $C^{\infty}\left(\mathbb{R}^{r}\right)$, we have

$$
\left(L(\phi \circ f)-\sum_{i=1}^{r}\left(D_{i} \phi \circ f\right)\left(L f^{i}\right)+\sum_{i, j=1}^{r}\left(D_{i} D_{j} \phi \circ f\right) \Gamma\left(f^{i}, f^{j}\right)\right)(X)=0,
$$

so that $L$ is a second order PDO by Lemma 1.1.19. The second claim uses

$$
d[f(X), g(X)]=2 \Gamma(f, g)(X) d t, \quad f, g \in C^{\infty}(M)
$$

### 1.2. Construction of Stochastic Flows

Flows to vector fields are classically constructed as solutions of ordinary differential equations on manifolds. In the same way, stochastic flows can be constructed as solutions to stochastic differential equations (SDE) on manifolds. We start by recalling same basic facts about stochastic differential equations on $\mathbb{R}^{n}$.

### 1.2.1. Stochastic differential equations on Euclidean space.

EXAMPLE 1.2.1 (SDE on $\mathbb{R}^{n}$ ). Given $\beta: \mathbb{R}_{+} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ and in addition a function

$$
\sigma: \mathbb{R}_{+} \times \mathbb{R}^{n} \rightarrow \operatorname{Hom}\left(\mathbb{R}^{r}, \mathbb{R}^{n}\right) \equiv \operatorname{Matr}(n \times r ; \mathbb{R})
$$

Let $B$ be a Brownian motion on $\mathbb{R}^{r}$. Now one wants to find a continuous semimartingale $Y$ on $\mathbb{R}^{n}$ such that

$$
d Y_{t}=\beta\left(t, Y_{t}\right) d t+\sigma\left(t, Y_{t}\right) d B_{t}
$$

in the sense of Itô, i.e.

$$
\begin{equation*}
Y_{t}=Y_{0}+\int_{0}^{t} \beta\left(s, Y_{s}\right) d s+\int_{0}^{t} \sigma\left(s, Y_{s}\right) d B_{s} \tag{1.2.1}
\end{equation*}
$$

In Eq. (1.2.1) the first term describes the "systematic part" (drift term) in the evolution of $Y$, whereas the second integral represents the "fluctuating part" (diffusion term).

DEFInItion 1.2.2. An $\mathbb{R}^{n}$-valued stochastic process $\left(Y_{t}\right)_{t \geq 0}$ is called Itô process if it has a representation as

$$
Y_{t}=Y_{0}+\int_{0}^{t} K_{s} d s+\int_{0}^{t} H_{s} d B_{s}
$$

where

- $Y_{0}$ is $\mathscr{F}_{0}$-measurable;
- $K_{s}$ and $H_{s}$ are adapted processes taking values in $\mathbb{R}^{n}$, resp. $\operatorname{Hom}\left(\mathbb{R}^{r}, \mathbb{R}^{n}\right)$;
- $\mathbb{E}\left[\int_{0}^{t}\left|K_{s}\right| d s\right]<\infty$ and $\mathbb{E}\left[\int_{0}^{t} H_{s}^{2} d s\right]<\infty$ for each $t \geq 0$.

PROPOSITION 1.2.3. Let $\beta: \mathbb{R}_{+} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ and $\sigma: \mathbb{R}_{+} \times \mathbb{R}^{n} \rightarrow \operatorname{Hom}\left(\mathbb{R}^{r}, \mathbb{R}^{n}\right)$ be continuous functions. For a continuous semimartingale $Y$ on $\mathbb{R}^{n}$, defined up to some predictable stopping time $\tau$ (i.e. there exists a sequence of stopping times $\tau_{n}<\tau$ with $\tau_{n} \uparrow \tau$ ), the following conditions are equivalent:
(a) $Y$ is a solution of the SDE

$$
\begin{equation*}
d Y_{t}=\beta\left(t, Y_{t}\right) d t+\sigma\left(t, Y_{t}\right) d B_{t} \quad \text { on }[0, \tau[ \tag{1.2.2}
\end{equation*}
$$

i.e.,

$$
Y_{t}=Y_{0}+\int_{0}^{t} \beta\left(s, Y_{s}\right) d s+\int_{0}^{t} \sigma\left(s, Y_{s}\right) d B_{s}, \quad \forall 0 \leq t<\tau \text { a.s. }
$$

(b) For all $f \in C^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
d(f(Y))=(L f)(t, Y) d t+\sum_{k=1}^{n} \sum_{i=1}^{r} \sigma_{k i}(t, Y) D_{k} f(Y) d B^{i} \quad \text { on }[0, \tau[
$$

where

$$
L=\sum_{k=1}^{n} \beta_{k} D_{k}+\frac{1}{2} \sum_{k, \ell=1}^{n}\left(\sigma \sigma^{*}\right)_{k \ell} D_{k} D_{\ell}
$$

where $\sigma^{*}$ is the transpose of $\sigma$, and $\left(\sigma \sigma^{*}\right)_{k \ell}=\sum_{i=1}^{r} \sigma_{k i} \sigma_{\ell i}$. In particular, every solution of (1.2.2) is an L-diffusion on $[0, \tau[$ in the sense that

$$
d(f(Y))-L f(t, Y) d t=d(\text { local martingale }) \text { on }[0, \tau[
$$

Proof. (a) $\Rightarrow$ (b) Let $Y$ be a solution of SDE (1.2.2). Then

$$
d Y^{k} d Y^{\ell} \equiv d\left[Y^{k}, Y^{\ell}\right]=\left(\sigma \sigma^{*}\right)_{k \ell}(t, Y) d t
$$

where $\left[Y^{k}, Y^{\ell}\right]$ represents the quadratic covariation of $Y^{k}$ and $Y^{\ell}$. By Itô's formula we get

$$
\begin{aligned}
d(f(Y))= & \sum_{k=1}^{n} D_{k} f(Y)\left(\beta_{k}(t, Y) d t+\sum_{i=1}^{r} \sigma_{k i}(t, Y) d B^{i}\right) \\
& +\frac{1}{2} \sum_{k, \ell=1}^{n} D_{k} D_{\ell} f(Y) \underbrace{\left(\sigma \sigma^{*}\right)_{k \ell}(t, Y) d t}_{=d\left[Y^{k}, Y^{\ell}\right]}
\end{aligned}
$$

$$
\begin{aligned}
& =L f(t, Y) d t+\sum_{k=1}^{n} \sum_{i=1}^{r} \sigma_{k i}(t, Y) D_{k} f(t, Y) d B^{i} \\
& =L f(t, Y) d t+d(\text { local martingale })
\end{aligned}
$$

(b) $\Rightarrow$ (a) Take $f(x)=x_{\ell}$. Then $D_{k} f=\delta_{k \ell}$ and $L f=\beta_{\ell}$, thus

$$
d Y^{\ell}=\beta_{\ell}(t, Y) d t+\sum_{i=1}^{r} \sigma_{\ell i}(t, Y) d B^{i} \quad \text { for each } \ell=1, \ldots, n
$$

This shows that $Y$ solves $\operatorname{SDE}(1.2 .2)$ on $[0, \tau[$.
Proposition 1.2.4 (Itô SDE on $\mathbb{R}^{n}$; case of global Lipschitz conditions). Let $Z$ be a continuous semimartingale on $\mathbb{R}^{r}$ and

$$
\alpha: \mathbb{R}^{n} \rightarrow \operatorname{Hom}\left(\mathbb{R}^{r}, \mathbb{R}^{n}\right)(=\operatorname{Matr}(n \times r ; \mathbb{R}))
$$

such that

$$
\exists L>0, \quad|\alpha(y)-\alpha(z)| \leq L|y-z| \forall y, z \in \mathbb{R}^{n} \quad \text { (global Lipschitz conditions). }
$$

Then, for each $\mathscr{F}_{0}$-measurable $\mathbb{R}^{n}$-valued random variable $x_{0}$, there exists a unique continuous semimartingale $\left(X_{t}\right)_{t \geq 0}$ on $\mathbb{R}^{n}$ such that

$$
\begin{equation*}
d X=\alpha(X) d Z \text { and } X_{0}=x_{0} \tag{1.2.3}
\end{equation*}
$$

Uniqueness holds in the following sense: suppose that $Y$ is another continuous semimartingale such that $d Y=\alpha(Y) d Z$ and $Y_{0}=x_{0}$, then $X_{t}=Y_{t}$ for all $t$ a.s.

Proof. The proof is standard in Stochastic Analysis, see for instance [37] or [22].

Proposition 1.2.5 (Itô SDEs on $\mathbb{R}^{n}$ : case of the local Lipschitz coefficients). Let $Z$ be a continuous semimartingale on $\mathbb{R}^{r}$ and let

$$
\alpha: \mathbb{R}^{n} \rightarrow \operatorname{Hom}\left(\mathbb{R}^{r}, \mathbb{R}^{n}\right)
$$

be locally Lipschitz, i.e. for each compact $K \subset \mathbb{R}^{n}$ there exists a constant $L_{K}>0$ such that

$$
\forall y, z \in K, \quad|\alpha(y)-\alpha(z)| \leq L_{K}|y-z|
$$

Then, for any $x_{0} \mathscr{F}_{0}$-measurable, there exists a unique maximal solution $X \mid[0, \zeta[$ of the SDE

$$
d X=\alpha(X) d Z, \quad X_{0}=x_{0}
$$

Uniqueness holds in the sense that if $Y \mid\left[0, \xi\left[\right.\right.$ is another solution and $y_{0}=x_{0}$, then $\xi \leq \zeta$ a.s. and $X \mid[0, \xi[=Y$.

Proof. The proof is reduced to Proposition 1.2 .4 by a standard truncation method. We briefly sketch the argument, since it will be used several times in the sequel. Let $B(0, R)=\left\{x \in \mathbb{R}^{n}:|x| \leq R\right\}$ where $R=1,2, \ldots$ and choose test functions $\phi_{R} \in$ $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ such that $\phi_{R} \mid B(0, R) \equiv 1$. For $R>0$ consider the "truncated SDE"

$$
\begin{equation*}
d X^{R}=\alpha^{R}\left(X^{R}\right) d Z, \quad X_{0}^{R}=x_{0} \tag{1.2.4}
\end{equation*}
$$

where $\alpha^{R}:=\phi_{R} \alpha$ is now global Lipschitz. By Proposition 1.2.4 there is a unique solution $X^{R}$ to (1.2.4). Then

$$
X \mid\left[0, \tau_{R}\left[:=X^{R} \mid\left[0, \tau_{R}[\right.\right.\right.
$$

is well-defined by uniqueness, where

$$
\tau_{R}=\inf \left\{t \geq 0: X_{t}^{R} \notin B(0, R)\right\}
$$

This finally defines $X$ on the stochastic interval $\left[0, \zeta\left[\right.\right.$ where $\zeta=\sup _{R} \tau_{R}$. Uniqueness of $X$ is deduced from the uniqueness of $X \mid\left[0, \tau_{R}[\right.$.

Example 1.2.6. Consider the following Itô $\operatorname{SDE}$ on $\mathbb{R}^{n}$ :

$$
\begin{equation*}
d X=\underbrace{\beta(X)}_{n \times 1} d t+\underbrace{\sigma(X)}_{n \times r} \underbrace{d B}_{r \times 1} \tag{1.2.5}
\end{equation*}
$$

where $B$ is Brownian motion on $\mathbb{R}^{r}$. Then the space-time process $Z_{t}=\left(t, B_{t}\right)$ is a semimartingale on $\mathbb{R}^{r+1}$ and $\operatorname{SDE}(1.2 .5)$ can be written as

$$
d X=\binom{\beta(X)}{\sigma(X)}\binom{d t}{d B}=\alpha(X) d Z
$$

where $\alpha(X):=\binom{\beta(X)}{\sigma(X)}$. Thus, under a local Lipschitz condition on the coefficients $\beta$ and $\sigma$, the SDE

$$
\begin{equation*}
d X=\beta(X) d t+\sigma(X) d B \tag{1.2.6}
\end{equation*}
$$

has a unique strong solution for every given initial condition $x_{0}$. By Proposition 1.2.3, maximal solutions of Eq. (1.2.6) are $L$-diffusions to the operator

$$
L=\sum_{i=1}^{n} \beta_{i} \partial_{i}+\frac{1}{2} \sum_{i, j=1}^{n}\left(\sigma \sigma^{*}\right)_{i j} \partial_{i} \partial_{j}
$$

where $\partial_{i}=\partial / \partial x_{i}$ is the derivative in direction $i$.
Definition 1.2.7 (PDO in Hörmander form). For a vector field $A \in \Gamma(T M)$ on $M$ (read as a derivation) let $A^{2}(f):=A(A(f)), f \in C^{\infty}(M)$. A map $L: C^{\infty}(M) \rightarrow$ $C^{\infty}(M)$ is called a partial differential operator (PDO) in Hörmander form, if there exist vector fields $A_{0}, A_{1}, \ldots, A_{r}$ on $M$ such that $L$ can be written as

$$
L=A_{0}+\sum_{i=1}^{r} A_{i}^{2}
$$

In the special case $M=\mathbb{R}^{n}$ and $A_{i}:=D_{i}=\frac{\partial}{\partial x^{i}}(i=1, \ldots, n)$ for instance, $\Delta=\sum_{i=1}^{n} A_{i}^{2}$ is the Euclidean Laplacian.

### 1.2.2. Stratonovich differentials.

DEFINITION 1.2.8. For continuous real-valued semimartingales $X$ and $Y$ let

$$
X \circ d Y:=X d Y+\frac{1}{2} d[X, Y]
$$

be the Stratonovich differential. Here $X d Y$ is the usual Itô differential and $d[X, Y]=$ $d X d Y$ the differential of the quadratic covariation of $X$ and $Y$. The integral

$$
\begin{equation*}
\int_{0}^{t} X \circ d Y=\int_{0}^{t} X d Y+\frac{1}{2}[X, Y]_{t} \tag{1.2.7}
\end{equation*}
$$

is called Stratonovich integral of $X$ with respect to $Y$.
Formula (1.2.7) gives the relation between the Stratonovich integral and the usual Itô integral. Since Stratonovich integrals can always be converted back to Itô integrals, their use in our context will be only formal and for the sake of convenient notations.

REMARK 1.2.9. We have the following properties of Stratonovich differential, respectively Stratonovich integrals.

1. (Associativity) $X \circ(Y \circ d Z)=(X Y) \circ d Z$, i.e.,

$$
X \circ d\left(\int_{0}^{\bullet} Y \circ d Z\right)=(X Y) \circ d Z
$$

Indeed, we have

$$
\begin{aligned}
X \circ(Y \circ d Z) & =X \circ d\left(\int_{0}^{\bullet} Y \circ d Z\right) \\
& =X d\left(\int_{0}^{\bullet} Y \circ d Z\right)+\frac{1}{2} d X d\left(\int_{0}^{\bullet} Y \circ d Z\right) \\
& =X(Y d Z)+\frac{1}{2} X d Y d Z+\frac{1}{2} d X\left(Y d Z+\frac{1}{2} d Y d Z\right) \\
& =(X Y) d Z+\frac{1}{2}(X d Y+Y d X+d X d Y) d Z \\
& =(X Y) d Z+\frac{1}{2} d(X Y) d Z \\
& =(X Y) \circ d Z
\end{aligned}
$$

2. $($ Product rule $) d(X Y)=X \circ d Y+Y \circ d X$

Proof. By Itô's formula we have

$$
d(X Y)=X d Y+Y d X+d X d Y=X \circ d Y+Y \circ d X
$$

Proposition 1.2.10 (Itô-Stratonovich formula). Let $X$ be a continuous $\mathbb{R}^{n}$-valued semimartingale and $f \in C^{3}\left(\mathbb{R}^{n}\right)$. Then

$$
\begin{equation*}
d(f(X))=\sum_{i=1}^{n}\left(D_{i} f\right)(X) \circ d X^{i} \equiv\langle\nabla f(X), \circ d X\rangle \tag{1.2.8}
\end{equation*}
$$

Proof. By Itô's formula, we have

$$
d\left(D_{i} f(X)\right)=\sum_{k=1}^{n}\left(D_{i} D_{k} f\right)(X) d X^{k}+\frac{1}{2} \sum_{k, \ell=1}^{n}\left(D_{i} D_{k} D_{\ell} f\right)(X) d X^{k} d X^{\ell}
$$

Hence we get

$$
\begin{aligned}
\sum_{i=1}^{n}\left(D_{i} f\right)(X) \circ d X^{i} & =\sum_{i=1}^{n}\left(D_{i} f\right)(X) d X^{i}+\frac{1}{2} \sum_{i=1}^{n} d\left(D_{i} f(X)\right) d X^{i} \\
& =\sum_{i=1}^{n}\left(D_{i} f\right)(X) d X^{i}+\frac{1}{2} \sum_{i, k=1}^{n}\left(D_{i} D_{k} f(X)\right) d X^{k} d X^{i} \\
& =d(f(X))
\end{aligned}
$$

Formula (1.2.8) shows the main advantage of the Stratonovich differential: it converts Itô's formula into the usual chain rule of classical analysis. Hence, at least formally, classical differential calculus can be applied in calculations involving Stratonovich differentials.

PROPOSITION 1.2.11. Let $\beta: \mathbb{R}_{+} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be continuous, $\sigma: \mathbb{R}_{+} \times \mathbb{R}^{n} \rightarrow$ $\operatorname{Hom}\left(\mathbb{R}^{r}, \mathbb{R}^{n}\right)$ be $C^{1}$. Furthermore, let $B$ be a Brownian motion on $\mathbb{R}^{r}$. For a semimartingale $Y$ on $\mathbb{R}^{n}$ (defined up to some predictable stopping time $\tau$ ) the following conditions are equivalent:
(i) The semimartingale $Y$ is a solution of the Stratonovich SDE

$$
\begin{equation*}
d Y=\beta(t, Y) d t+\sigma(t, Y) \circ d B \tag{1.2.9}
\end{equation*}
$$

i.e.

$$
Y_{t}=Y_{0}+\int_{0}^{t} \beta\left(s, Y_{s}\right) d s+\int_{0}^{t} \sigma\left(s, Y_{s}\right) \circ d B_{s}, \quad \text { for } 0 \leq t<\tau \text { a.s. }
$$

(ii) For all $f \in C^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
d(f(Y))=(L f)(t, Y) d t+\sum_{k=1}^{r}\left(A_{k} f\right)(t, Y) d B^{k} \quad \text { on }[0, \tau[
$$

where

$$
L=A_{0}+\frac{1}{2} \sum_{k=1}^{r} A_{k}^{2}
$$

with the vector fields $A_{i} \in \Gamma\left(T \mathbb{R}^{n}\right)$ defined as

$$
\begin{equation*}
A_{0}=\sum_{i=1}^{n} \beta_{i} D_{i}, \quad A_{k}=\sum_{i=1}^{n} \sigma_{i k} D_{i}, \quad k=1, \ldots, r . \tag{1.2.10}
\end{equation*}
$$

Proof. (i) $\Rightarrow$ (ii) By the Itô-Stratonovich formula (Proposition 1.2.10) we have

$$
\begin{aligned}
d(f(Y)) & =\sum_{i=1}^{n}\left(D_{i} f\right)(Y) \circ d Y^{i} \\
& =\sum_{i=1}^{n}\left(D_{i} f\right)(Y) \beta_{i}(t, Y) d t+\sum_{i=1}^{n}\left(D_{i} f\right)(Y)\left(\sum_{k=1}^{r} \sigma_{i k}(t, Y) \circ d B^{k}\right) \\
& =\left(A_{0} f\right)(t, Y) d t+\sum_{k=1}^{r}\left(A_{k} f\right)(t, Y) \circ d B^{k} \\
& =\left(A_{0} f\right)(t, Y) d t+\sum_{k=1}^{r}\left(A_{k} f\right)(t, Y) d B^{k}+\frac{1}{2} \sum_{k=1}^{r} d\left(\left(A_{k} f\right)(t, Y)\right) d B^{k}
\end{aligned}
$$

Since

$$
d\left(A_{k} f(t, Y)\right)=\partial_{t}\left(A_{k} f\right)(t, Y) d t+\left(A_{0} A_{k} f\right)(t, Y) d t+\sum_{\ell=1}^{r}\left(A_{\ell} A_{k} f\right)(t, Y) \circ d B^{\ell}
$$

we observe that

$$
d\left(A_{k} f(t, Y)\right) d B^{k}=\left(A_{k}^{2} f\right)(t, Y) d t
$$

and hence

$$
d(f(Y))=\underbrace{\left(\left(A_{0} f\right)(t, Y)+\frac{1}{2} \sum_{k=1}^{r}\left(A_{k}^{2} f\right)(t, Y)\right)}_{=(L f)(t, Y)} d t+\sum_{k=1}^{r}\left(A_{k} f\right)(t, Y) d B^{k}
$$

(ii) $\Rightarrow$ (i) It is sufficient to take $f(x)=x_{\ell}$.

Corollary 1.2.12. Solutions to the Stratonovich SDE

$$
d Y=\beta(t, Y) d t+\sigma(t, Y) \circ d B
$$

define L-diffusions for the operator

$$
L=A_{0}+\frac{1}{2} \sum_{i=1}^{r} A_{i}^{2} \quad \text { with } A_{0}, A_{1}, \ldots, A_{r} \text { as in Eq. (1.2.10), }
$$

in the sense that

$$
d(f \circ Y)-(L f)(t, Y) d t \stackrel{\mathrm{~m}}{=} 0
$$

for all $f \in C^{\infty}\left(\mathbb{R}^{n}\right)$.
1.2.3. Stochastic differential equations on manifolds. In this section we describe the construction of $L$-diffusions as solutions of stochastic differential equations on manifolds [11, 17].

Definition 1.2.13 (Stochastic differential equation on $M$ ). Let $M$ be a differentiable manifold, $\pi: T M \rightarrow M$ its tangent bundle and $E$ a finite dimensional vector space (without restrictions $E=\mathbb{R}^{r}$ ). A stochastic differential equation on $M$ is a pair $(A, Z)$ where

1. $Z$ is a semimartingale taking values in $E$;
2. $A: M \times E \rightarrow T M$ is a smooth homomorphism of vector bundles over $M$, i.e.

$$
(x, e) \longmapsto A(x) e:=A(x, e)
$$



REMARK 1.2.14. Formally the homomorphism $A$ may be considered as section $A \in$ $\Gamma\left(E^{*} \otimes T M\right)$. In particular, we have

$$
\begin{cases}\forall x \in M, & A(x) \in \operatorname{Hom}\left(E, T_{x} M\right), \\ \forall e \in E, & A(\cdot) e \in \Gamma(T M)\end{cases}
$$

Notation 1.2.15. For the $\operatorname{SDE}(A, Z)$ we also write

$$
d X=A(X) \circ d Z
$$

or

$$
d X=\sum_{i=1}^{r} A_{i}(X) \circ d Z^{i}
$$

where $A_{i}=A(\cdot) e_{i} \in \Gamma(T M)$ and $e_{1}, \ldots, e_{r}$ is a basis of $E$.
DEFINITION 1.2.16 (Solution of a stochastic differential equation). Let $(A, Z)$ be an SDE on $M$ and let $x_{0}: \Omega \rightarrow M$ be $\mathscr{F}_{0}$-measurable. An adapted continuous process $X \mid\left[0, \zeta\left[\equiv\left(X_{t}\right)_{t<\zeta}\right.\right.$ taking values in $M$, defined up to the stopping time $\zeta$, is called solution to the SDE

$$
\begin{equation*}
d X=A(X) \circ d Z \tag{1.2.11}
\end{equation*}
$$

with initial condition $X_{0}=x_{0}$, if for all $f \in C_{c}^{\infty}(M)$ the following conditions are satisfied:
(i) $f(X)$ is a semimartingale;
(ii) for any stopping time $\tau$ such that $0 \leq \tau<\zeta$, we have

$$
\begin{equation*}
f\left(X_{\tau}\right)=f\left(X_{0}\right)+\int_{0}^{\tau}(d f)_{X_{s}} A\left(X_{s}\right) \circ d Z_{s} \tag{1.2.12}
\end{equation*}
$$

We call $X$ maximal solution of the $\operatorname{SDE}(1.2 .11)$ if

$$
\{\zeta<\infty\} \subset\left\{\lim _{t \uparrow \zeta} X_{t}=\infty \text { in } \hat{M}=M \dot{\cup}\{\infty\}\right\} \text { a.s. }
$$

Note: The integral in (1.2.12) is defined using the linear functional

$$
E \xrightarrow{A(x)} T_{x} M \xrightarrow{(d f)_{x}} \mathbb{R}, \quad x \in M
$$

REMARK 1.2.17. We adopt the convention $X_{t}(\omega):=\infty$ for $\zeta(\omega) \leq t<\infty$ and $f(\infty)=0$ for $f \in C_{c}^{\infty}(M)$. Then we may write, for all $t \geq 0$,

$$
\begin{aligned}
f\left(X_{t}\right) & =f\left(X_{0}\right)+\int_{0}^{t}(d f)_{X_{s}} A\left(X_{s}\right) \circ d Z_{s} \\
& =f\left(X_{0}\right)+\sum_{i=1}^{r} \int_{0}^{t}(d f)_{X_{s}} A_{i}\left(X_{s}\right) \circ d Z_{s}^{i} \\
& =f\left(X_{0}\right)+\sum_{i=1}^{r} \int_{0}^{t}\left(A_{i} f\right)\left(X_{s}\right) \circ d Z_{s}^{i} \quad \text { with } A_{i}=A(\cdot) e_{i}
\end{aligned}
$$

EXAMPLE 1.2.18. Let $E=\mathbb{R}^{r+1}$ and $Z=\left(t, Z^{1}, \ldots, Z^{r}\right)$ where $\left(Z^{1}, \ldots, Z^{r}\right)$ is a Brownian motion on $\mathbb{R}^{r}$. Denote the standard basis of $\mathbb{R}^{r+1}$ by $\left(e_{0}, e_{1}, \ldots, e_{r}\right)$. Letting

$$
A: M \times E \rightarrow T M
$$

be a homomorphism of vector bundles over $M$, we consider the vector fields

$$
A_{i}:=A(\cdot) e_{i} \in \Gamma(T M), \quad i=0,1, \ldots, r .
$$

Then the SDE

$$
\begin{equation*}
d X=A(X) \circ d Z \tag{1.2.13}
\end{equation*}
$$

writes as

$$
d X=A_{0}(X) d t+\sum_{i=1}^{r} A_{i}(X) \circ d Z^{i}
$$

and for each $f \in C_{c}^{\infty}(M)$ we have

$$
\begin{aligned}
d(f(X)) & =(d f)_{X} A(X) \circ d Z \\
& =\sum_{i=0}^{r}(d f)_{X} A(X) e_{i} \circ d Z^{i} \\
& =\sum_{i=0}^{r}(d f)_{X} A_{i}(X) \circ d Z^{i} \\
& =\sum_{i=0}^{r}\left(A_{i} f\right)(X) \circ d Z^{i} \\
& =\left(A_{0} f\right)(X) d t+\sum_{i=1}^{r}\left(A_{i} f\right)(X) \circ d Z^{i}
\end{aligned}
$$

$$
=\left(A_{0} f\right)(X) d t+\sum_{i=1}^{r}\left(\left(A_{i} f\right)(X) d Z^{i}+\frac{1}{2} d\left(\left(A_{i} f\right)(X)\right) d Z^{i}\right)
$$

Taking into account that

$$
d\left(\left(A_{i} f\right)(X)\right)=\sum_{j=1}^{r}\left(A_{j} A_{i} f\right)(X) d Z^{j}+d(\text { terms of bounded variation })
$$

we see that

$$
d\left(\left(A_{i} f\right)(X)\right) d Z^{i}=\left(A_{i}^{2} f\right)(X) d t
$$

where we used that $d Z^{i} d Z^{j}=\delta_{i j} d t$ for $1 \leq i, j \leq r$. Hence we get

$$
\begin{aligned}
d(f(X)) & =\left(A_{0} f\right)(X) d t+\frac{1}{2} \sum_{j=1}^{r}\left(A_{i}^{2} f\right)(X) d t+\sum_{i=1}^{r}\left(A_{i} f\right)(X) d Z^{i} \\
& =(L f)(X) d t+\sum_{i=1}^{r}\left(A_{i} f\right)(X) d Z^{i}
\end{aligned}
$$

Corollary 1.2.19. Let $L=A_{0}+\frac{1}{2} \sum_{i=1}^{r} A_{i}^{2}$ and let $X$ be a solution to Eq. (1.2.13). Then, for all $f \in C_{c}^{\infty}(M)$,

$$
d(f(X))-(L f)(X) d t \stackrel{\mathrm{~m}}{=} 0
$$

where $\stackrel{m}{=}$ denotes equality modulo differentials of martingales. In other words, maximal solutions to the SDE

$$
d X=A(X) \circ d Z
$$

are $L$-diffusions to the operator $L=A_{0}+\frac{1}{2} \sum_{i=1}^{r} A_{i}^{2}$.
THEOREM 1.2.20 (SDE: Existence and uniqueness of solutions; $M=\mathbb{R}^{n}$ ). Let $(A, Z)$ be an SDE on $M=\mathbb{R}^{n}$ and $x_{0}$ an $\mathscr{F}_{0}$-measurable random variable taking values in $\mathbb{R}^{n}$. Then there exists a unique maximal solution $X$ (with maximal lifetime $\zeta>0$ a.s.) of the SDE

$$
\begin{equation*}
d X=A(X) \circ d Z \tag{1.2.14}
\end{equation*}
$$

with initial condition $X_{0}=x_{0}$. Uniqueness holds in the following sense: if $Y \mid[0, \xi[$ is another solution of (1.2.14) to the same initial condition, then $\xi \leq \zeta$ a.s. and $X \mid[0, \xi[=Y$ a.s.

Proof. As in the proof of Proposition 1.2.5 let $B(0, R)=\left\{x \in \mathbb{R}^{n}:|x| \leq R\right\}$ where $R=1,2, \ldots$ and choose test functions $\phi_{R} \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ such that $\phi_{R} \mid B(0, R) \equiv 1$. Since

$$
A \in \Gamma\left(\operatorname{Hom}\left(\mathbb{R}^{r}, T M\right)\right)
$$

we have for each $x \in \mathbb{R}^{n}$ the linear map

$$
A(x): \mathbb{R}^{r} \rightarrow T_{x} M
$$

In this way $A$ gives rise to a smooth map $\mathbb{R}^{n} \rightarrow \operatorname{Matr}(n \times r ; \mathbb{R})$.
Consider now the "truncated SDE"

$$
\begin{equation*}
d X^{R}=A^{R}\left(X^{R}\right) \circ d Z \tag{1.2.15}
\end{equation*}
$$

where $A^{R}=\phi_{R} A$. By Proposition 1.2.4, the truncated $\operatorname{SDE}$ (1.2.15) has a unique global solution $X^{R}$ with initial condition $X_{0}^{R}=x_{0}$, i.e., for each $R$ there exists a continuous
$\mathbb{R}^{n}$-valued semimartingale $\left(X_{t}^{R}\right)_{t \geq 0}$ satisfying $X_{0}^{R}=x_{0}$ such that (1.2.15) holds in the Itô-Stratonovich sense. In terms of the stopping times

$$
\tau_{R}:=\inf \left\{t \geq 0: X_{t}^{R} \notin B(0, R)\right\}
$$

we have for $R<R^{\prime}$,

$$
X^{R^{\prime}} \mid\left[0, \tau_{R}\left[=X^{R} \mid\left[0, \tau_{R}[\quad \text { a.s. }\right.\right.\right.
$$

Hence a stochastic process $X$ (with lifetime $\zeta=\lim _{R \uparrow \infty} \tau_{R}$ ) is well-defined via

$$
X \mid\left[0, \tau_{R}\left[=X^{R} \mid\left[0, \tau_{R}[.\right.\right.\right.
$$

For each $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ such that $\operatorname{supp}(f) \subset B(0, R)$ (with $R$ sufficiently large), we have

$$
\begin{aligned}
d(f(X)) & =d\left(f\left(X^{R}\right)\right) \\
& =\sum_{k=1}^{n}\left(D_{k} f\left(X^{R}\right)\right) \circ d\left(X^{R}\right)^{k} \quad \text { (using Itô-Stratonovich formula) } \\
& =\left\langle\nabla f\left(X^{R}\right), \circ d X^{R}\right\rangle \\
& =\left\langle\nabla f\left(X^{R}\right), \phi_{R}\left(X^{R}\right) A\left(X^{R}\right) \circ d Z\right\rangle \\
& =\langle\nabla f(X), A(X) \circ d Z\rangle \\
& =\sum_{i=1}^{r}\left\langle\nabla f(X), A_{i}(X) \circ d Z^{i}\right\rangle \\
& =\sum_{i=1}^{r}(d f)_{X} A_{i}(X) \circ d Z^{i} \\
& =(d f)_{X} A(X) \circ d Z
\end{aligned}
$$

Hence, $X$ is the unique solution to Eq. (1.2.14) with initial condition $X_{0}=x_{0}$. Note that $X$ is a solution of $d X=A(X) \circ d Z$ in the Itô-Stratonovich sense (in $\mathbb{R}^{n}$ ) if and only if $\forall f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
d(f(X))=(d f)_{X} A(X) \circ d Z
$$

THEOREM 1.2.21 (SDE: Existence and uniqueness of solutions; general case). Let $(A, Z)$ be an SDE on a differentiable manifold $M$ and let $x_{0}: \Omega \rightarrow M$ be $\mathscr{F}_{0}$-measurable . There exists a unique maximal solution $X \mid[0, \zeta[$ (where $\zeta>0$ a.s.) of the SDE

$$
d X=A(X) \circ d Z
$$

with initial condition $X_{0}=x_{0}$. Uniqueness holds in the sense that if $Y \mid[0, \xi[$ is another solution with $Y_{0}=x_{0}$, then $\xi \leq \zeta$ a.s. and $X \mid[0, \xi[=Y$ a.s.

We shall reduce Theorem 1.2.21 to Theorem 1.2 .20 via embedding the manifold $M$ into a high-dimensional Euclidean space.

Whitney's embedding theorem. Each manifold $M$ of dimension $n$ can be embedded into $\mathbb{R}^{n+k}$ as a closed submanifold (for $k$ sufficiently large, e.g. $k=n+1$ ), i.e.,

$$
M \hookrightarrow \iota(M) \subset \mathbb{R}^{n+k}
$$

where $\iota: M \rightarrow \iota(M)$ is a diffeomorphism and $\iota(M) \subset \mathbb{R}^{n+k}$ a closed submanifold.
Proof (of Theorem 1.2.21). We choose a Whitney embedding (in general not intrinsic)

$$
M \underset{\text { diffeom. }}{\stackrel{\iota}{\longrightarrow}} \iota(M) \subset \mathbb{R}^{n+k}
$$

and identify $M$ and $\iota(M)$; in particular for each $x \in M$ the tangent space $T_{x} M$ is then a linear subspace of $\mathbb{R}^{n+k}$ according to

$$
T_{x} M \stackrel{\text { 抆 }}{\longrightarrow} T_{x} \mathbb{R}^{n+k} \equiv \mathbb{R}^{n+k}
$$

Vector fields $A_{1}, \ldots, A_{r} \in \Gamma(T M)$ can be extended to vector fields

$$
\bar{A}_{1}, \ldots, \bar{A}_{r} \in \Gamma\left(T \mathbb{R}^{n+k}\right) \equiv C^{\infty}\left(\mathbb{R}^{n+k} ; \mathbb{R}^{n+k}\right) \quad \text { with } \bar{A}_{i} \mid M=A_{i}
$$

i.e. $\bar{A}_{i} \circ \iota=d \iota \circ A_{i}$. Hence a given bundle map

$$
A: M \times \mathbb{R}^{r} \rightarrow T M, \quad(x, z) \mapsto A(x) z=\sum_{i=1}^{r} A_{i}(x) z^{i}
$$

has a continuation

$$
\bar{A}: \mathbb{R}^{n+k} \times \mathbb{R}^{r} \rightarrow \mathbb{R}^{n+k} \times \mathbb{R}^{n+k}, \quad(x, z) \mapsto \bar{A}(x) z=\sum_{i=1}^{r} \bar{A}_{i}(x) z^{i}
$$

The idea is to consider in place of the original SDE

$$
\begin{equation*}
d X=A(X) \circ d Z \text { on } M \tag{*}
\end{equation*}
$$

the SDE

$$
\begin{equation*}
d X=\bar{A}(X) \circ d Z \text { on } \mathbb{R}^{n+k} \tag{*}
\end{equation*}
$$

First of all it is clear that any solution of $(*)$ in $M$ provides a solution of $(\bar{*})$ in $\mathbb{R}^{n+k}$. More precisely: If $X$ is a solution to $(*)$ with starting value $X_{0}=x_{0}$, then $\bar{X}:=\iota \circ X$ solves equation ( $\bar{*}$ ) with starting value $\bar{X}_{0}=\iota \circ x_{0}$. Indeed if $\bar{f} \in C_{c}^{\infty}\left(\mathbb{R}^{n+k}\right)$, then $f:=\bar{f} \mid M=\bar{f} \circ \iota \in C_{c}^{\infty}(M)$, and we have:

$$
\begin{aligned}
d(\bar{f}(\bar{X}))=d(f(X)) & =\sum_{i=1}^{r}(d f)_{X} A_{i}(X) \circ d Z^{i} \\
& =\sum_{i=1}^{r}(d \bar{f})_{\bar{X}}(d \iota)_{X} A_{i}(X) \circ d Z^{i} \\
& =\sum_{i=1}^{r}(d \bar{f})_{\bar{X}} \bar{A}_{i}(\iota \circ X) \circ d Z^{i} \\
& =\sum_{i=1}^{r}(d \bar{f})_{\bar{X}} \bar{A}_{i}(\bar{X}) \circ d Z^{i} .
\end{aligned}
$$

This implies in particular uniqueness of solutions to $(*)$, since equation $(\bar{*})$ has a unique solution to a given initial condition.

To establish existence of solutions to $(*)$ we first remark that any test function $f \in$ $C_{c}^{\infty}(M)$ has a continuation $\bar{f} \in C_{c}^{\infty}\left(\mathbb{R}^{n+k}\right)$ such that $\bar{f} \mid M \equiv \bar{f} \circ \iota=f$. We make the following important observation.

Each solution $X \mid\left[0, \zeta\left[\right.\right.$ of $(\bar{*})$ in $\mathbb{R}^{n+k}$ with $X_{0}=x_{0}$ which stays on $M$ for $t<\zeta$ (where $x_{0}$ is an $M$-valued $\mathscr{F}_{0}$-measurable random variable) gives a solution of $(*)$.

Hence, to complete the proof it is sufficient to show the following lemma.
Lemma 1.2.22. If $X \mid\left[0, \zeta\left[\right.\right.$ is the maximal solution of $(\bar{*})$ in $\mathbb{R}^{n+k}$ with $X_{0}=x_{0}$, then

$$
\{t<\zeta\} \subset\left\{X_{t} \in M\right\}, \quad \text { for all } t \text { a.s. }
$$

Observe that it is enough to verify Lemma 1.2.22 for one specific continuation $\bar{A}$ of $A$.

Proof (of Lemma 1.2.22). Let

$$
\perp M=\left\{(x, v) \in M \times \mathbb{R}^{n+k} \mid v \in\left(T_{x} M\right)^{\perp}\right\}
$$

be the normal bundle of $M$ and consider $M$ embedded into $\perp M$ as zero section:

$$
M \hookrightarrow \perp M, \quad x \mapsto(x, 0)
$$



Figure 1.2.1. Normal bundle $\perp M$
Fact: There is a smooth function $\varepsilon: M \rightarrow] 0, \infty[$ such that the map

$$
\begin{aligned}
\tau_{\varepsilon}(M):=\{(x, v) \in \perp M:|v|<\varepsilon(x)\} & \cong \\
(x, v) & \longmapsto x+v
\end{aligned}
$$

is a diffeomorphism from the tubular neighbourhood $\tau_{\varepsilon}(M)$ of $M$ of radius $\varepsilon$ onto the indicated part in $\mathbb{R}^{n+k}$. This follows from the local inversion theorem since the given map has full rank along the zero section of $\perp M$.

Note that both

$$
\left.\begin{array}{rl}
\pi: \tau_{\varepsilon}(M) & \rightarrow M, \quad(x, v) \\
\operatorname{dist}^{2}(\cdot, M): \tau_{\varepsilon}(M) & \rightarrow \mathbb{R}, \quad(x, v)
\end{array}\right)|v|^{2},
$$

are smooth maps.
Now letting $R>0$ be sufficiently large such that

$$
M \cap B(0, R+1) \neq \varnothing
$$

then

$$
\varepsilon_{R}=\inf \{\varepsilon(x) \mid x \in M \cap B(0, R+1)\}>0
$$

We choose a decreasing smooth function $\lambda:[0, \infty[\rightarrow[0,1]$ of the form


Figure 1.2.2. Cut-off function $\lambda$
and a test function $0 \leq \varphi \in C_{c}^{\infty}\left(\mathbb{R}^{n+k}\right)$ such that $\varphi \mid B(0, R) \equiv 1$ and $\operatorname{supp}(\varphi) \subset$ $B(0, R+1)$. Consider the map

$$
\begin{aligned}
& \bar{A}^{R}: \mathbb{R}^{n+k} \times \mathbb{R}^{r} \rightarrow \mathbb{R}^{n+k} \times \mathbb{R}^{n+k}, \\
& \bar{A}^{R}(y, z):= \begin{cases}\varphi(y) \lambda\left(\operatorname{dist}^{2}(y, M)\right) A(\pi(y)) z & \text { if } y \in \tau_{\varepsilon}(M), \\
0 & \text { if } y \notin \tau_{\varepsilon}(M)\end{cases}
\end{aligned}
$$



Figure 1.2.3. Extended coefficients of the SDE
Let $X$ be the solution of

$$
d X=\bar{A}^{R}(X) \circ d Z, \quad X_{0}=x_{0}
$$

Consider the test function $f \in C_{c}^{\infty}\left(\mathbb{R}^{n+k}\right)$ given as

$$
f(y)=\varphi(y) \lambda\left(\operatorname{dist}^{2}(y, M)\right)
$$

Then

$$
\begin{aligned}
d(f(X)) & =(d f)_{X} \bar{A}^{R}(X) \circ d Z \\
& =\left\langle\nabla f(X), \bar{A}^{R}(X) \circ d Z\right\rangle \\
& =0 \quad \text { on }\left[0, \tau_{R}[,\right.
\end{aligned}
$$

where $\tau_{R}:=\inf \left\{t \geq 0: X_{t} \notin B(0, R)\right\}$. Indeed, $f$ is constant on each submanifold of the form

$$
\{\operatorname{dist}(\cdot, M)=s\} \cap B(0, R), \quad s<\varepsilon_{R},
$$

whereas $\bar{A}^{R}(y, z)$ is tangent to such submanifolds. Thus, for all $y \in B(0, R)$ and $z \in \mathbb{R}^{r}$,

$$
\nabla f(y) \perp \bar{A}^{R}(y) z
$$

Hence, for any solution $X$ of $(\bar{*})$, we obtain that

$$
f(X) \equiv \text { constant on }\left[0, \tau_{R}[\text { a.s. }\right.
$$

Since $R$ is arbitrary, this completes the proof of the Lemma.
Solutions to an SDE on $M$ of the type (1.2.11) are by definition semimartingales on $M$ as defined above: A continuous adapted process $X$ with values in $M$ is a semimartingale on $M$ if, for each $f \in C_{c}^{\infty}(M)$, the composition $f \circ X$ provides a continuous real-valued
semimartingale. It is easy to see that each $M$-valued semimartingale can be obtained as solution of an SDE on $M$.

THEOREM 1.2.23 (Manifold-valued semimartingales as solutions of an SDE). Every semimartingale on a manifold $M$ is given as solution of an SDE of the type (1.2.11).

Proof. Let $X$ be an arbitrary semimartingale on $M$. Without loss of generality (after an eventual change of time), we may assume that $X$ has infinite lifetime. Choosing a Whitney embedding $\iota: M \hookrightarrow \mathbb{R}^{n+k}$ we may consider the semimartingale $Z:=\iota \circ X$ taking values in $E:=\mathbb{R}^{n+k}$. Let $A: M \times E \rightarrow T M$ be the bundle homomorphism which is fiberwise the orthogonal projection $A(x): \mathbb{R}^{n+k} \rightarrow T_{x} M$ of $\mathbb{R}^{n+k}$ onto $T_{x} M \subset$ $T_{x} \mathbb{R}^{n+k}=\mathbb{R}^{n+k}$. We show that $X$ solves the equation

$$
d X=A(X) \circ d Z
$$

Let $f \in C_{c}^{\infty}(M)$ be given. We choose a continuation $\bar{f} \in C_{c}^{\infty}\left(\mathbb{R}^{n+k}\right)$ where $\bar{f} \circ \iota=f$ such that $\bar{f}$ is constant locally about $M$ on the normal subspaces $\perp_{x} M$ (this is $\bar{f}(y)=f(x)$ for $y \in \perp_{x} M$ sufficiently small). Now let $x \in M$ and $z \in \mathbb{R}^{n+k}$. By decomposing $z=z_{0}+z^{\perp}$ where $z_{0} \in T_{x} M$ and $z^{\perp} \in \perp_{x} M$, we obtain:

$$
(d f)_{x} A(x) z=(d \bar{f})_{\iota(x)}(d \iota)_{x} A(x) z=(d \bar{f})_{\iota(x)} z_{0}=(d \bar{f})_{\iota(x)} z
$$

But then

$$
\begin{aligned}
d(f(X)) & =d(\bar{f}(\iota(X)))=\sum_{i=1}^{n+k}\left(D_{i} \bar{f}\right)(\iota(X)) \circ d Z^{i} \\
& =\sum_{i=1}^{n+k}(d f)_{X} A(X) e_{i} \circ d Z^{i}=(d f)_{X} A(X) \circ d Z
\end{aligned}
$$

which gives the claim.
Remark 1.2.24. Let $M$ and $N$ be differentiable manifolds. For semimartingales $X$ on $M$, respectively $X^{\prime}$ on $N$, both adapted to the same filtration, consider the product semimartingale $\tilde{X}:=\left(X, X^{\prime}\right)$ taking values in $M \times N$. Suppose that

$$
\begin{equation*}
d X=A(X) \circ d Z, \quad \text { resp. } \quad d X^{\prime}=A^{\prime}\left(X^{\prime}\right) \circ d Z^{\prime} \tag{1.2.16}
\end{equation*}
$$

with bundle maps $A: M \times \mathbb{R}^{k} \rightarrow T M$ over $M$, respectively $A^{\prime}: N \times \mathbb{R}^{k^{\prime}} \rightarrow T N$ over $N$. Then $\tilde{X}$ solves the "composed" SDE

$$
\begin{equation*}
d \tilde{X}=\tilde{A}(\tilde{X}) \circ d \tilde{Z} \tag{1.2.17}
\end{equation*}
$$

driven by the $\mathbb{R}^{k} \times \mathbb{R}^{k^{\prime}}$-valued semimartingale $\tilde{Z}:=\left(Z, Z^{\prime}\right)$ where

$$
\tilde{A}\left(x, x^{\prime}\right)\left(z, z^{\prime}\right):=\left(A(x) z, A^{\prime}\left(x^{\prime}\right) z^{\prime}\right) \in T_{x} M \oplus T_{x^{\prime}} N \equiv T_{\left(x, x^{\prime}\right)}(M \times N)
$$

defines a bundle map $\tilde{A}:(M \times N) \times\left(\mathbb{R}^{k} \times \mathbb{R}^{k^{\prime}}\right) \rightarrow T(M \times N)$ over $M \times N$.
Proof. Let $\iota: M \hookrightarrow \mathbb{R}^{\ell}$ and $\iota^{\prime}: N \hookrightarrow \mathbb{R}^{\ell^{\prime}}$ be Whitney embeddings. Any function $f \in C^{\infty}(M \times N)$ factorizes as $f=\bar{f} \circ\left(\iota, \iota^{\prime}\right)$ for some $\bar{f} \in C^{\infty}\left(\mathbb{R}^{\ell} \times \mathbb{R}^{\ell^{\prime}}\right)$. Let $\bar{X}=\iota(X)$ and $\bar{X}^{\prime}=\iota^{\prime}\left(X^{\prime}\right)$. Then for $f \in C^{\infty}(M \times N)$, the semimartingale $f(\tilde{X})=\bar{f}\left(\bar{X}, \bar{X}^{\prime}\right)$ satisfies

$$
\begin{aligned}
d(f(\tilde{X})) & =d\left(\bar{f}\left(\bar{X}, \bar{X}^{\prime}\right)\right)=(d \bar{f})\left(\bar{X}, \bar{X}^{\prime}\right) \circ d\left(\bar{X}, \bar{X}^{\prime}\right) \\
& =(d \bar{f})\left(\bar{X}, \bar{X}^{\prime}\right) \circ\left(d \bar{X}, d \bar{X}^{\prime}\right)=(d \bar{f})\left(\bar{X}, \bar{X}^{\prime}\right) \circ\left(d(\iota(X)), d\left(\iota\left(X^{\prime}\right)\right)\right) \\
& =\bar{f}_{*}\left(\iota_{*} A(X) \circ d Z, \iota_{*}^{\prime} A^{\prime}\left(X^{\prime}\right) \circ d Z^{\prime}\right)=\left(\bar{f}_{*}\left(\iota, \iota^{\prime}\right)_{*}\right) \tilde{A}(\tilde{X}) \circ d \tilde{Z}
\end{aligned}
$$

$$
=f_{*} \tilde{A}(\tilde{X}) \circ d \tilde{Z} \equiv \sum_{i} f_{*} \tilde{A}_{i}(\tilde{X}) \circ d \tilde{Z}^{i}
$$

which proves the claim.

### 1.3. Quadratic Variation and Integration of one-forms

In this section we give canonical constructions related to continuous semimartingales on a manifold $M$, including the quadratic variation of continuous semimartingales with respect to bilinear forms on $T M$ and the integral of one-forms on $M$ along semimartingales, see [12] for more details. In the particular case $M=\mathbb{R}^{n}$ endowed with the Euclidean metric this notion of the quadratic variation reduces to the usual quadratic variation of a semimartingale.

Both notions (quadratic variation and integration of one-forms) can be deduced from a unified construction principle within the framework of second order differential geometry. We postpone this point of view and develop the theory first only as far as needed for martingale theory on manifolds.

We start with an elementary technical lemma on continuous processes, which is quite useful as it allows a spatial localization of continuous adapted processes, besides the usual localization in time through a localizing sequence of stopping times. The lemma basically reduces to properties of continuous paths.

Lemma 1.3.1. Let $\left(V_{k}\right)_{k \in \mathbb{N}}$ be a countable covering of $M$ by open sets $V_{k}$ and $X$ be a continuous adapted $M$-valued process. Then there exists a non-decreasing sequence $\left(\tau_{n}\right)_{n \geq 0}$ of stopping times with $\tau_{0}=0$ and $\sup _{n} \tau_{n}=\infty$, such that on each of the intervals $\left[\tau_{n}, \tau_{n+1}\right] \cap\left(\mathbb{R}_{+} \times\left\{\tau_{n}<\tau_{n+1}\right\}\right)$ the process $X$ takes values only in one of the $V_{k}$.

Proof of the Lemma. First of all, we choose a refinement $\left(W_{k}\right)_{k \in \mathbb{N}}$ to $\left(V_{k}\right)_{k \in \mathbb{N}}$ such that for each $k \in \mathbb{N}$ the closure $\bar{W}_{k}$ of $W_{k}$ is still contained in one of the $V_{n(k)}$. We construct a sequence $\left(\tau_{n}^{k}\right)_{0 \leq k \leq n, n \geq 0}$ of stopping times which after a suitable renumbering will satisfy the claimed assertions. Let $\tau_{0}^{0}:=0$. Suppose that $\tau_{n}^{k}$ is already constructed up to a certain $n$, then let

$$
\tau_{n+1}^{0}:=\tau_{n}^{n}, \text { and } \tau_{n+1}^{k}:=\inf \left\{t \geq \tau_{n+1}^{k-1}: X_{t} \notin W_{k}\right\} \quad \text { for } k=1, \ldots, n+1
$$

It remains to verify that $\sup _{n \geq 0} \sup _{k \leq n} \tau_{n}^{k}=\infty$. Let's suppose that there exists $\omega \in \Omega$ such that $t_{0}:=\sup _{n \geq 0} \sup _{k \leq n} \tau_{n}^{k}(\omega)<\infty$. Then we know $X_{t_{0}}(\omega) \in W_{\ell}$ for some $\ell$, and by continuity even $X_{t}(\omega) \in W_{\ell}$ for all $t \in\left[t_{0}-\varepsilon, t_{0}+\varepsilon\right]$ with some sufficiently small $\varepsilon>0$. By definition of $t_{0}$ there exists $n_{0} \in \mathbb{N}, n_{0} \geq \ell$ such that $\tau_{n_{0}}^{0}(\omega)>t_{0}-\varepsilon$, with the consequence that then $\tau_{n_{0}}^{\ell}(\omega) \geq t_{0}+\epsilon$ which gives a contradiction.

Given a filtered probability space $\left(\Omega, \mathscr{F}, \mathbb{P} ;\left(\mathscr{F}_{t}\right)_{t \in \mathbb{R}_{+}}\right)$satisfying the usual conditions, we denote by $\mathscr{S}$ be the vector space of real-valued continuous semimartingales:

$$
\mathscr{S}=\mathscr{M} \oplus \mathscr{A}
$$

where $\mathscr{M}$ denotes the space of continuous local martingales and $\mathscr{A}$ the space of continuous adapted processes, starting at 0 almost surely, which are pathwise locally of bounded variation.

We start by stating an elementary but useful representation lemma.
Lemma 1.3.2. Let $M$ be an arbitrary differentiable manifold. There exist finitely many functions $h^{1}, \ldots, h^{\ell} \in C^{\infty}(M)$ such that the following properties hold:
(i) Each function $f \in C^{\infty}(M)$ factorizes through $\left(h^{1}, \ldots, h^{\ell}\right)$ as $f=\bar{f} \circ\left(h^{1}, \ldots, h^{\ell}\right)$ for some $\bar{f} \in C^{\infty}\left(\mathbb{R}^{\ell}\right)$.
(ii) Each section $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$ can be written as $b=\sum_{i, j=1}^{\ell} b_{i j} d h^{i} \otimes d h^{j}$ with functions $b_{i j} \in C^{\infty}(M)$.
(iii) Each differential form $\alpha \in \Gamma\left(T^{*} M\right)$ can be written as $\alpha=\sum_{i=1}^{\ell} \alpha_{i} d h^{i}$ with functions $\alpha_{i} \in C^{\infty}(M)$.
(iv) If $X$ is a semimartingale on $M$, then every continuous adapted $T^{*} M \otimes T^{*} M$-valued process $B$ above $X$ (i.e., $B_{t} \in T_{X_{t}}^{*} M \otimes T_{X_{t}}^{*} M$ fort $\in \mathbb{R}_{+}$) which is a semimartingale in the sense that $B_{t}(V, U)$ is a real semimartingale for any vector fields $V, U \in$ $\Gamma(T M)$, has a representation of the form $B=\sum_{i, j=1}^{\ell} B_{i j}\left(d h^{i} \otimes d h^{j}\right) \circ X$ with continuous adapted real-valued processes $B_{i j}$.
(v) If $X$ is a semimartingale on $M$, then every continuous adapted $T^{*} M$-valued process $J$ above $X$ (i.e., $J_{t} \in T_{X_{t}}^{*} M$ for $t \in \mathbb{R}_{+}$) which is a semimartingale in the sense that $J_{t}(V)$ is a real semimartingale for any vector fields $V \in \Gamma(T M)$, has a representation of the form $J=\sum_{i=1}^{\ell} J_{i}\left(d h^{i} \circ X\right)$ with continuous adapted real-valued processes $J_{i}$.

Proof. We represent $M$ via a Whitney embedding $h: M \longleftrightarrow \mathbb{R}^{\ell}$ as a closed submanifold of some $\mathbb{R}^{\ell}$. Then there exists a differentiable partition $\left(\varphi_{\lambda}\right)_{\lambda \in \Lambda}$ of the unity on $M$ and a family $\left(I_{\lambda}\right)_{\lambda \in \Lambda}$ of subsets $I_{\lambda} \subset\{1, \ldots, \ell\}$ with the following property: for each $\lambda \in \Lambda$ the family $\left(h^{i}\right)_{i \in I_{\lambda}}$ define a chart for $M$ on some open neighbourhood of $\operatorname{supp}\left(\varphi_{\lambda}\right)$.

Part (i) is evident: One defines $\bar{f} \mid h(M)$ through $f=\bar{f} \circ h$ and extends $\bar{f}$ constantly along the normal subspaces $\perp_{x} M$ to an open neighbourhood of $M \cong h(M)$, and finally smoothens $\bar{f}$ by multiplication with a function identical to 1 locally about $h(M)$ and vanishing outside a suitable larger tubular neighbourhood.

To part (ii): Note that $\varphi_{\lambda} b=\sum_{i, j=1}^{\ell} b_{i j}^{\lambda} d h^{i} \otimes d h^{j}$ with $b_{i j}^{\lambda} \in C^{\infty}(M)$ such that $\operatorname{supp}\left(b_{i j}^{\lambda}\right) \subset \operatorname{supp}\left(\varphi_{\lambda}\right)$ and $b_{i j}^{\lambda}:=0$ for $\{i, j\} \not \subset I_{\lambda}$, but then

$$
b=\sum_{i, j=1}^{\ell} b_{i j} d h^{i} \otimes d h^{j} \quad \text { where } b_{i j}:=\sum_{\lambda} b_{i j}^{\lambda}
$$

The proof of part (iii) is analogous to (ii).
To (iv): Analogously to (ii) we first write $\varphi_{\lambda}(X) B=\sum_{i, j=1}^{\ell} B_{i j}^{\lambda}\left(d h^{i} \otimes d h^{j}\right) \circ X$ with appropriate continuous $\mathbb{R}$-valued processes $B_{i j}^{\lambda}$, namely $B_{i j}^{\lambda}:=\varphi_{\lambda}(X) B\left(\frac{\partial}{\partial h^{i}}, \frac{\partial}{\partial h^{j}}\right)$ for $\{i, j\} \subset I_{\lambda}$ and $B_{i j}^{\lambda}:=0$ for $\{i, j\} \not \subset I_{\lambda}$. Summation over $\lambda$ then gives the claim.

The proof of (v) is again carried out analogously.
THEOREM 1.3.3. Let $X$ be an $M$-valued semimartingale. There exists a unique linear mapping $\Gamma\left(T^{*} M \otimes T^{*} M\right) \rightarrow \mathscr{A}, b \mapsto \int b(d X, d X)$, such that for all $f, g \in C^{\infty}(M)$,

$$
\begin{align*}
& d f \otimes d g \mapsto[f(X), g(X)]  \tag{1.3.1}\\
& f b \mapsto \int(f(X)) b(d X, d X) \tag{1.3.2}
\end{align*}
$$

Here, by definition, $b(d X, d X):=d \int b(d X, d X)$ and $[f(X), g(X)]$ in item (1.3.1) is the quadratic covariation process of $f(X)$ and $g(X)$.

DEFINITION 1.3.4 ( $b$-quadratic variation). The process $\int b(d X, d X)$ is called integral of $b$ along $X$ or $b$-quadratic variation of $X$. The random variable $\left(\int b(d X, d X)\right)_{t}$ giving its value at time $t$ is written as $\int_{0}^{t} b(d X, d X)$.

Proof (of Theorem 1.3.3). By Lemma 1.3.2 (ii) each section $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$ can be represented as $b=\sum b_{i j} d h^{i} \otimes d h^{j}$. We define

$$
\begin{equation*}
\int b(d X, d X):=\sum \int\left(b_{i j}(X)\right) d\left[h^{i}(X), h^{j}(X)\right] \tag{1.3.3}
\end{equation*}
$$

Then uniqueness is obvious; to prove existence it remains to show that (1.3.3) is welldefined. To this end assume that

$$
b=\sum_{\text {finite }} u_{\nu} d f^{\nu} \otimes d g^{\nu}=0
$$

We need to check that

$$
\sum_{\nu} u_{\nu}(X) d\left[f^{\nu}(X), g^{\nu}(X)\right]=0
$$

as well. Without loss of generality, by means of Lemma 1.3.1, we may assume that $h$ is already a global chart for $M$. According to Lemma 1.3.2 (i) we write $u_{\nu}=\bar{u}_{\nu} \circ h$, $f^{\nu}=\bar{f}^{\nu} \circ h$ and $g^{\nu}=\bar{g}^{\nu} \circ h$ in terms of appropriate extensions $\bar{u}_{\nu}, \bar{f}^{\nu}, \bar{g}^{\nu} \in C^{\infty}\left(\mathbb{R}^{\ell}\right)$. Defining $\bar{X}=h \circ X$, the claim then follows from the following calculation:

$$
\begin{aligned}
\sum_{\nu} & u_{\nu}(X) d\left[f^{\nu}(X), g^{\nu}(X)\right]=\sum_{\nu} \bar{u}_{\nu}(\bar{X}) d\left[\bar{f}^{\nu}(\bar{X}), \bar{g}^{\nu}(\bar{X})\right] \\
& =\sum_{i, j} \sum_{\nu} \bar{u}_{\nu}(\bar{X})\left(D_{i} \bar{f}^{\nu}\right)(\bar{X})\left(D_{j} \bar{g}^{\nu}\right)(\bar{X}) d\left[\bar{X}^{i}, \bar{X}^{j}\right] \\
& =\sum_{i, j}\left(\sum_{\nu} u_{\nu} d f^{\nu} \otimes d g^{\nu}\right)\left(\left(\frac{\partial}{\partial h^{i}}\right)_{X},\left(\frac{\partial}{\partial h^{j}}\right)_{X}\right) d\left[\bar{X}^{i}, \bar{X}^{j}\right]=0 .
\end{aligned}
$$

COROLLARY 1.3.5. The b-quadratic variation $\int b(d X, d X)$ depends only on the symmetric part of $b$. In particular, $\int b(d X, d X)=0$ if $b$ is antisymmetric.

Proof. Defining $\bar{b}(v, w):=b(w, v)$, the assignment $b \mapsto \int \bar{b}(d X, d X)$ has the defining properties (1.3.1) and (1.3.2) as well.

The next remark is again an immediate consequence of the defining properties (1.3.1) and (1.3.2) of the $b$-quadratic variation.

REMARK 1.3.6. The $b$-quadratic variation of a semimartingale commutes with timechange. More precisely, the following holds: Let $X$ be an $M$-valued semimartingale, let $\left(\tau_{t}\right)_{t \geq 0}$ be a continuous finite time-change, and consider the time-changed semimartingale $\hat{X}$ defined by $\hat{X}_{t}=X_{\tau_{t}}$ (w.r.t. the time-changed filtration $\left(\hat{\mathscr{F}}_{t}\right)_{t \geq 0}:=\left(\mathscr{F}_{\tau_{t}}\right)_{t \geq 0}$ ). Then

$$
\int_{0}^{t} b(d \hat{X}, d \hat{X})=\int_{\tau_{0}}^{\tau_{t}} b(d X, d X)
$$

In particular, for an arbitrary stopping time $\tau$, if we denote by $X_{t}^{\tau}=X_{t \wedge \tau}$ the semimartingale stopped at the random time $\tau$, then the formula $\int b\left(d X^{\tau}, d X^{\tau}\right)=\left(\int b(d X, d X)\right)^{\tau}$ where on the right-hand side the process $\int b\left(d X^{\tau}, d X^{\tau}\right)$ is stopped at time $\tau$.

REMARK 1.3.7. (i) (Induced form) Let $\phi: M \rightarrow N$ be a differentiable map between manifolds, $E$ be a vector bundle over $N$ and $s \in \mathbb{N} \cup\{0\}$. Each multilinear form $L \in$ $\Gamma\left(T^{*} N^{\otimes s} \otimes E\right)$ taking values in $E$ induces via

$$
\left(\phi^{*} L\right)_{p}\left(w_{1}, \ldots, w_{s}\right):=L_{\phi(p)}\left(d \phi_{p} w_{1}, \ldots, d \phi_{p} w_{s}\right), \quad w_{i} \in T_{p} M, p \in M
$$

a multilinear form $\phi^{*} L \in \Gamma\left(T^{*} M^{\otimes s} \otimes \phi^{*} E\right)$ with values in $\phi^{*} E$, called pullback of $L$ via $\phi$. In particular, to each $X \in \Gamma(E)$ there is the induced section $\phi^{*} X \in \Gamma\left(\phi^{*} E\right)$ with $\left(\phi^{*} X\right)_{p}=X_{\phi(p)}, p \in M$.
(ii) (Induced frame) Let $e_{1}, \ldots, e_{m} \in \Gamma(E / U)$ be a local frame for $E$. Then

$$
\phi^{*} e_{1}, \ldots, \phi^{*} e_{m} \in \Gamma\left(\phi^{*} E / \phi^{-1}(U)\right)
$$

is a local frame for $\phi^{*} E$. Hence, to each section $Y \in \Gamma\left(\phi^{*} E\right)$, there exist uniquely determined functions $b^{1}, \ldots, b^{m} \in C^{\infty}\left(\phi^{-1}(U)\right)$ such that $Y \mid \phi^{-1}(U)=\sum b^{i} \phi^{*} e_{i}$.

THEOREM 1.3.8 (Pullback formula for the $b$-quadratic variation). Let $\phi: M \rightarrow N$ be a differentiable map and $b \in \Gamma\left(T^{*} N \otimes T^{*} N\right)$. Then, for any semimartingale $X$ on $M$,

$$
\begin{equation*}
\int\left(\phi^{*} b\right)(d X, d X)=\int b(d(\phi \circ X), d(\phi \circ X)) \tag{1.3.4}
\end{equation*}
$$

Proof. The left-hand side of (1.3.4) satisfies the defining properties for the $b$-quadratic variation of $\phi(X)$.

We now turn to the problem of integrating one-forms on $M$ along $M$-valued semimartingales, see [21].

THEOREM 1.3.9. Let $X$ be a semimartingale taking values in $M$. There is a unique linear mapping

$$
\Gamma\left(T^{*} M\right) \equiv A^{1}(M) \rightarrow \mathscr{S}, \quad \alpha \mapsto \int \alpha(\circ d X) \equiv \int_{X} \alpha
$$

such that for all $f \in C^{\infty}(M)$,

$$
\begin{align*}
d f & \mapsto f(X)-f\left(X_{0}\right)  \tag{1.3.5}\\
f \alpha & \mapsto \int f(X) \circ \alpha(\circ d X) \tag{1.3.6}
\end{align*}
$$

On the right-hand side of (1.3.6) we have the Stratonovich integral of the process $f(X)$ with respect to the semimartingale $\int \alpha(\circ d X)$, thus

$$
f(X) \circ \alpha(\circ d X) \equiv f(X) \circ d\left(\int \alpha(\circ d X)\right)
$$

DEFINITION 1.3.10 (Stratonovich integral of one-forms along semimartingales). The process $\int \alpha(\circ d X)$ is called the Stratonovich integral of $\alpha$ along $X$. We also use the notation $\int_{X} \alpha$ for $\int \alpha(\circ d X)$.

Proof (of Theorem 1.3.9). By Lemma 1.3 .2 (iii) differential forms $\alpha \in \Gamma\left(T^{*} M\right)$ can be represented as $\alpha=\sum_{i} \alpha_{i} d h^{i}$ with functions $\alpha_{i} \in C^{\infty}(M)$. We define

$$
\begin{equation*}
\int_{X} \alpha:=\sum_{i} \int \alpha_{i}(X) \circ d\left(h^{i}(X)\right) \tag{1.3.7}
\end{equation*}
$$

Uniqueness is again obvious; it is thus sufficient to show that formula (1.3.7) is welldefined. To this end, we have to verify that if $\alpha=\sum_{\text {finite }} u_{\nu} d f^{\nu}=0$ then

$$
\sum_{\nu} u_{\nu}(X) \circ d\left(f^{\nu}(X)\right)=0
$$

holds as well. Proceeding as in the proof of Theorem 1.3.3, without loss of generality, we assume again that $h$ is already a global chart for $M$. But then we have

$$
\sum_{\nu} u_{\nu}(X) \circ d\left(f^{\nu}(X)\right)=\sum_{\nu} \bar{u}_{\nu}(\bar{X}) \circ d\left(\bar{f}^{\nu}(\bar{X})\right)
$$

$$
\begin{aligned}
& =\sum_{i} \sum_{\nu} \bar{u}_{\nu}(\bar{X}) \circ\left(D_{i} \bar{f}^{\nu}(\bar{X}) \circ d \bar{X}^{i}\right) \\
& =\sum_{i}\left(\left(\sum_{\nu} u_{\nu} d f^{\nu}\right)\left(\frac{\partial}{\partial h^{i}}\right)_{X}\right) \circ d \bar{X}^{i}=0
\end{aligned}
$$

which gives the claim.
EXAMPLE 1.3.11. In the special case of a deterministic $C^{1}$ curve $X$ in $M$, say $X_{t}=$ $x(t)$, which is trivially a semimartingale, we obtain

$$
\begin{equation*}
\int_{X} \alpha=\int \alpha(\dot{x}(t)) d t, \quad \alpha \in \Gamma\left(T^{*} M\right) \tag{1.3.8}
\end{equation*}
$$

Indeed, the right-hand side of (1.3.8) obviously has the defining properties of $\int_{X} \alpha$.
REMARK 1.3.12. Stratonovich integration of differential forms $\alpha$ along semimartingales commutes with time-change. More precisely, the following holds: Let $X$ be a semimartingale taking values in $M,\left(\tau_{t}\right)_{t \geq 0}$ a continuous finite time-change, and consider the time-changed semimartingale $\hat{X}$ defined by $\hat{X}_{t}:=X_{\tau_{t}}$ (with respect to the time-changed filtration $\left.\left(\hat{\mathscr{F}}_{t}\right)_{t \geq 0}:=\left(\mathscr{F}_{\tau_{t}}\right)_{t \geq 0}\right)$. Then

$$
\int_{0}^{t} \alpha(\circ d \hat{X})=\int_{\tau_{0}}^{\tau_{t}} \alpha(\circ d X)
$$

In particular, for an arbitrary stopping time $\tau$, if we denote by $X_{t}^{\tau}=X_{t \wedge \tau}$ the semimartingale stopped at the random time $\tau$, then the formula

$$
\int_{X^{\tau}} \alpha=\left(\int_{X} \alpha\right)^{\tau}
$$

holds where on the right-hand side the semimartingale $\int_{X} \alpha$ is stopped at time $\tau$.
THEOREM 1.3.13 (Pullback formula for the Stratonovich integral of a one-form). Let $\phi: M \rightarrow N$ be a differentiable map and $\alpha \in A^{1}(N) \equiv \Gamma\left(T^{*} N\right)$. Then, for any semimartingale $X$ on $M$,

$$
\begin{equation*}
\int_{X} \phi^{*} \alpha=\int_{\phi \circ X} \alpha \tag{1.3.9}
\end{equation*}
$$

Proof. The left-hand side of Eq. (1.3.9) satisfies the defining properties for the Stratonovich integral of $\alpha$ along $\phi \circ X$.

Remark 1.3.14. Let $\alpha, \beta \in \Gamma\left(T^{*} M\right)$. Then $\alpha \otimes \beta \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$ and for the quadratic covariation process of $\int_{X} \alpha$ and $\int_{X} \beta$ we have the formula:

$$
\begin{equation*}
\left[\int_{X} \alpha, \int_{X} \beta\right]=\int(\alpha \otimes \beta)(d X, d X) \tag{1.3.10}
\end{equation*}
$$

We continue with the observation that Theorems 1.3 .3 and 1.3.9 can be slightly extended in an obvious way. In Eqs. (1.3.2) and (1.3.6), instead of $f(X)$ where $f \in C^{\infty}(M)$, more generally, continuous adapted $\mathbb{R}$-valued processes $K$ may serve as multipliers.

THEOREM 1.3.15. Let $X$ be an $M$-valued semimartingale and let $\mathbb{B}$ be the real vector space of continuous adapted $T^{*} M \otimes T^{*} M$-valued processes $B$ over $X$ such that $B_{t}(V, U)$ are real semimartingales for any vector fields $V, U \in \Gamma(T M)$. There exists exactly one linear mapping

$$
\mathbb{B} \rightarrow \mathscr{A}, \quad B \mapsto \int B(d X, d X)
$$

with the following properties:

$$
\begin{aligned}
b \circ X & \mapsto \int b(d X, d X) \quad \text { for any } b \in \Gamma\left(T^{*} M \otimes T^{*} M\right) \\
K B & \mapsto \int K B(d X, d X) \text { for any continuous adapted real-valued processes } K .
\end{aligned}
$$

Here $\int K B(d X, d X):=\int K d\left(\int B(d X, d X)\right)$.
Proof. According to Lemma 1.3.2 (iv) each continuous adapted $T^{*} M \otimes T^{*} M$-valued process $B$ over $X$ has a representation as a finite sum of the form

$$
B=\sum_{\nu} B_{\nu}\left(d f^{\nu} \otimes d g^{\nu}\right) \circ X
$$

We set

$$
\int B(d X, d X):=\sum_{\nu} B_{\nu} d\left[f^{\nu}(X), g^{\nu}(X)\right]
$$

Well-definedness is verified as in the proof of Theorem 1.3.3.
THEOREM 1.3.16. Let $X$ be an $M$-valued semimartingale and let $\mathbb{D}$ be the real vector space of continuous adapted $T^{*} M$-valued processes $J$ over $X$ such that $J_{t}(V)$ are real semimartingales for any vector field $V \in \Gamma(T M)$. There exists exactly one linear mapping $\mathbb{D} \rightarrow \mathscr{S}, \quad J \mapsto \int J(\circ d X) \equiv \int_{X} J$, with the following properties:
$\alpha \circ X \mapsto \int \alpha(\circ d X)=\int_{X} \alpha \quad$ for any $\alpha \in \Gamma\left(T^{*} M\right)$,
$K J \mapsto \int K \circ J(\circ d X) \quad$ for any continuous adapted $\mathbb{R}$-valued process $K$.
Here $\int K \circ J(\circ d X):=\int K \circ d\left(\int J(\circ d X)\right)$.
Proof. According to Lemma 1.3.2 (v) each continuous adapted $T^{*} M$-valued process $J$ over $X$ has a representation as a finite sum of the form

$$
J=\sum_{\nu} J_{\nu}\left(d f^{\nu} \circ X\right)
$$

We set

$$
\int J(\circ d X):=\sum_{\nu} J_{\nu} \circ d\left(f^{\nu}(X)\right)
$$

Well-definedness is verified with the same calculation as in the proof of Theorem 1.3.9.
The pullback formulas (1.3.4) and (1.3.9) carry over in an obvious way.
REMARK 1.3.17 (Pullback formulas). Let $\phi: M \rightarrow N$ be a differentiable map and $X$ be a semimartingale on $M$.
(i) For a continuous adapted $T^{*} N \otimes T^{*} N$-valued process $B$ over $\phi \circ X$ we have:

$$
\int\left(\phi^{*} B\right)(d X, d X)=\int B(d(\phi(X)), d(\phi(X)))
$$

(ii) For a continuous adapted $T^{*} N$-valued process $J$ over $\phi \circ X$ we have:

$$
\int_{X} \phi^{*} J=\int_{\phi(X)} J
$$

REMARK 1.3.18. Under a complex differential form $\alpha$ on a differentiable manifold $M$ we understand a section $\alpha \in \Gamma\left(T^{*} M \otimes \mathbb{C}\right)$. Decomposing $\alpha$ into its real and imaginary part, i.e., $\alpha=\alpha_{1}+i \alpha_{2}$ where $\alpha_{i} \in \Gamma\left(T^{*} M\right)$ are real differential forms on $M$, we extend the Stratonovich integral of differential forms along $M$-valued semimartingales via

$$
\int_{X} \alpha:=\int_{X} \alpha_{1}+i \int_{X} \alpha_{2}
$$

to complex differential forms.
As an example for Stratonovich integration of one-forms we consider the winding of semimartingales in the plane. This notion generalizes the classical winding number of a (closed) differentiable curve in $\mathbb{C} \backslash\{0\}$, as defined in elementary function theory, to semimartingales in the plane. We identify the complex plane $\mathbb{C}$ with the Euclidean space $\mathbb{R}^{2}$.

REMARK 1.3.19 (Winding of a semimartingale in the plane). Let $Z$ be a continuous $\mathbb{C}$-valued semimartingale such that $Z_{0} \neq 0$ and $Z$ does not hit the origin almost surely. Integration of the complex differential form $\alpha=d z / z$ on $\mathbb{C} \backslash\{0\}$ along $Z$,

$$
\int_{Z} \alpha=\int \frac{1}{Z} \circ d Z \in \mathscr{S}+i \mathscr{S}
$$

gives a continuous version of a logarithm along the paths of $Z$ via

$$
\log _{\omega}\left(Z_{t}(\omega)\right)-\log _{\omega}\left(Z_{0}(\omega)\right):=\left(\int_{Z} \frac{d z}{z}\right)_{t}(\omega), \quad t \geq 0, \quad \mathbb{P} \text {-almost all } \omega \in \Omega
$$

In other words, writing

$$
Z_{t} \equiv\left|Z_{t}\right| e^{i \Theta_{t}}, \quad t \geq 0
$$

with a (pathwise) continuous version $\Theta_{t}$ of the argument of $Z_{t}$, then

$$
\Theta_{t}=\Theta_{0}+\operatorname{Im}\left(\int_{Z} \frac{d z}{z}\right)_{t}
$$

The process $\operatorname{Im} \int_{Z} \frac{d z}{z}$ is called winding of the semimartingale $Z$ about the origin.
Proof. It is sufficient to verify that, modulo indistinguishability,

$$
\exp \left(\int_{0} \frac{1}{Z} \circ d Z\right)=\frac{Z}{Z_{0}}
$$

But using the abbreviation $L:=\int_{Z} d z / z \equiv \int Z^{-1} \circ d Z$, then

$$
d e^{L}=e^{L} \circ d L=\left(e^{L} / Z\right) \circ d Z
$$

and hence

$$
d\left(\frac{e^{L}}{Z}\right)=e^{L}\left(-\frac{1}{Z^{2}}\right) \circ d Z+\frac{1}{Z}\left(\frac{e^{L}}{Z}\right) \circ d Z=0
$$

In the sequel let $\mathscr{M}(\mathbb{C})$ denote the class of $\mathbb{C}$-valued local martingales. A local martingale $Z=X+i Y \in \mathscr{M}(\mathbb{C})$ is said to be conformal if $[X, X]=[Y, Y]$ and $[X, Y]=0$, or equivalently, if $d Z d Z=0$.

REMARK 1.3.20. Stratonovich integrals of holomorphic differential forms along conformal martingales give local martingales. More precisely: Let $Z$ be a conformal local martingale and $D \subset \mathbb{C}$ be a domain not left by $Z$ a.s. For any complex differential form $\alpha=f(z) d z$ on $D$ (where $f: D \rightarrow \mathbb{C}$ is a holomorphic function) the process

$$
\int_{Z} \alpha \equiv \int f(Z) \circ d Z=\int f(Z) d Z \in \mathscr{M}(\mathbb{C})
$$

is a conformal local martingale. On the other hand, a local martingale $Z$ in $\mathbb{C}$ is already a conformal local martingale if $\int_{Z} \alpha \in \mathscr{M}(\mathbb{C})$ for $\alpha=z d z$.

Proof. Indeed we have

$$
\begin{equation*}
f(Z) \circ d Z=f(Z) d Z+\frac{1}{2}\left(f^{\prime}(Z)\right) d Z d Z=f(Z) d Z \tag{1.3.11}
\end{equation*}
$$

where the first equality in (1.3.11) results from the Itô formula for complex semimartingales (e.g. [16] Corollary to Theorem $4.46^{\prime}$ ), whereas the second equality is a consequence of the conformity of $Z$. In addition local martingales of the type $N=\int f(Z) d Z$ are automatically conformal, since $d N d N=f(Z)^{2} d Z d Z=0$. The last statement follows with $f=\mathrm{id}$.

In particular, if in the situation of Remark 1.3.20 the conformal local martingale $Z$ is a Brownian motion on $\mathbb{C}$, then for each holomorphic function $f$ the process $\int_{Z} f(z) d z$ is a conformal local martingale, and thus there exist independent one-dimensional Brownian motions $B$ and $\beta$ such that

$$
\operatorname{Re} \int_{Z} f(z) d z=B_{T_{t}}, \quad \operatorname{Im} \int_{Z} f(z) d z=\beta_{T_{t}}
$$

where the time-change is given by $T_{t}:=\int_{0}^{t}\left(|f|^{2} \circ Z_{s}\right) d s, t \geq 0$. If $f \not \equiv 0$ then $T_{\infty} \equiv \infty$ $\mathbb{P}$-a.s., as is easily verified by using recurrence and the strong Markov property of the 2-dimensional Brownian motion:

Lemma 1.3.21. Let $Z$ be a Brownian motion on $\mathbb{R}^{2}$ and $f: \mathbb{R}^{2} \rightarrow \mathbb{R}$ be a continuous, not identically vanishing function. Then $\int_{0}^{\infty}|f|^{2}\left(Z_{s}\right) d s \equiv \infty, \mathbb{P}$-a.s.

We are going to summarize the results above in the case $f(z) d z=d z / z$.
Corollary 1.3.22. Let $Z=X+i Y$ be a BM in $\mathbb{C}$ starting from some point $z_{0} \neq 0$. Then $Z_{t}=\left|Z_{t}\right| e^{i \Theta_{t}}$ where

$$
\begin{aligned}
\log \left|Z_{t}\right|-\log \left|Z_{0}\right| & =\operatorname{Re} \int_{0}^{t} \frac{d Z}{Z}
\end{aligned}=\int_{0}^{t} \frac{X d X+Y d Y}{|Z|^{2}}, ~=\Theta_{0}=\operatorname{Im} \int_{0}^{t} \frac{d Z}{Z}=\int_{0}^{t} \frac{X d Y-Y d X}{|Z|^{2}} .
$$

In addition there exist independent one-dimensional Brownian motions $B$ and $\beta$ such that

$$
\int_{0}^{t} \frac{d Z}{Z}=B_{T_{t}}+i \beta_{T_{t}}
$$

with the time-change $T_{t}$ given by $T_{t}:=\int_{0}^{t}\left|Z_{s}\right|^{-2} d s$.
Since $T_{\infty}=\infty \mathbb{P}$-a.s., one concludes from $\Theta_{t}-\Theta_{0}=\beta_{T_{t}}$ that $\operatorname{BM}(\mathbb{C})$ winds with probability 1 arbitrary often clockwise and anti-clockwise about any given point, but unwinds again almost surely. On the other hand, $|Z|$ and $B$ generate the same $\sigma$-algebra, hence

$$
\mathscr{B}_{\infty}:=\sigma\left\{\left|Z_{s}\right|: s \in \mathbb{R}_{+}\right\}=\sigma\left\{B_{s}: s \in \mathbb{R}_{+}\right\} \quad \text { modulo } \mathbb{P} \text {-nullsets; }
$$

indeed first of all $\sigma\left\{\left|Z_{s}\right|: s \leq t\right\}=\sigma\left\{\log \left|Z_{s}\right|: s \leq t\right\}=\sigma\left\{B_{T_{s}}: s \leq t\right\}$; on the other hand, the time-change $\left(T_{t}\right)_{t \geq 0}$ may be described in terms of $B$, as is seen from the formula

$$
\begin{equation*}
T_{t}=\inf \left\{s \geq 0:\left|z_{0}\right|^{2} \int_{0}^{s} \exp \left(2 B_{r}\right) d r>t\right\} \tag{1.3.12}
\end{equation*}
$$

which is easily verified with the substitution $r=T_{u}$. As a consequence, the $\mathrm{BM} \beta$ describing the angular process is independent of the whole radial process, and hence independent
of $\mathscr{B}_{\infty} \equiv \sigma\left\{\left|Z_{s}\right|: s \in \mathbb{R}_{+}\right\}$and in particular of the time-change $\left(T_{t}\right)_{t \geq 0}$. Thus for any $\xi \in \mathbb{R}$ :

$$
\mathbb{E}^{\mathscr{B}_{\infty}}\left[\exp \left(i \xi\left(\Theta_{t}-\Theta_{0}\right)\right)\right]=\exp \left(-\xi^{2} / 2 T_{t}\right) \quad \mathbb{P} \text {-a.s. }
$$

This formula allows to calculate the distribution of $\Theta_{t}$ for fixed $t$, and is moreover a useful tool for many explicit calculations related to the stochastic behaviour of BM in the plane (e.g. [48], [49]).

### 1.4. Linear Connections and Martingales on Manifolds

The aim of this section is to introduce martingales on manifolds. This task requires on the manifold a linear connection as additional geometric structure. We start the discussion by recalling basic notions from differential geometry; for more background on these topics the reader may consult $[\mathbf{1 3}, \mathbf{2 5}, \mathbf{2 6}, 27]$.

From a geometrical point of view we want to deal with the following situation. Let $\pi: E \rightarrow M$ be a vector bundle over a manifold $M$, for instance the tangent bundle $T M$ of $M$, and let $\alpha:[0,1] \rightarrow M$ be a differentiable curve such that $\alpha(0)=p$ and $\alpha(1)=q$. We look for a canonical procedure to translate vectors $v \in E_{p}$ to $E_{q}$ along the curve $\alpha$.

If in addition $E$ is endowed with a metric, in the sense that each fiber $E_{x}$ carries a scalar product depending smoothly on $x$, then it is natural to demand in addition that angles are preserved by the translation along curves.


Figure 1.4.1. Parallel transport
The fibers of a vector bundle are all isomorphic to a fixed finite-dimensional vector space which however does not mean that there is a canonical way to identify them. The additional structure needed to relate fibers among each other in an intrinsic way is a "linear connection" in $E$. Such a structure encodes the information necessary to transport elements of one fiber of $E$ along some curve to another fiber.

There are different (but equivalent) ways to introduce linear connections in a vector bundle $E$, for instance, as parallel transport, as covariant derivative, or horizontal splitting of $T E$. The most intuitive way is the concept of a parallel transport.

DEFINITION 1.4.1 (Parallel transport). Let $\pi: E \rightarrow M$ be a vector bundle over a differentiable manifold $M$. A parallel transport $L$ in $E$ is an assignment of a linear isomorphism $L_{\alpha}: E_{p} \rightarrow E_{q}$ to each differentiable path $\alpha$ from $p$ to $q$ in $M$ such that the following properties hold:
(i) (Invariance under reparametrization) If $\alpha:[a, b] \rightarrow M$ is a differentiable curve then $L_{\alpha \circ \varphi}=L_{\alpha}$ for any differentiable reparametrization $\varphi:\left[a^{\prime}, b^{\prime}\right] \rightarrow[a, b]$ such that $\varphi\left(a^{\prime}\right)=a$ and $\varphi\left(b^{\prime}\right)=b$.
(ii) (Transitivity) If $\alpha:[a, b] \rightarrow M$ and $a \leq c \leq b$ then $L_{\alpha}=L_{\alpha \mid[c, b]} \circ L_{\alpha \mid[a, c]}$.
(iii) (Behaviour under back-transport) $L_{\alpha^{-}}=L_{\alpha}^{-1}$ for $\alpha^{-}:[a, b] \rightarrow M, t \mapsto \alpha(a+b-t)$.
(iv) (Dependence on parameters) If $\alpha$ depends differentiably on parameters (e.g. if $\alpha$ is a differentiable family of curves), then $L_{\alpha}$ depends differentiably on these parameters as well.
(v) (First-Order-Axiom) For any $X \in \Gamma(E)$ and $v \in T_{p} M$ the covariant derivative $\nabla_{v} X$ of $X$ in direction $v$,

$$
\begin{array}{ll}
\nabla_{v} X:=\nabla_{D}(X \circ \alpha)(0) \in E_{p} & \text { for } \alpha:[-\varepsilon, \varepsilon] \rightarrow M C^{\infty} \text {-curve } \\
& \text { with } \alpha(0)=p \text { and } \dot{\alpha}(0)=v
\end{array}
$$

is well-defined and independent of the choice of the curve $\alpha$.
In (v) we use the following notion: for a differentiable curve $\alpha:[a, b] \rightarrow M$ and a $C^{\infty}$ section $\sigma \in \Gamma\left(\alpha^{*} E\right)$, the covariant derivative $\nabla_{D} \sigma \in \Gamma\left(\alpha^{*} E\right)$ of $\sigma$ along $\alpha$ with respect to $L$ is defined as

$$
\left(\nabla_{D} \sigma\right)(t):=\left.\frac{d}{d \varepsilon}\right|_{\varepsilon=0} L_{\alpha \mid[t, t+\varepsilon]}^{-1} \sigma(t+\varepsilon) \in E_{\alpha(t)}
$$

if well-defined.
Before introducing the abstract notion of a covariant derivative on a vector bundle, we state the following Lemma.

LEMMA 1.4.2. Let $E, F$ be vector bundles over $M$, further $K: \Gamma(E) \rightarrow \Gamma(F)$ a $C^{\infty}(M)$-linear map and $p \in M$. Then $K(A)_{p}=K(B)_{p}$ for all sections $A, B \in \Gamma(E)$ with $A_{p}=B_{p}$. Thus $K$ provides a section of the bundle $\operatorname{Hom}(E, F) \cong E^{*} \otimes F$.

Proof. It is sufficient to show that $A_{p}=0$ already implies $K(A)_{p}=0$. Fixing a local frame $e_{1}, \ldots, e_{m} \in \Gamma(E / U)$ at $p$ there exist uniquely determined functions $a^{1}, \ldots, a^{m} \in C^{\infty}(U)$ such that $A \mid U=\sum_{i} a^{i} e_{i}$. In particular, we have $a^{1}(p)=\cdots=$ $a^{m}(p)=0$. Now let $\psi \in C^{\infty}(M)$ such that $\psi(p)=1$ and $\operatorname{supp} \psi \subset U$. In particular, $\bar{e}_{i}:=\psi e_{i} \in \Gamma(E / U)$ and $\bar{a}^{i}:=\psi a^{i} \in C^{\infty}(U)$ extend smoothly to global sections, resp. functions on $M$ (being equal to 0 outside of $U$ ). Then $\psi^{2} A=\sum_{i} \bar{a}^{i} \bar{e}_{i}$, and thus $K(A)_{p}=\psi(p)^{2} K(A)_{p}=K\left(\psi^{2} A\right)_{p}=\sum_{i} \bar{a}^{i}(p) K\left(\bar{e}_{i}\right)_{p}=0$.

Definition 1.4.3 (Covariant derivative). Let $E$ be a vector bundle over a manifold $M$. A covariant derivative on $E$ is an $\mathbb{R}$-linear mapping

$$
\nabla: \Gamma(E) \rightarrow \Gamma\left(T^{*} M \otimes E\right)
$$

satisfying the product rule

$$
\nabla(f X)=d f \otimes X+f \nabla X, \quad X \in \Gamma(E), f \in C^{\infty}(M)
$$

Sections $X \in \Gamma(E)$ with the property that $\nabla X=0$ are called parallel.
REMARK 1.4.4. Since according to Lemma 1.4.2,

$$
\Gamma\left(T^{*} M \otimes E\right) \cong \operatorname{Hom}_{C^{\infty}(M)}(\Gamma(T M), \Gamma(E))
$$

a covariant derivative $\nabla$ on $E$ can equally be seen as $\mathbb{R}$-bilinear mapping

$$
\Gamma(T M) \times \Gamma(E) \rightarrow \Gamma(E), \quad(A, X) \mapsto \nabla_{A} X:=(\nabla X) A
$$

In this notation a covariant derivative is $C^{\infty}(M)$-linear in the first argument and (as a consequence of the product rule) derivative in the second argument, i.e.,
(i) $\nabla_{f A} X=f \nabla_{A} X$;
(ii) $\nabla_{A}(f X)=(A f) X+f \nabla_{A} X$,
for all $A \in \Gamma(T M), X \in \Gamma(E)$ and $f \in C^{\infty}(M)$.
REmARK 1.4.5. (i) Given $A \in \Gamma(T M), X \in \Gamma(E)$ and $p \in M$, by Lemma 1.4.2, $\left(\nabla_{A} X\right)_{p}$ depends only on $A_{p} \in T_{p} M$. Hence, for $v \in T_{p} M$ choosing $A \in \Gamma(T M)$ such that $A_{p}=v$, then $\nabla_{v} X:=\left(\nabla_{A} X\right)_{p} \in E_{p}$ is well-defined and is called covariant derivative of $X$ in direction $v$.
(ii) For $v \in T_{p} M, X \in \Gamma(E)$ the covariant derivative $\nabla_{v} X$ depends only on the germ of $X$ at $p$.

Proof. For $p \in U \subset M$ open and $X \mid U \equiv 0$ we have to show that $\nabla_{v} X=0$. To this end let $\psi \in C^{\infty}(M)$ such that $\operatorname{supp} \psi \subset U$ and $\psi(p)=1$. But then $\psi X \equiv 0$, and thus $0=\nabla_{v}(\psi X)=v(\psi) X_{p}+\psi(p) \nabla_{v} X=\nabla_{v} X$.

Notation 1.4.6 (Christoffel symbols). Let $E$ be a vector bundle of rank $m$ over an $n$ dimensional manifold $M$. Let $\nabla$ be a covariant derivative on $E$ and $e_{1}, \ldots, e_{m} \in \Gamma(E / U)$ a local frame for $E$. If in addition $(h, U)$ is a local chart for $M$ and $\partial_{i}=\frac{\partial}{\partial h^{i}} \in \Gamma(T M / U)$ the corresponding local coordinate vector fields, then $\left(\partial_{1}, \ldots, \partial_{n}\right)$ defines a local frame for $T M$, and the sections $\nabla_{\partial_{i}} e_{j} \in \Gamma(E / U)$ are well-defined by Remark 1.4.5. The uniquely determined functions $\Gamma_{i j}^{k} \in C^{\infty}(U)$ such that

$$
\nabla_{\partial_{i}} e_{j}=\sum_{k=1}^{m} \Gamma_{i j}^{k} e_{k}
$$

are called the Christoffel symbols of $\nabla$ with respect to $(h, U)$ and $e_{1}, \ldots, e_{m} \in \Gamma(E / U)$. They determine the covariant derivative $\nabla$ on $E / U$.

A covariant derivative on a vector bundle $E$ in the sense of Definition 1.4.3 induces canonically a notion of covariant derivative of sections along maps. For a precise statement we come back to the notion of induced forms and frames as introduced in Remark 1.3.7.

DEFINITION 1.4.7 (Induced covariant derivative). Let $f: M \rightarrow N$ be a differentiable map between manifolds and $\nabla$ a covariant derivative on a vector bundle $E$ over $N$. There exists exactly one covariant derivative on the induced bundle $f^{*} E$ over $M$ (called the covariant derivative on $f^{*} E$ induced by $f$ and denoted again by $\nabla$ ) such that

$$
\begin{equation*}
\nabla_{w}\left(f^{*} X\right)=\nabla_{d f_{p} w} X \in E_{f(p)}, \quad X \in \Gamma(E), w \in T_{p} M, p \in M \tag{1.4.1}
\end{equation*}
$$

Indeed, let $e_{1}, \ldots, e_{m} \in \Gamma(E / U)$ be a local frame for $E$ and $Y \in \Gamma\left(f^{*} E\right)$ a global section. By Remark 1.3.7 (ii), $Y$ has a unique representation on $f^{-1}(U)$ of the form $Y \mid f^{-1}(U)=\sum_{i} b^{i} f^{*} e_{i}$ where $b^{i} \in C^{\infty}\left(f^{-1}(U)\right)$. For $w \in T_{p} M, p \in f^{-1}(U)$ we deduce from the product rule and property (1.4.1) of $\nabla$ that

$$
\begin{equation*}
\nabla_{w} Y=\sum_{i=1}^{m}\left(w\left(b^{i}\right)\left(e_{i}\right)_{f(p)}+b^{i}(p) \nabla_{d f_{p} w} e_{i}\right) \in E_{f(p)} \tag{1.4.2}
\end{equation*}
$$

This shows uniqueness of the induced covariant derivative. On the other hand, Eq. (1.4.2) defines a covariant derivative on $f^{*} E$ which establishes existence.

If $X$ is a section of $E$ along a curve $\alpha$ on $M$, then Definition 1.4.7 gives in particular a notion of a covariant derivative $X$ along $\alpha$.

DEfinition 1.4.8. Let $\nabla$ be a covariant derivative on a vector bundle $E$ over $M$ and $\alpha: I \rightarrow M$ a differentiable curve defined on some real interval.
(i) (Covariant derivative for sections along curves) For sections $X \in \Gamma\left(\alpha^{*} E\right)$ along $\alpha$ the vector field $\nabla_{D} X \in \Gamma\left(\alpha^{*} E\right)$ is called the covariant derivative of $X$ along $\alpha$; here $D$ denotes the canonical vector field on $I$.
(ii) (Parallel sections along curves) A section $X \in \Gamma\left(\alpha^{*} E\right)$ along $\alpha$ is said to be parallel along $\alpha$ (with respect to $\nabla$ ) if $\nabla_{D} X=0$. The linear subspace of $\Gamma\left(\alpha^{*} E\right)$ of parallel sections along $\alpha$ is denoted $\Gamma_{\text {par }}\left(\alpha^{*} E\right)$.
DEFINITION 1.4.9 (Geodesics). Let $M$ be a manifold and $\nabla$ a covariant derivative on $T M$. A differentiable curve $\gamma: I \rightarrow M$ is said to be a geodesic if $\dot{\gamma} \in \Gamma\left(\gamma^{*} T M\right)$ is parallel along $\gamma$ with respect to $\nabla$, in other words, if $\nabla_{D} \dot{\gamma}=0$.

REMARK 1.4.10 (Covariant derivative in coordinates). Let $\nabla$ be a covariant derivative on a vector bundle $E$ over $M$ and $e_{1}, \ldots, e_{m} \in \Gamma(E / U)$ be a local frame for $E$. Let $(h, U)$ be a local chart for $M$ and $\partial_{i}=\frac{\partial}{\partial h^{i}}$ for $i=1, \ldots, n$. Then

$$
\nabla_{\partial_{i}} e_{j}=\sum_{k} \Gamma_{i j}^{k} e_{k} \quad \text { locally on } U
$$

We consider a section $X \in \Gamma\left(\alpha^{*} E\right)$ along a differentiable curve $\alpha: I \rightarrow M$. Fixing $t_{0} \in I$ such that $\alpha\left(t_{0}\right) \in U$, then $X=\sum_{j=1}^{m} X^{j} \alpha^{*} e_{j}$ locally about $t_{0}$. By Definition 1.4.7 we get for $t$ locally about $t_{0}$ (since $\left.\dot{\alpha}(t)=\sum_{i=1}^{n} \dot{\alpha}^{i}(t)\left(\partial_{i}\right)_{\alpha(t)}\right)$ :

$$
\begin{aligned}
\left(\nabla_{D} X\right)(t) & =\sum_{j=1}^{m}\left(\dot{X}^{j}(t)\left(e_{j}\right)_{\alpha(t)}+X^{j}(t) \nabla_{\dot{\alpha}(t)} e_{j}\right) \\
& =\sum_{j=1}^{m}\left(\dot{X}^{j}(t)\left(e_{j}\right)_{\alpha(t)}+\sum_{i=1}^{n} X^{j}(t) \dot{\alpha}^{i}(t) \nabla_{\left(\partial_{i}\right)_{\alpha(t)}} e_{j}\right)
\end{aligned}
$$

Thus locally about $t_{0}$ :

$$
\begin{equation*}
\nabla_{D} X=\sum_{k}\left(D\left(X^{k}\right)+\sum_{i, j} X^{j} D\left(\alpha^{i}\right)\left(\Gamma_{i j}^{k} \circ \alpha\right)\right) \alpha^{*} e_{k} \tag{1.4.3}
\end{equation*}
$$

THEOREM 1.4.11. Let $\nabla$ be a covariant derivative on a vector bundle $E$ over $M$, further let $\alpha: I \rightarrow M$ be a differentiable curve, $t_{0} \in I$ and $e \in E_{\alpha\left(t_{0}\right)}$. There exists exactly one section $X \in \Gamma_{\mathrm{par}}\left(\alpha^{*} E\right)$ along $\alpha$ such that $X\left(t_{0}\right)=e$.

Proof. The claim is reduced to the existence and uniqueness theorem for linear differential equations. Since it is sufficient to consider the local situation, we may assume the existence of a global chart $(h, M)$ for $M$ and a global frame $e_{1}, \ldots, e_{m} \in \Gamma(E)$ for $E$. Then there are uniquely determined coefficients $b^{i} \in \mathbb{R}$ such that $e=\sum_{i=1}^{m} b^{i} e_{i}$. Defining $c_{k j}:=-\sum_{i=1}^{n} \dot{\alpha}^{i}\left(\Gamma_{i j}^{k} \circ \alpha\right) \in C^{\infty}(I)$, by Eq. (1.4.3) the requirement $\nabla_{D} X=0$ together with $X_{t_{0}}=e$ is seen to be equivalent to the system of linear differential equations

$$
\begin{equation*}
\dot{X}^{k}=-\sum_{j} c_{k j} X^{j}, \quad X^{k}\left(t_{0}\right)=b^{k}, \quad k=1, \ldots, m \tag{1.4.4}
\end{equation*}
$$

It remains to recall that the unique solution to Eq. (1.4.4) is defined on all of $I$.
DEFINITION 1.4.12 (Parallel transport; induced by $\nabla$ ). Let $\nabla$ be a covariant derivative on a vector bundle $E$ over a manifold $M$ and $\alpha: I \rightarrow M$ a differentiable curve. For $s, t \in I$ there is an isomorphism

$$
/ /_{s, t}: E_{\alpha(s)} \rightarrow E_{\alpha(t)}
$$

explained by $/ /_{s, t} e:=X(t)$ where $X \in \Gamma_{\text {par }}\left(\alpha^{*} E\right)$ is the unique parallel section along $\alpha$ such that $X(s)=e$. The isomorphism $/ / s, t$ is called the parallel transport of $E_{\alpha(s)}$ to $E_{\alpha(t)}$ along $\alpha$.

REMARK 1.4.13. We have $/ /_{s, t}^{-1}=/ / t, s$ and $/ /_{t, t}=\operatorname{id}_{E_{\alpha(t)}}$. Each basis $e_{1}, \ldots, e_{m}$ of $E_{\alpha(s)}$ can be extended to a global frame $\bar{e}_{1}, \ldots, \bar{e}_{m} \in \Gamma\left(\alpha^{*} E\right)$ for $\alpha^{*} E$ via $\bar{e}_{i, t}:=/ / s, t e_{i}$.

REMARK 1.4.14. The parallel transport associated to a covariant derivative $\nabla$ according to Definition 1.4.12 defines a parallel transport in $E$ in the sense of Definition 1.4.1. On the other hand the parallel transport determines again the underlying covariant derivative: If $X \in \Gamma(E), v \in T_{p} M$ and $\alpha: I \rightarrow M$ a differentiable curve such that $\dot{\alpha}(0)=v$, then

$$
\nabla_{v} X=\left.\frac{d}{d t}\right|_{t=0}\left(/ / /_{0, t}^{-1} X_{\alpha(t)}\right) \in E_{p}
$$

Proof. Let $e_{1}, \ldots, e_{m}$ be a basis of $E_{p}$ and $\bar{e}_{1}, \ldots, \bar{e}_{m} \in \Gamma\left(\alpha^{*} E\right), \bar{e}_{i, t}:=/ /_{0, t} e_{i}$, an extension to a global frame for $\alpha^{*} E$. Furthermore let $a^{i} \in C^{\infty}(I)$ be such that $\alpha^{*} X=$ $\sum a^{i} \bar{e}_{i}$. Then $/ /_{0, t}^{-1}\left(\alpha^{*} X\right)_{t}=\sum a^{i}(t) e_{i}$, and hence $\nabla_{v} X=\nabla_{D}\left(\alpha^{*} X\right)_{0}=\sum\left(\dot{a}^{i}(0) e_{i}+\right.$ $\left.a^{i}(0)\left(\nabla_{D} \bar{e}_{i}\right)_{0}\right)=\left.\frac{d}{d t}\right|_{t=0}\left(/ /{ }_{0, t}^{-1} X_{\alpha(t)}\right)$.

Thus a parallel transport in $E$ and a covariant derivation on $E$ provide identical structures on $E$. We continue with a third equivalent point of view.

DEFINITION 1.4.15 (Horizontal splitting of $T E$ ). Let $\pi: E \rightarrow M$ be a vector bundle over a manifold $M$. A subbundle $H \subset T E$ is said to be a horizontal splitting of $T E$ if the following two conditions hold:
(i) $T E=H \oplus \pi^{*} E$ (this is $T_{e} E=H_{e} \oplus E_{\pi(e)}$ for $e \in E$ )
(ii) For $s \in \mathbb{R}^{*} \equiv \mathbb{R} \backslash\{0\}$ the subbundle $H$ is compatible with the operation $\rho_{s}: E \rightarrow E$, $e \mapsto s e$, in the sense that $\left(\rho_{s}\right)_{*} H_{e}=H_{s e}$ for $e \in E$ and $s \in \mathbb{R}^{*}$.

REMARK 1.4.16. Let $\pi: E \rightarrow M$ be a vector bundle over $M$ and $e \in E$ such that $p=\pi(e)$. The projection $\pi: E \rightarrow M$ is submersive at $e$, i.e. $(d \pi)_{e}: T_{e} E \rightarrow T_{p} M$ is surjective, with $\operatorname{ker}(d \pi)_{e}=T_{e}\left(\pi^{-1} p\right)=T_{e}\left(E_{p}\right) \cong E_{p} \subset T_{e} E$. Thus there is an exact sequence of vector bundles over $E$

and the decomposition $T E=H \oplus \pi^{*} E$ induces a splitting of (1.4.5): $d \pi \circ h=\mathrm{id}$ where $h=(d \pi \mid H)^{-1}$. The differentiable splitting $h: \pi^{*} T M \xrightarrow{\sim} H \subset T E$ of the sequence (1.4.5) of vector bundles over $M$ is called horizontal lift. Fiberwise, we have linear isomorphisms $(d \pi)_{e} \mid H_{e}: H_{e} \xrightarrow{\sim} T_{p} M$ and $h_{e}: T_{p} M \xrightarrow{\sim} H_{e}$.

Notation 1.4.17. Let $T E=H \oplus V$ with $V:=\pi^{*} E$ be a horizontal splitting of $T E$; further let $w \in T_{e} E$ where $e \in E$ and $p=\pi(e) \in M$. We call $w$ horizontal if $w \in H_{e}$, and vertical (in the sense "tangential to the submanifold $E_{p}$ of $E$ ") if $w \in V_{e} \equiv E_{p}$.

In each fiber $T_{e} E$ of $T E$ the vertical space $V_{e}$ is canonically given, however in general there is no canonical choice of a horizontal space $H_{e}$ : a horizontal splitting provides exactly a selection of a horizontal complement $H_{e}$ to $V_{e}$ at each $e \in E$ in an $\mathbb{R}^{*}$-invariant way (differentiable depending on $e$ ).


Figure 1.4.2. Horizontal splitting
THEOREM 1.4.18. Horizontal splittings of TE and covariant derivatives on $E$ are equivalent structures. For any vector bundle $\pi: E \rightarrow M$ over $M$, the following holds:
(i) A covariant derivative $\nabla$ on $E$ defines canonically a horizontal splitting $H$ of $T E$, namely for each $e \in E$ with $p:=\pi(e)$ via

$$
H_{e}:=\left\{X_{*} v: v \in T_{p} M, X \in \Gamma(E) \text { with } X(p)=e \text { and } \nabla_{v} X=0\right\} \subset T_{e} E .
$$

(ii) Inversely to (i) a horizontal splitting $H$ of TE gives rise to a covariant derivation $\nabla$ on E as follows:
(a) For $X \in \Gamma(E)$ the covariant derivative $\nabla X \in \Gamma\left(T^{*} M \otimes E\right)$ is explained through the following homomorphism of vector bundles over $M$
$T M \xrightarrow{X_{*}} X^{*} T E \equiv X^{*} H \oplus X^{*} V \xrightarrow{\text { pr }_{V}} X^{*} V=X^{*} \pi^{*} E=E$, where $\mathrm{pr}_{V}$ denotes the projection onto the vertical subspace.
(b) For $\sigma \in \Gamma\left(\alpha^{*} E\right)$ where $\alpha$ is a differentiable curve in $M$ the covariant derivative $\nabla_{D} \sigma \in \Gamma\left(\alpha^{*} E\right)$ is given as follows: $\left(\nabla_{D} \sigma\right)\left(t_{0}\right)$ is the image of $\left(\frac{\partial}{\partial t}\right)_{t_{0}}$ under

$$
\mathbb{R}=T \mathbb{R} \xrightarrow{\sigma_{*}} \sigma^{*} T E \equiv \sigma^{*} H \oplus \sigma^{*} V \xrightarrow{\mathrm{pr}_{V}} \sigma^{*} V=\sigma^{*} \pi^{*} E=E ;
$$

here $\left(\frac{\partial}{\partial t}\right)_{t_{0}}$ is first mapped to $\dot{\sigma}\left(t_{0}\right)$ and then projected on the vertical component.
The constructions in (i) and (ii) are inverse to each other.
Proof. (i) First of all, $H_{e}$ as defined in (i) is a vector space. Indeed, to each $v \in T_{p} M$ there exists exactly one $w \in H_{e}$ such that $\pi_{*} w=v$, in other words, if $v \in T_{p} M$ and $X, \tilde{X} \in \Gamma(E)$ such that $X(p)=\tilde{X}(p)=e$ and $\nabla_{v} X=\nabla_{v} \tilde{X}=0$, then $X_{*} v=\tilde{X}_{*} v$. Since this is a local statement at $p$, it is sufficient to consider the situation $E=U \times \mathbb{R}^{m}$ with $p \in U \subset \mathbb{R}^{n}$, where then $X: U \rightarrow \mathbb{R}^{m}, \tilde{X}: U \rightarrow \mathbb{R}^{m}$ and $\tilde{X}=A X$ for some differentiable map $A: U \rightarrow \operatorname{GL}(m ; \mathbb{R}), A(p)=$ identity matrix. For $v \in T_{p} M$, one obtains from $\nabla_{v} \tilde{X}=0$ together with $\nabla_{v} X=0$ the equation $v(A) X_{p}=0$ where $v$ in $v(A)$ is applied as derivation componentwise to the matrix function $A$. This shows, as claimed, $d(A X)_{p} v=v(A) X_{p}+A(p) v(X)=A(p) v(X)=v(X)=(d X)_{p} v$. In particular, this shows

$$
\begin{equation*}
H_{e}=X_{*} T_{p} M=(d X)_{p} T_{p} M \tag{1.4.6}
\end{equation*}
$$

in terms of a fixed section $X \in \Gamma(E)$ such that

$$
\begin{equation*}
X_{p}=e \quad \text { and } \quad(\nabla X)_{p}=0 \tag{1.4.7}
\end{equation*}
$$

from where the vector space structure of $H_{e}$ is obvious. Existence of a section $X$ with property (1.4.7) is immediate: it is sufficient to construct $X$ locally about $p$ and to extend it
then to a smooth global section but locally in coordinates about $p$ condition (1.4.7) reduces to find a function with prescribed 1-jet at the single point $p$.

Injectivity of $(d X)_{p}$ follows from $\pi_{*}(d X)_{p}=\mathrm{id} \mid T_{p} M$ and implies in particular $\operatorname{dim} H_{e}=\operatorname{dim} M$. Also $H_{e} \cap V_{e}=\{0\}$ is obvious since $w \in H_{e}$, say $w=X_{*} v$, implies $\pi_{*} w=(\pi \circ X)_{*} v=v$ whereas $w \in V_{e}$ just means that $\pi_{*} w=0$. This proves $T_{e} E=H_{e} \oplus V_{e}$.

It remains to check that $H$ defines a subbundle of $T E$, i.e., that $H_{e}$ depends differentiably on $e \in E$. To this end, we fix a local chart $(h, U)$ for $M$ and assume without restriction of generality that $E \cong U \times \mathbb{R}^{m}$. If $\partial_{i}=\frac{\partial}{\partial h^{i}}$ is one of the basis vector fields over $U$ and $X \in \Gamma(E)$ a non-vanishing section on $U$, i.e., $X: U \rightarrow \mathbb{R}^{m}$ differentiable and $X(p) \neq 0$ for all $p \in U$, then there exists a $C^{\infty}$ function $A: U \rightarrow \mathrm{GL}(m ; \mathbb{R})$ such that

$$
\begin{equation*}
\nabla_{\partial_{i}}(A X)=0 \text { on } U . \tag{1.4.8}
\end{equation*}
$$

Note that condition (1.4.8) is equivalent to

$$
\begin{equation*}
\partial_{i}(A) X+A \nabla_{\partial_{i}} X=0 \tag{1.4.9}
\end{equation*}
$$

which gives a differential equation for $A$. For fixed $p \in U$ and $g \in \operatorname{GL}(m ; \mathbb{R})$ let now $A=A_{i, p, g}: U \rightarrow \operatorname{GL}(m ; \mathbb{R})$ denote the solution to (1.4.9) satisfying $A(p)=g X(p)$. Furthermore, choose for each $e \in E$ a matrix $g(e) \in \mathrm{GL}(m ; \mathbb{R})$ depending differentiably on $e$ such that $e=g(e) X_{\pi(e)}$. This construction gives to each $e \in E$ vector fields

$$
X_{i}^{(e)}:=A_{i, \pi(e), g(e)} X \in \Gamma(T U), \quad 1 \leq i \leq n
$$

and induced vector fields on $E$, namely

$$
\bar{\partial}_{i} \in \Gamma(E), \quad\left(\bar{\partial}_{i}\right)_{e}:=d\left(X_{i}^{(e)}\right)_{\pi(e)}\left(\partial_{i}\right)_{\pi(e)}, \quad 1 \leq i \leq n
$$

such that $\left(\left(\bar{\partial}_{1}\right)_{e}, \ldots,\left(\bar{\partial}_{n}\right)_{e}\right)$ gives a basis for $H_{e}$ for each $e \in E$.
Finally it is easy to see that $H$ is compatible with the operation $\mathbb{R}^{*}$ which completes the proof of part (i) of Theorem 1.4.18.
(ii) The second part can be checked in an elementary way; verification of the product rule requires the $\mathbb{R}^{*}$-invariance of $H$ (condition (ii) in Definition 1.4.15).

According to Theorem 1.4 .18 (ii) a section $X \in \Gamma(E)$ is parallel (i.e., $\nabla_{v} X=0$ for all $v \in T M)$ if and only if $X_{*} v$ is horizontal for any $v \in T M$. In the same way, a section $\sigma \in \Gamma\left(\alpha^{*} E\right)$ along $\alpha: I \rightarrow M$ is parallel (i.e., $\nabla_{D} \sigma=0$ ) if and only if $\dot{\sigma}(t) \in H_{\sigma(t)}$ for all $t \in I$. Hence, as consequence of Theorem 1.4.11, we have the following result.

THEOREM 1.4.19. Let $\pi: E \rightarrow M$ be a vector bundle over a differentiable manifold $M$ and $H$ a horizontal splitting of TE. Furthermore let $\alpha: I \rightarrow M$ be a differentiable curve and $e \in E_{\alpha\left(t_{0}\right)}$ for some $t_{0} \in I$. Then there exists exactly one lift of $\alpha$ to a "horizontal curve" $u: I \rightarrow E$ above $\alpha$ with $u\left(t_{0}\right)=e$, i.e. such that $\pi \circ u=\alpha, u\left(t_{0}\right)=e$ and $\dot{u}(t) \in H_{u(t)}$ for $t \in I$.

DEFINITION 1.4.20 (Linear connection). Let $\pi: E \rightarrow M$ be a vector bundle over a differentiable manifold $M$. A linear connection in $E$ is a covariant derivative on $E$ (or equivalently, a parallel transport in $E$ or a horizontal splitting of $T E$ ). Linear connections in $T M$ are simply called linear connections on $M$.

Let $M$ be an $n$-dimensional manifold equipped with a linear connection $\nabla$ in $T M$. By Definition 1.4.9, geodesics are curves with the property that their velocity field along
the curve is parallel. According to Remark 1.4.10, in local coordinates $(h, U)$, for a differentiable curve $\gamma: I \rightarrow M$ the condition $\nabla_{D} \dot{\gamma}=0$ means that

$$
\begin{equation*}
\ddot{\gamma}^{k}+\sum_{i, j}\left(\Gamma_{i j}^{k} \circ \gamma\right) \dot{\gamma}^{i} \dot{\gamma}^{j}=0, \quad k=1, \ldots, n \tag{1.4.10}
\end{equation*}
$$

with $\gamma^{k}=h^{k} \circ \gamma$ and $\Gamma_{i j}^{k}$ the Christoffel symbols determined by $\nabla_{\partial_{i}} \partial_{j}=\sum_{k} \Gamma_{i j}^{k} \partial_{k}$. According to Theorem 1.4.18, condition $\nabla_{D} \dot{\gamma}=0$ is equivalent to $\ddot{\gamma}(t) \in H_{\dot{\gamma}(t)}$ for each $t \in I$. The horizontal lift $h: \pi^{*} T M \sim H \subset T T M$ induced by $\nabla$ according to diagram (1.4.5), defines a canonical (horizontal) vector field $\xi$ on $T M$, namely for $v \in T M$ via

$$
T_{\pi(v)} M \xrightarrow{h_{v}} H_{v} \longleftrightarrow T_{v} T M, \quad v \longmapsto \xi(v) .
$$

Obviously, $\xi$ is not only a second order differential equation (i.e., a vector field on $T M$ such that $(d \pi)_{v} \xi(v)=v$ for $\left.v \in T M\right)$, it is even a spray which means that in addition $\xi(s v)=\left(d \varrho_{s}\right)(s \xi(v))$ holds for all $s \in \mathbb{R}^{*}$ and the multiplication $\varrho_{s}: T M \rightarrow T M$, $v \mapsto s v$. The vector field $\xi$ on $T M$ is called the geodesic spray to the linear connection $\nabla$ on $T M$. In general, for a second order differential equation $\xi$, curves $\gamma: I \rightarrow M$ such that $\ddot{\gamma}(t)=\xi(\dot{\gamma}(t))$ are called integral curves of $\xi$. Since $\pi \circ \dot{\gamma}=\gamma$, the relation $\pi_{*} \ddot{\gamma}=\dot{\gamma}$ holds trivially, so that in the case of the geodesic spray $\xi$ :

$$
\nabla_{D} \dot{\gamma}=0 \Longleftrightarrow \ddot{\gamma}(t)=\xi(\dot{\gamma}(t)) \quad \text { for any } t \in I
$$

Thus, a curve $\gamma$ in $M$ is a geodesic if and only if $\gamma$ is an integral curve of the corresponding geodesic spray.

COROLLARY 1.4.21. Let $M$ be a manifold and $\nabla$ a linear connection in $T M$. Given $p \in M$ and $v \in T_{p} M$, there exists a unique geodesic $\gamma=\gamma_{v}$ such that $\gamma(0)=p$ and $\dot{\gamma}(0)=v$. In addition, $\gamma_{s v}(t)=\gamma_{v}(s t)$ for $s, t \in \mathbb{R}$, if one of the two sides is defined.

Proof. In general, given a second order differential equation $\xi$, there exists to each $v \in T_{p} M$ a unique maximal integral curve $\gamma$ such that $\dot{\gamma}(0)=v$ (in particular then $\gamma(0)=$ $p)$; this is an immediate consequence of the existence and uniqueness result for integral curves to vector fields. The condition that $\xi$ is even a spray guarantees the addition: the integral curve to $s$-times the initial velocity corresponds to the original integral curve run through $s$-times as fast.

DEFINITION 1.4.22. Let $M$ be a smooth manifold. A linear connection $\nabla$ in $T M$ is called metrically complete if every maximal geodesic is defined on all of $\mathbb{R}$.

DEfinition 1.4.23 (Tensor field). Let $T M$ and $T^{*} M$ be the tangent bundle, resp., cotangent bundle of a differentiable manifold $M$. For $r, s \in \mathbb{N} \cup\{0\}$ the elements of $\Gamma\left(T^{*} M^{\otimes s} \otimes T M^{\otimes r}\right)$ are called tensor fields of type $(r, s)$ or $(r, s)$-tensors in short. For $s \in \mathbb{N}$, in terms of the canonical $C^{\infty}(M)$ - linear isomorphism
$\Gamma\left(T^{*} M^{\otimes s} \otimes T M\right) \cong \Gamma\left(\operatorname{Mult}_{\mathbb{R}}\left(T M^{s} ; T M\right)\right) \cong \operatorname{Mult}_{C^{\infty}(M)}\left(\Gamma(T M)^{s} ; \Gamma(T M)\right)$,
tensor fields of type $(1, s)$ correspond to $C^{\infty}(M)$-multilinear maps $\Gamma(T M)^{s} \rightarrow \Gamma(T M)$.
DEFINITION 1.4.24 (Torsion, curvature). Let $M$ be a differentiable manifold and $\nabla$ a linear connection on $M$.
(i) The map

$$
T: \Gamma(T M)^{2} \rightarrow \Gamma(T M), \quad(X, Y) \mapsto \nabla_{X} Y-\nabla_{Y} X-[X, Y]
$$

is $C^{\infty}(M)$-bilinear and represents a (1,2)-tensor $T \in \Gamma\left(T^{*} M^{\otimes 2} \otimes T M\right)$, the so-called torsion tensor of $\nabla$. Recall that for two vector fields $X, Y \in \Gamma(T M)$ the Lie product $[X, Y] \in \Gamma(T M)$ is defined as derivation via

$$
[X, Y] f=X(Y f)-Y(X f), \quad f \in C^{\infty}(M)
$$

The connection $\nabla$ is said to be torsion-free or symmetric if $T \equiv 0$.
(ii) The map

$$
R: \Gamma(T M)^{3} \rightarrow \Gamma(T M), \quad(X, Y, Z) \mapsto \nabla_{X} \nabla_{Y} Z-\nabla_{Y} \nabla_{X} Z-\nabla_{[X, Y]} Z
$$

is $C^{\infty}(M)$-trilinear and represents a $(1,3)$-tensor $R \in \Gamma\left(T^{*} M^{\otimes 3} \otimes T M\right)$, the curvature tensor of the connection $\nabla$. The tensor $R$ may be written as $C^{\infty}(M)$-bilinear map

$$
R: \Gamma(T M)^{2} \rightarrow \operatorname{End}_{C^{\infty}(M)} \Gamma(T M) \cong \Gamma(\operatorname{End} T M)
$$

and gives then a section $R \in \Gamma\left(T^{*} M^{\otimes 2} \otimes \operatorname{End} T M\right)$. This leads to the common notation $R(X, Y) Z \equiv R(X, Y, Z)$.

REMARK 1.4.25. A linear connection $\nabla$ in $T M$ is torsion-free if and only if $\Gamma_{i j}^{k}=\Gamma_{j i}^{k}$ for all Christoffel symbols. In particular, a linear connection $\nabla$ in $T M$ can be "symmetrized" to a torsion-free connection by passing from $\nabla$ to $\bar{\nabla}$ with the new Christoffel symbols $\bar{\Gamma}_{i j}^{k}:=\frac{1}{2}\left(\Gamma_{i j}^{k}+\Gamma_{j i}^{k}\right)$.

Proof. If $(h, U)$ is a local chart for $M$ and $\partial_{1}, \ldots, \partial_{n}$ the corresponding coordinate vector fields, then on $U$ the Christoffel symbols are determined by $\nabla_{\partial_{i}} \partial_{j}=\sum_{k} \Gamma_{i j}^{k} \partial_{k}$. By the $C^{\infty}$-bilinearity of $T$, the condition $T\left(\partial_{i}, \partial_{j}\right)=0$ for all coordinate vector fields $\partial_{i}, \partial_{j}$ implies already $T(X, Y) \mid U=0$ for any $X, Y \in \Gamma(T M)$. Since $\left[\partial_{i}, \partial_{j}\right] \equiv 0$ we have $T\left(\partial_{i}, \partial_{j}\right)=\nabla_{\partial_{i}} \partial_{j}-\nabla_{\partial_{j}} \partial_{i}$ which proves the claim.

REMARK 1.4.26. Let $\nabla$ be a linear connection in $T M$. The symmetrized connection $\bar{\nabla}$ defined in Remark 1.4.25 is given by

$$
\bar{\nabla}_{X} Y=\nabla_{X} Y-\frac{1}{2} T(X, Y)=\frac{1}{2}\left(\nabla_{X} Y+\nabla_{Y} X+[X, Y]\right), \quad X, Y \in \Gamma(T M)
$$

Let $f: M \rightarrow N$ be a differentiable map between manifolds. Reading the curvature tensor $R$ of a linear connection in $T N$ as an element of $\Gamma\left(T^{*} N^{\otimes 2} \otimes \operatorname{End} T N\right)$ and taking into account that $f^{*} \operatorname{End} T N \cong \operatorname{End}\left(f^{*} T N\right)$, we obtain

$$
f^{*} R \in \Gamma\left(T^{*} M^{\otimes 2} \otimes f^{*} \operatorname{End} T N\right) \cong \operatorname{Bil}_{C^{\infty}(M)}\left(\Gamma(T M), \Gamma(T M) ; \operatorname{End}_{C^{\infty}(M)} f^{*} T N\right)
$$

In explicit terms, for $A, B \in \Gamma(T M), Y \in \Gamma\left(f^{*} T N\right)$ and $p \in M$,

$$
\left(\left(f^{*} R\right)(A, B) Y\right)_{p}=R_{f(p)}\left(d f_{p} A_{p}, d f_{p} B_{p}, Y_{p}\right) \in T_{f(p)} N
$$

THEOREM 1.4.27 (Cartan's structural equations). Let $f: M \rightarrow N$ be a differentiable map between manifolds and $\nabla$ a linear connection in $T N$. Then, for $A, B \in \Gamma(T M)$, $Y \in \Gamma\left(f^{*} T N\right)$,

$$
\begin{aligned}
\left(f^{*} T\right)(A, B) & =\nabla_{A}(d f B)-\nabla_{B}(d f A)-d f[A, B] \in \Gamma\left(f^{*} T N\right) \\
\left(f^{*} R\right)(A, B) Y & =\nabla_{A} \nabla_{B} Y-\nabla_{B} \nabla_{A} Y-\nabla_{[A, B]} Y \in \Gamma\left(f^{*} T N\right)
\end{aligned}
$$

(On the right-hand sides $\nabla$ corresponds to the induced covariant derivative on $f^{*} T M$; see Definition 1.4.7).

Proof. It is sufficient to verify the two equations locally. To this end, let $(h, U)$ be a chart for $N$ and $\partial_{1}, \ldots, \partial_{d} \in \Gamma(T N / U)$ the corresponding local frame for $T N$. Then also $f^{*} \partial_{1}, \ldots, f^{*} \partial_{d} \in \Gamma\left(f^{*} T N / f^{-1}(U)\right)$ is a local frame for $f^{*} T N$ over $M$ and on $f^{-1}(U)$ we have

$$
\begin{array}{r}
d f A=\sum A\left(h^{i} \circ f\right) f^{*} \partial_{i}, \quad d f B=\sum B\left(h^{i} \circ f\right) f^{*} \partial_{i} \\
d f[A, B]=\sum\left(A\left(B\left(h^{i} \circ f\right)\right)-B\left(A\left(h^{i} \circ f\right)\right)\right) f^{*} \partial_{i} .
\end{array}
$$

From this one obtains furthermore (always on $f^{-1}(U)$ )

$$
\begin{aligned}
& \nabla_{A}(d f B)=\sum\left(A\left(B\left(h^{i} \circ f\right)\right) f^{*} \partial_{i}+B\left(h^{i} \circ f\right) \nabla_{A} f^{*} \partial_{i}\right) \\
& \nabla_{B}(d f A)=\sum\left(B\left(A\left(h^{i} \circ f\right)\right) f^{*} \partial_{i}+A\left(h^{i} \circ f\right) \nabla_{B} f^{*} \partial_{i}\right) .
\end{aligned}
$$

On the other hand, we have $T\left(\partial_{i}, \partial_{j}\right)=\nabla_{\partial_{i}} \partial_{j}-\nabla_{\partial_{j}} \partial_{i}$ and hence

$$
\begin{aligned}
\left(f^{*} T\right)(A, B) & =\sum_{i, j} A\left(h^{i} \circ f\right) B\left(h^{j} \circ f\right) f^{*}\left(T\left(\partial_{i}, \partial_{j}\right)\right) \\
& =\sum_{i}\left(B\left(h^{i} \circ f\right) \nabla_{A} f^{*} \partial_{i}-A\left(h^{i} \circ f\right) \nabla_{B} f^{*} \partial_{i}\right)
\end{aligned}
$$

This shows the first structural equation. The verification of the second equation is similar.

DEFINITION 1.4.28 (Covariant derivative of a differential form). Let $M$ be a manifold and $\nabla$ a linear connection in $T M$. For a differential form $\alpha \in A^{1}(M) \equiv \Gamma\left(T^{*} M\right)$ and a vector field $A \in \Gamma(T M)$ let $\nabla_{A} \alpha \in A^{1}(M)$ be defined as

$$
\begin{equation*}
\left(\nabla_{A} \alpha\right)(B):=A(\alpha B)-\alpha\left(\nabla_{A} B\right), \quad B \in \Gamma(T M) \tag{1.4.11}
\end{equation*}
$$

Note that $\nabla_{A} \alpha$ is well-defined by Lemma 1.4.2, as the right-hand side of (1.4.11) is $C^{\infty}(M)$-linear in $B$. We may write

$$
\nabla \alpha \in \Gamma\left(T^{*} M \otimes T^{*} M\right) \equiv \Gamma(\operatorname{Bil}(T M, T M) ; \mathbb{R}), \quad \nabla \alpha(A, B):=\left(\nabla_{A} \alpha\right)(B)
$$

DEfinition 1.4.29 (Hessian). For $f \in C^{\infty}(M)$ the covariant derivative of $\alpha=d f$,

$$
\operatorname{Hess}(f):=\nabla d f \in \Gamma\left(T^{*} M \otimes T^{*} M\right), \quad(\nabla d f)(A, B)=A B f-\left(\nabla_{A} B\right) f
$$

is called second fundamental form (Hessian) of $f$.
REMARK 1.4.30. Let $M$ be a manifold and $\nabla$ a linear connection in $T M$. Then

$$
\nabla d f \in \operatorname{Bil}_{C \infty(M)}\left(\Gamma(T M), \Gamma(T M) ; C^{\infty}(M)\right), \quad(A, B) \mapsto(\nabla d f)(A, B)
$$

is symmetric for each $f \in C^{\infty}(M)$ if and only if $\nabla$ is torsion-free, i.e.,

$$
T(A, B) \equiv \nabla_{A} B-\nabla_{B} A-[A, B]=0
$$

for all $A, B \in \Gamma(T M)$.
Example 1.4.31. If $M=\mathbb{R}^{n}$ and $\nabla$ the canonical connection on $\mathbb{R}^{n}$ defined by $\nabla_{D_{i}} D_{j}=0$, then $(\nabla d f)\left(D_{i}, D_{j}\right)=D_{i} D_{j} f$.

We now turn to a central concept of the stochastic calculus on manifolds, the notion of manifold-valued martingales.

DEFINITION 1.4.32 ( $\nabla$-martingale). Let $M$ be a manifold and $\nabla$ be a linear connection in $T M$. Further let $X$ be an $M$-valued semimartingale defined on some filtered probability space $\left(\Omega ; \mathscr{F} ; \mathbb{P} ;\left(\mathscr{F}_{t}\right)_{t \geq 0}\right)$. Then $X$ is called $\nabla$-martingale (or simply martingale) if for any $f \in C^{\infty}(M)$ :

$$
\begin{equation*}
d(f(X)) \stackrel{m}{=} \frac{1}{2}(\nabla d f)(d X, d X) \tag{1.4.12}
\end{equation*}
$$

where $\stackrel{\underline{m}}{=}$ means equality modulo differentials of local martingales.

REmark 1.4.33. Since $(\nabla d f)(d X, d X)$ in Eq. (1.4.12) only depends on the symmetric part of $\nabla d f$, we may always assume that the linear connection $\nabla$ is torsion-free. Symmetrization of the connection does not change the class of $\nabla$-martingales.

A priori, martingales on $M$ may be defined only up to some predictable stopping time. Since the concept of a martingale is invariant under time transformation (see Remark 1.3.6) and since by an appropriate time transformation infinite (or deterministic finite) lifetime can be achieved, we neglect this point in the notation.

EXAMPLE 1.4.34. In the special case of $M=\mathbb{R}^{n}$ equipped with the canonical linear connection $\nabla$, we have $(\nabla d f)\left(D_{i}, D_{j}\right)=D_{i} D_{j} f$, and hence $\nabla$-martingales in the sense of Definition 1.4.32 coincide with the usual class of continuous local martingales on $\mathbb{R}^{n}$. Indeed, according to Itô's formula, a continuous $\mathbb{R}^{n}$-valued semimartingale $X$ is a local martingale if and only if

$$
d(f(X))-\frac{1}{2} \sum_{i, j}\left(D_{i} D_{j} f\right)(X) d\left[X^{i}, X^{j}\right] \in d \mathscr{M}
$$

for all $f \in C^{\infty}\left(\mathbb{R}^{n}\right)$ which is exactly condition (1.4.12) of Definition 1.4.32.
REMARK 1.4.35 ( $\nabla$-martingales as solutions of SDEs). Let $\nabla$ be a linear connection on $T M$ which without loss of generality is torsion-free. Let $A_{0}, A_{1}, \ldots, A_{r} \in \Gamma(T M)$ and suppose that $X$ solves the SDE

$$
\begin{equation*}
d X=A_{0}(X) d t+\sum_{i=1}^{r} A_{i}(X) \circ d Z^{i} \tag{1.4.13}
\end{equation*}
$$

Here $Z$ may be an arbitrary continuous $\mathbb{R}^{r}$-valued semimartingale. Then for $f \in C^{\infty}(M)$ we have

$$
d(f(X))=\left(A_{0} f\right)(X) d t+\sum_{i=1}^{r}\left(A_{i} f\right)(X) d Z^{i}+\frac{1}{2} \sum_{i, j=1}^{r}\left(A_{i} A_{j} f\right)(X) d\left[Z^{i}, Z^{j}\right]
$$

Since $(\nabla d f)\left(A_{i}, A_{j}\right)=A_{i} A_{j} f-\left(\nabla_{A_{i}} A_{j}\right) f$ and since on the other hand

$$
(\nabla d f)(d X, d X)=\sum_{i, j=1}^{r}(\nabla d f)\left(A_{i}, A_{j}\right)(X) d\left[Z^{i}, Z^{j}\right]
$$

we obtain

$$
\begin{aligned}
& d(f(X))-\frac{1}{2}(\nabla d f)(d X, d X) \\
& \quad=\left(A_{0} f\right)(X) d t+\sum_{i=1}^{r}\left(A_{i} f\right)(X) d Z^{i}+\frac{1}{2} \sum_{i, j=1}^{r}\left(\nabla_{A_{i}} A_{j} f\right)(X) d\left[Z^{i}, Z^{j}\right]
\end{aligned}
$$

Denoting the drift of the semimartingale $Z$ by $Z^{\text {drift }}$, we obtain that $X$ is a $\nabla$-martingale if and only if for any $f \in C^{\infty}(M)$,

$$
\left(A_{0} f\right)(X) d t+\sum_{i=1}^{r}\left(A_{i} f\right)(X) d\left(Z^{\mathrm{drift}}\right)^{i}+\frac{1}{2} \sum_{i, j=1}^{r}\left(\nabla_{A_{i}} A_{j} f\right)(X) d\left[Z^{i}, Z^{j}\right]=0
$$

In the special case when $Z$ is a Brownian motion on $\mathbb{R}^{r}$ we find that solutions $X$ to the SDE (1.4.13) are $\nabla$-martingales if

$$
A_{0}=-\frac{1}{2} \sum_{i=1}^{r} \nabla_{A_{i}} A_{i}
$$

### 1.5. Riemannian Metrics and Brownian Motions

The measurement of the distance between points and the length of curves on a manifold requires as additional structure a metric on the tangent bundle. Manifolds equipped with a metric are called Riemannian manifolds. Such a structure is also needed for the notion of Brownian motions on manifolds.

Definition 1.5.1 (Riemannian metric). Let $E$ be a vector bundle over M. A Riemannian metric on $E$ is a section

$$
g \in \Gamma\left(E^{*} \otimes E^{*}\right) \cong \Gamma(\operatorname{Bil}(E, E ; \mathbb{R})) \cong \operatorname{Bil}_{C^{\infty}(M)}\left(\Gamma(E), \Gamma(E) ; C^{\infty}(M)\right)
$$

such that $g_{x} \in \operatorname{Bil}\left(E_{x}, E_{x} ; \mathbb{R}\right)$ is symmetric and positive definite for any $x \in M$.
We often write $\langle\cdot, \cdot\rangle$ instead of $g$ and then $\langle\cdot, \cdot\rangle_{x}$ for the scalar product $g_{x}$ on the fiber $E_{x}$ (depending differentiably on $x$ in bundle charts). For a section $A \in \Gamma(E)$ we use the notation $|A|$ for $\sqrt{g(A, A)}$ (and write $|e|_{x}$ instead of $\sqrt{g_{x}(e, e)}$ for $e \in E_{x}$ ).

DEFINITION 1.5.2 (Riemannian manifold). A Riemannian manifold is a pair $(M, g)$ consisting of a differentiable manifold $M$ and a Riemannian metric $g$ on the tangent bundle $T M$.

DEFINITION 1.5.3 (Length of curves). Let $\alpha:[a, b] \rightarrow M$ be a piecewise differentiable curve on a Riemannian manifold $(M, g)$ such that $\alpha \mid\left[t_{i-1}, t_{i}\right]$ is differentiable for some partition $a=t_{0}<t_{1}<\cdots<t_{r}=b$ of the interval $[a, b]$. Then

$$
L(\alpha):=\sum_{i=1}^{r} \int_{t_{i-1}}^{t_{i}}|\dot{\alpha}(t)|_{\alpha(t)} d t
$$

is well-defined and called the length of $\alpha$.
DEFINITION 1.5.4 (Isometry; local isometry). Let $(M, g)$ and ( $N, h$ ) be Riemannian manifolds. A differentiable map $f: M \rightarrow N$ is called local isometry if $g=f^{*} h$, i.e., if $g_{p}(u, v)=\left(f^{*} h\right)_{p}(u, v) \equiv h_{f(p)}\left((d f)_{p} u,(d f)_{p} v\right)$ for all $u, v \in T_{p} M$. If in addition $f$ is a diffeomorphisms, then $f$ is called an isometry.

The condition $g=f^{*} h$ means that, for any $p \in M$, the map $(d f)_{p}:\left(T_{p} M, g_{p}\right) \rightarrow$ $\left(T_{f(p)} N, h_{f(p)}\right)$ is an isometry of Euclidean vector spaces. In particular, local isometries let the length of curves invariant.

DEFINITION 1.5.5 (Riemannian connection). Let $(M, g)$ be a Riemannian manifold and $\nabla$ a linear connection in $T M$. Then $\nabla$ is called a Riemannian connection if all parallel transports

$$
/ /_{s, t}:\left(T_{\alpha(s)} M, g_{\alpha(s)}\right) \rightarrow\left(T_{\alpha(t)} M, g_{\alpha(t)}\right)
$$

along differentiable curves $\alpha$ are isometries.
THEOREM 1.5.6 (Characterization of Riemannian connections). Let $(M, g)$ be a Riemannian manifold and $\nabla$ a linear connection in TM. The following items equivalent:
(i) $\nabla$ is a Riemannian connection.
(ii) (Ricci identity) For all $X, Y, Z \in \Gamma(T M)$,

$$
Z\langle X, Y\rangle=\left\langle\nabla_{Z} X, Y\right\rangle+\left\langle X, \nabla_{Z} Y\right\rangle
$$

(iii) If $f: N \rightarrow M$ is a differentiable map, then for $X, Y \in \Gamma\left(f^{*} T M\right), A \in \Gamma(T N)$,

$$
A\langle X, Y\rangle=\left\langle\nabla_{A} X, Y\right\rangle+\left\langle X, \nabla_{A} Y\right\rangle
$$

Proof. (i) $\Rightarrow$ (ii) Let $p \in M$ and $\alpha: I \rightarrow M$ a differentiable curve such that $\dot{\alpha}(0)=$ $Z_{p}$. In terms of the parallel transport $/ /_{s, t}$ along $\alpha$ we calculate

$$
\begin{aligned}
Z_{p}\langle X, Y\rangle & =\left.\frac{d}{d t}\right|_{t=0}\left\langle X_{\alpha(t)}, Y_{\alpha(t)}\right\rangle_{\alpha(t)} \\
& =\left.\frac{d}{d t}\right|_{t=0}\left\langle\sigma_{0, t}^{-1} X_{\alpha(t)}, \sigma_{0, t}^{-1} Y_{\alpha(t)}\right\rangle_{p} \\
& =\left\langle\left.\frac{d}{d t}\right|_{t=0}\left(\sigma_{0, t}^{-1} X_{\alpha(t)}\right), Y_{p}\right\rangle_{p}+\left\langle X_{p},\left.\frac{d}{d t}\right|_{t=0}\left(\sigma_{0, t}^{-1} Y_{\alpha(t)}\right)\right\rangle_{p} \\
& =\left\langle\left(\nabla_{Z} X\right)_{p}, Y_{p}\right\rangle_{p}+\left\langle X_{p},\left(\nabla_{Z} Y\right)_{p}\right\rangle_{p}
\end{aligned}
$$

(ii) $\Rightarrow$ (iii) Let $X, Y \in \Gamma\left(f^{*} T M\right)$ and $A \in \Gamma(T N)$. Since for $\phi \in C^{\infty}(N)$,

$$
\begin{gathered}
A\langle\phi X, Y\rangle=(A \phi)\langle X, Y\rangle+\phi A\langle X, Y\rangle \quad \text { and } \\
\left\langle\nabla_{A}(\phi X), Y\right\rangle=(A \phi)\langle X, Y\rangle+\phi\left\langle\nabla_{A} X, Y\right\rangle
\end{gathered}
$$

it is sufficient to verify the statement for $X=f^{*} U, Y=f^{*} V$ where $U, V \in \Gamma(T M)$. But with $q \in N$ and $w:=A_{q}$, we obtain

$$
\begin{aligned}
A_{q}\langle X, Y\rangle & =w(\langle U, V\rangle \circ f)=\left(d f_{q} w\right)\langle U, V\rangle \\
& =\left\langle\nabla_{d f_{q} w} U, V_{f(q)}\right\rangle+\left\langle U_{f(q)}, \nabla_{d f_{q} w} V\right\rangle=\left\langle\nabla_{w} X, Y_{q}\right\rangle+\left\langle X_{q}, \nabla_{w} Y\right\rangle
\end{aligned}
$$

(iii) $\Rightarrow$ (i) Let $X, Y$ be parallel vector fields along a differentiable curve $\alpha$ in $M$. Then $D\langle X, Y\rangle=\left\langle\nabla_{D} X, Y\right\rangle+\left\langle X, \nabla_{D} Y\right\rangle=0$ which shows that $\langle X, Y\rangle$ is constant.

Geodesics with respect to Riemannian connections are parameterized proportionally to arc length. Indeed, we have $D\langle\dot{\gamma}, \dot{\gamma}\rangle=2\left\langle\nabla_{D} \dot{\gamma}, \dot{\gamma}\right\rangle=0$, and hence $|\dot{\gamma}|$ is constant.

THEOREM 1.5.7 (of Levi-Civita). On a Riemannian manifold $(M, g)$ there exists a unique torsion-free Riemannian connection $\nabla$ in $T M$.

Proof. For uniqueness it is sufficient to show that $\left\langle\nabla_{X} Y, Z\right\rangle$ is uniquely determined for $X, Y, Z \in \Gamma(T M)$. Indeed, from the Ricci identity we obtain

$$
\begin{aligned}
X\langle Y, Z\rangle & =\left\langle\nabla_{X} Y, Z\right\rangle+\left\langle Y, \nabla_{X} Z\right\rangle \\
Y\langle X, Z\rangle & =\left\langle\nabla_{Y} X, Z\right\rangle+\left\langle X, \nabla_{Y} Z\right\rangle \\
Z\langle X, Y\rangle & =\left\langle\nabla_{Z} X, Y\right\rangle+\left\langle X, \nabla_{Z} Y\right\rangle .
\end{aligned}
$$

Adding the first two equations and subtracting the last one, along with the torsion-freeness of $\nabla$, gives

$$
\begin{align*}
2\left\langle\nabla_{X} Y, Z\right\rangle= & X\langle Y, Z\rangle+Y\langle X, Z\rangle-Z\langle X, Y\rangle \\
& +\langle[X, Y], Z\rangle-\langle[X, Z], Y\rangle-\langle[Y, Z], X\rangle \tag{1.5.1}
\end{align*}
$$

The right-hand side of this equation is $C^{\infty}(M)$-linear in $Z$ and determines the vector field $\nabla_{X} Y \in \Gamma(T M)$. It is straightforward to check that it defines a torsion-free Riemannian connection in $T M$.

DEFINITION 1.5.8. The unique torsion-free Riemannian connection in $T M$ for a Riemannian manifold $(M, g)$ according to Theorem 1.5 .7 is called Levi-Civita connection in TM and the associated parallel transport the Levi-Civita parallelism.

REMARK 1.5.9. Eq. (1.5.1) can be used to express the Levi-Civita connection of a Riemannian manifold $(M, g)$ directly via the metric $g$. To this end, let $(h, U)$ be a local chart for $M$ and $\partial_{i}=\frac{\partial}{\partial h^{i}} \in \Gamma(T M / U)$ for $i=1, \ldots, n$. Consider

$$
g_{i j}:=\left\langle\frac{\partial}{\partial h^{i}}, \frac{\partial}{\partial h^{j}}\right\rangle \in C^{\infty}(U),
$$

and $g^{i j} \in C^{\infty}(U)$ with $\sum_{j} g^{i j} g_{j k}=\delta_{i k}$. Then, by means of Eq. (1.5.1)

$$
2\left\langle\nabla_{\partial_{i}} \partial_{j}, \partial_{k}\right\rangle=\partial_{i}\left\langle\partial_{j}, \partial_{k}\right\rangle+\partial_{j}\left\langle\partial_{k}, \partial_{i}\right\rangle-\partial_{k}\left\langle\partial_{i}, \partial_{j}\right\rangle
$$

i.e., $2 \sum_{m} \Gamma_{i j}^{m} g_{m k}=\partial_{i} g_{j k}+\partial_{j} g_{k i}-\partial_{k} g_{i j}$ from where the wanted relation follows:

$$
\begin{equation*}
\Gamma_{i j}^{\ell}=\frac{1}{2} \sum_{k} g^{k \ell}\left(\partial_{i} g_{j k}+\partial_{j} g_{k i}-\partial_{k} g_{i j}\right) \tag{1.5.2}
\end{equation*}
$$

EXAMPLE 1.5.10 (Levi-Civita connection on $\left.\mathbb{R}^{n}\right)$. Let $(M, g)=\left(\mathbb{R}^{n}\right.$, eucl) with the canonical Riemannian metric

$$
\operatorname{eucl}(A, B) \equiv\langle A, B\rangle=\sum_{i=1}^{n} A^{i} B^{i}
$$

for vector fields $A, B$ on $\mathbb{R}^{n}$. Vector fields on $\mathbb{R}^{n}$ are interpreted equally as functions in $C^{\infty}\left(\mathbb{R}^{n}, \mathbb{R}^{n}\right)$ and as derivations on $C^{\infty}\left(\mathbb{R}^{n}\right)$ : the constant maps $\left(x \mapsto e_{i}\right) \in C^{\infty}\left(\mathbb{R}^{n}, \mathbb{R}^{n}\right)$, $i=1, \ldots, n$, correspond to the derivations $D_{1}, \ldots, D_{n}$ where $D_{i}=\frac{\partial}{\partial x^{i}}$.
(i) According to Eq. (1.5.2), the Levi-Civita connection $\nabla$ on $\left(\mathbb{R}^{n}\right.$, eucl) is determined by $\nabla_{D_{i}} D_{j}=0$ which for vector fields $A, B$ on $\mathbb{R}^{n}$ means that

$$
\begin{equation*}
\nabla_{A} B=\sum_{i} A\left(B^{i}\right) D_{i} \equiv \sum_{i} A\left(B^{i}\right) e_{i}=\left(A\left(B^{1}\right), \ldots, A\left(B^{n}\right)\right) \tag{1.5.3}
\end{equation*}
$$

This connection is also denoted by $D$, i.e. $\nabla_{A} B=D_{A} B$, since according to (1.5.3), $\nabla_{v} B=D_{v} B$ coincides with the directional derivative of $B \in C^{\infty}\left(\mathbb{R}^{n}, \mathbb{R}^{n}\right)$ in direction $v \in \mathbb{R}^{n}$.
(ii) Linear connections induced by the Levi-Civita connection on $\left(\mathbb{R}^{n}\right.$, eucl) are easy to determine. For instance, let $M$ be a manifold and $f: M \rightarrow \mathbb{R}^{n}$ a differentiable map. To describe the induced covariant derivative on $f^{*} T \mathbb{R}^{n}$, we note that each section $A \in$ $\Gamma\left(f^{*} T \mathbb{R}^{n}\right)$ writes as $\sum_{i=1}^{n} A^{i} f^{*} D_{i}$ with $A^{i} \in C^{\infty}(M)$. Then, for $X \in \Gamma(T M)$, the covariant derivative $\nabla_{X} A=\sum_{i=1}^{n}\left(\nabla_{X} A\right)^{i} f^{*} D_{i} \in \Gamma\left(f^{*} T \mathbb{R}^{n}\right)$ with respect to the induced linear connection in $f^{*} T \mathbb{R}^{n}$ is given by

$$
\left(\nabla_{X} A\right)^{i}=X A^{i}, \quad i=1, \ldots, n
$$

Indeed, by the product rule, we have

$$
\nabla_{X} A=\sum\left(X A^{i}\right) f^{*} D_{i}+A^{i} \nabla_{X} f^{*} D_{i}
$$

where $\nabla_{X} f^{*} D_{i}=D_{d f X} D_{i}=0$.
PROPOSITION 1.5.11 (Levi-Civita connections on Riemannian submanifolds). Let $(M, g)$ be a Riemannian submanifold of a Riemannian manifold $(\tilde{M}, \tilde{g})$ in the sense that there is an embedding $M \stackrel{\iota}{\hookrightarrow}$ such that $g=\iota^{*} \tilde{g}$. The homomorphism $d \iota_{x}: T_{x} M \rightarrow$ $T_{x} \tilde{M}$ then is an isometry for each $x \in M$. For $x \in M$ let $\mathrm{pr}_{x}^{\top}: T_{x} \tilde{M} \rightarrow T_{x} M$ denote the orthogonal projection onto the linear subspace $T_{x} M \equiv d \iota_{x} T_{x} M \subset T_{x} \tilde{M}$. This gives $a$ homomorphism of vector bundles $\mathrm{pr}^{\top}: \iota^{*} T \tilde{M} \rightarrow T M$ over $M$ with $\mathrm{pr}^{\top} \circ d \iota=\mathrm{id}_{T M}$ :


Let now $\nabla$, respectively $\tilde{\nabla}$, denote the Levi-Civita connection in $T M$, respectively $T \tilde{M}$. Then, for $A, B \in \Gamma(T M)$,

$$
\begin{equation*}
\nabla_{A} B=\operatorname{pr}^{\top} \tilde{\nabla}_{A}(d \iota B) \tag{1.5.4}
\end{equation*}
$$

In particular, $\nabla$ is uniquely determined by $\tilde{\nabla}$.
Proof. First of all, we remark that the right-hand side of (1.5.4) defines a linear connection $\nabla^{\prime}$ in $T M$. By the uniqueness part of Theorem 1.5.7 it is hence enough to show that $\nabla^{\prime}$ is a torsion-free Riemannian connection. Denoting by $T^{\prime}$, respectively $\tilde{T}$, the torsion tensor of $\nabla^{\prime}$, respectively $\tilde{\nabla}$, then we have by Theorem 1.4.27 for $A, B \in \Gamma(T M)$ :

$$
\begin{aligned}
T^{\prime}(A, B) & =\nabla_{A}^{\prime} B-\nabla_{B}^{\prime} A-[A, B] \\
& =\operatorname{pr}^{\top}\left(\tilde{\nabla}_{A}(d \iota B)-\tilde{\nabla}_{B}(d \iota A)-d \iota[A, B]\right) \\
& =\operatorname{pr}^{\top}\left(\iota^{*} \tilde{T}\right)(A, B)=0
\end{aligned}
$$

Here we used that $\tilde{T} \equiv 0$ since $\tilde{\nabla}$ is torsion-free, and hence $\nabla^{\prime}$ is also torsion-free. Furthermore, by Theorem 1.5.6, we obtain for $A, B, C \in \Gamma(T M)$,

$$
\begin{aligned}
\left\langle\nabla_{C}^{\prime} A, B\right\rangle+\left\langle A, \nabla_{C}^{\prime} B\right\rangle & =\left\langle\tilde{\nabla}_{C}(d \iota A), d \iota B\right\rangle+\left\langle d \iota A, \tilde{\nabla}_{C}(d \iota B)\right\rangle \\
& =C\langle d \iota A, d \iota B\rangle=C\langle A, B\rangle .
\end{aligned}
$$

Thus $\nabla^{\prime}$ satisfies the Ricci identity and is therefore a Riemannian connection according to Theorem 1.5.6.

Example 1.5.12 (Riemannian submanifolds of ( $\mathbb{R}^{n}$, eucl)). We specialize Proposition 1.5 .11 to the case of a Riemannian submanifold of $\mathbb{R}^{n}$. Then $(M, g) \stackrel{\iota}{\longleftrightarrow}\left(\mathbb{R}^{n}\right.$, eucl) with $g=\iota^{*}$ eucl. By means of the fiberwise isometric bundle embedding $d \iota: T M \hookrightarrow$ $\iota^{*} T \mathbb{R}^{n}$ we have $T M \subset \iota^{*} T \mathbb{R}^{n} \cong M \times \mathbb{R}^{n}$ as a vector subbundle and then $\Gamma(T M) \subset$ $\Gamma\left(\iota^{*} T \mathbb{R}^{n}\right)=C^{\infty}\left(M ; \mathbb{R}^{n}\right)$ : vector fields on $M$ are hereby $\mathbb{R}^{n}$-valued $C^{\infty}$-maps on $M$. In terms of the orthogonal projection $\mathrm{pr}^{\top}: M \times \mathbb{R}^{n} \rightarrow T M,(x, v) \mapsto \mathrm{pr}_{x}^{\top} v$, we then have according to (1.5.3) and (1.5.4),

$$
\begin{equation*}
\nabla_{A} B=\operatorname{pr}^{\top}\left(A B^{1}, \ldots, A B^{n}\right), \quad A, B \in \Gamma(T M) \tag{1.5.5}
\end{equation*}
$$

REMARK 1.5.13. Let $(M, g)$ be a Riemannian manifold and $\nabla$ a Riemannian connection in $T M$. For $f \in C^{\infty}(M)$ consider $\operatorname{grad} f \in \Gamma(T M)$, the gradient of $f$, defined by

$$
\langle\operatorname{grad} f, A\rangle=A f, \quad A \in \Gamma(T M)
$$

Note that

$$
(\nabla d f)(A, B)=\left\langle\nabla_{A} \operatorname{grad} f, B\right\rangle
$$

Indeed, according to Theorem 1.5.6 (ii) (Ricci identity) we have

$$
A\langle\operatorname{grad} f, B\rangle=\left\langle\nabla_{A} \operatorname{grad} f, B\right\rangle+\left\langle\operatorname{grad} f, \nabla_{A} B\right\rangle
$$

and hence $\left\langle\nabla_{A} \operatorname{grad} f, B\right\rangle=A B f-\left(\nabla_{A} B\right) f$.
THEOREM 1.5.14 (Martingales on submanifolds of ( $\mathbb{R}^{n}$, eucl)). Let $M \hookrightarrow \mathbb{R}^{n}$ be a submanifold of $\mathbb{R}^{n}$ endowed with the induced metric $g=\iota^{*}$ eucl and let $\nabla$ be the LeviCivita connection on $M$. Suppose that $X$ is an $M$-valued semimartingale and $\bar{X}=\iota(X)$ its embedding into $\mathbb{R}^{n}$. Let

$$
\bar{X}=\bar{X}_{0}+N+C
$$

be the Doob-Meyer decomposition of $\bar{X}$ in $\mathbb{R}^{n}$, where $N \in \mathscr{M}_{0}\left(\mathbb{R}^{n}\right)$ and $C \in \mathscr{A}_{0}\left(\mathbb{R}^{n}\right)$. We consider $T_{X_{t}} M$ as a linear subspace of $\mathbb{R}^{n}$. Then $X$ is $a \nabla$-martingale if and only if

$$
\begin{equation*}
d C_{t} \perp T_{X_{t}} M \quad \text { for all } t \in \mathbb{R}_{+}, \text {a.s. } \tag{1.5.6}
\end{equation*}
$$

where the last condition is understood in the sense that $\left\langle H_{t}, d C_{t}\right\rangle=0$ for each piecewise continuous process $H$ such that $H_{t} \in T_{X_{t}} M$ a.s. In particular, each continuous local martingale on $\mathbb{R}^{n}$ taking values in $M$ is a $\nabla$-martingale on $M$.

Proof. For $f \in C^{\infty}(M)$ we denote by $\bar{f} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ a continuation of $f$ to $\mathbb{R}^{n}$, i.e. $f=\bar{f} \circ \iota$. For a vector field $A \in \Gamma(T M)$ let $\bar{A}=\iota_{*} A \in \Gamma\left(\iota^{*} T \mathbb{R}^{n}\right) \equiv C^{\infty}\left(M ; \mathbb{R}^{n}\right)$. Calculating the Hessian with respect to the Levi-Civita connection $\nabla$ on ( $M, \iota^{*}$ eucl), we then get for $A, B \in \Gamma(T M)$,

$$
\begin{equation*}
(\nabla d f)(A, B)=\left(\operatorname{Hess}_{\mathbb{R}^{n}} \bar{f}\right)(\bar{A}, \bar{B}) \tag{1.5.7}
\end{equation*}
$$

Indeed, for $f \in C^{\infty}(M)$, we note that

$$
\iota_{*} \operatorname{grad} f=\left(D_{1} \bar{f}, \ldots, D_{n} \bar{f}\right) \mid M \in C^{\infty}\left(M ; \mathbb{R}^{n}\right)
$$

Identifying $A \in \Gamma(T M)$ and $\iota_{*} A \in \Gamma\left(\iota^{*} T \mathbb{R}^{n}\right)=C^{\infty}\left(M ; \mathbb{R}^{n}\right)$ then gives according to Eq. (1.5.5),

$$
\begin{aligned}
(\nabla d f)(A, B) & =\left\langle\nabla_{A} \operatorname{grad} f, B\right\rangle_{T M}=\left\langle\operatorname{pr}^{\top}\left(A\left(D_{1} \bar{f}\right), \ldots, A\left(D_{n} \bar{f}\right)\right), B\right\rangle_{T M} \\
& =\left\langle\left(A\left(D_{1} \bar{f}\right), \ldots, A\left(D_{n} \bar{f}\right)\right), B\right\rangle_{\mathbb{R}^{n}} \\
& =\sum_{i, j} A^{i} B^{j} D_{i} D_{j} \bar{f}=\left(\operatorname{Hess}_{\mathbb{R}^{n}} \bar{f}\right)(\bar{A}, \bar{B})
\end{aligned}
$$

For $h \in C^{\infty}(M)$ with $\bar{h} \in C^{\infty}\left(\mathbb{R}^{n}\right)$, we thus obtain by Itô's formula along with pullback formula (1.3.4),

$$
\begin{aligned}
d(h(X))=d(\bar{h}(\bar{X})) & =\sum_{i=1}^{n} D_{i} \bar{h}(\bar{X}) d \bar{X}^{i}+\frac{1}{2} \sum_{i, j=1}^{n} D_{i} D_{j} \bar{h}(\bar{X}) d \bar{X}^{i} d \bar{X}^{j} \\
& =\left\langle\left(\operatorname{grad}_{\mathbb{R}^{n}} \bar{h}\right)(\bar{X}), d \bar{X}\right\rangle+\frac{1}{2}\left(\operatorname{Hess}_{\mathbb{R}^{n}} \bar{h}\right)(d \bar{X}, d \bar{X}) \\
& \stackrel{\mathrm{m}}{=}\left\langle\left(\operatorname{grad}_{\mathbb{R}^{n}} \bar{h}\right)(\bar{X}), d C\right\rangle+\frac{1}{2}(\nabla d h)(d X, d X)
\end{aligned}
$$

Hence $X$ is a $\nabla$-martingale if and only if $\left\langle\left(\operatorname{grad}_{\mathbb{R}^{n}} \bar{h}\right)(\bar{X}), d C\right\rangle=0$ for all continuations $\bar{h} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ of functions $h \in C^{\infty}(M)$. Applied to the coordinate functions $h^{1}, \ldots, h^{n}$ of the embedding $M \stackrel{\iota}{\longleftrightarrow} \mathbb{R}^{n}$ this gives the claim.

For the rest of this section let $(M, g)$ be a Riemannian manifold equipped with the Levi-Civita connection $\nabla$.

DEFINITION 1.5.15 (Riemannian quadratic variation). Let $X$ be a semimartingale taking values in a Riemannian manifold $(M, g)=(M,\langle\cdot, \cdot\rangle)$. The process

$$
[X, X]:=\int g(d X, d X)=\int\langle d X, d X\rangle
$$

is called Riemannian quadratic variation of $X$.
DEFINITION 1.5.16 (Laplace-Beltrami operator). Let $(M, g)$ be a Riemannian manifold and $\nabla$ the Levi-Civita connection on $M$. For $f \in C^{\infty}(M)$ let

$$
\Delta f:=\operatorname{trace} \nabla d f \in C^{\infty}(M)
$$

where $\nabla d f \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$ denotes the Hessian of $f$. In other words, $(\Delta f)(x)=$ $\sum_{i}(\nabla d f)\left(e_{i}, e_{i}\right)$ where $e_{1}, \ldots, e_{n}$ is some orthonormal basis for $T_{x} M$. The operator $\Delta$ is called Laplace-Beltrami operator on $M$.

In local coordinates $(h, U)$ for $M$,

$$
\begin{equation*}
\nabla d f \mid U=\sum_{i, j}\left(\partial_{i} \partial_{j} f-\sum_{k} \Gamma_{i j}^{k} \partial_{k} f\right) d h^{i} \otimes d h^{j} \tag{1.5.8}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\Delta f \mid U=\sum_{i, j} g^{i j}\left(\partial_{i} \partial_{j} f-\sum_{k} \Gamma_{i j}^{k} \partial_{k} f\right) \tag{1.5.9}
\end{equation*}
$$

where $g_{i j}=g\left(\partial_{i}, \partial_{j}\right)$ and $\left(g^{i j}\right) \equiv\left(g_{i j}\right)^{-1}$; here we used that $g\left(\partial_{i}, \cdot\right)=\sum_{j} g_{i j} d h^{j}$ or equivalently $\sum_{i} g^{j i} g\left(\partial_{i}, \cdot\right)=d h^{j}$. In particular, we see that $\Delta$ is a second order differential operator on $M$ (without constant term).

Definition 1.5.17 (Brownian motion on a Riemannian manifold). Let $(M, g)$ be a Riemannian manifold and $X$ an adapted $M$-valued process with maximal lifetime $\zeta$, defined on a filtered probability space $\left(\Omega ; \mathscr{F} ; \mathbb{P} ;\left(\mathscr{F}_{t}\right)_{t \in \mathbb{R}_{+}}\right)$satisfying the usual conditions. The process $X$ is called a Brownian motion on $(M, g)$ if, for any $f \in C^{\infty}(M)$, the realvalued process

$$
f(X)-\frac{1}{2} \int(\Delta f)(X) d t
$$

is a local martingale (with lifetime $\zeta$ ). The class of Brownian motions on $(M, g)$ is denoted by $\operatorname{BM}(M, g)$.

THEOREM 1.5.18 (Lévy's characterization of $M$-valued Brownian motions). Let $X$ be a semimartingale with maximal lifetime taking values in a Riemannian manifold $(M, g)$. The following conditions are equivalent:
(i) $X$ is $\mathrm{BM}(M, g)$.
(ii) $X$ is a $\nabla$-martingale with the property that $[f(X), f(X)]=\int|\operatorname{grad} f|^{2}(X) d t$ for every $f \in C^{\infty}(M)$.
(iii) $X$ is $a \nabla$-martingale with the property that $\int b(d X, d X)=\int($ trace $b)(X) d t$ for every $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$.
In particular, for the Riemannian quadratic variation $[X, X]=\int g(d X, d X) \equiv$ $\int\langle d X, d X\rangle$ of a Brownian motion $X$ on $M$, we get

$$
\int_{0}^{t} g(d X, d X)=n t
$$

where $n=\operatorname{dim} M$.
Proof. A. We verify that for an arbitrary $M$-valued semimartingale $X$ the following two conditions are equivalent:
(a) $[f(X), f(X)]=\int|\operatorname{grad} f|^{2}(X) d t$
(b) $\int b(d X, d X)=\int($ trace $b)(X) d t$ for every $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$.

In particular, this shows (ii) $\Longleftrightarrow$ (iii). Indeed, for $f, h \in C^{\infty}(M)$ we have

$$
\begin{aligned}
\operatorname{trace}(d f \otimes d h) & =\sum_{i}(d f \otimes d h)\left(e_{i}, e_{i}\right)=\sum_{i}(d f)\left(e_{i}\right)(d h)\left(e_{i}\right) \\
& =\sum_{i}\left\langle\operatorname{grad} f, e_{i}\right\rangle\left\langle\operatorname{grad} h, e_{i}\right\rangle=\langle\operatorname{grad} f, \operatorname{grad} h\rangle .
\end{aligned}
$$

Thus $(\mathrm{b}) \Rightarrow(\mathrm{a})$ is just the special case for $b=d f \otimes d f$. To show the converse direction (a) $\Rightarrow$ (b), first note that (a) implies by polarization

$$
[f(X), h(X)]=\int\langle\operatorname{grad} f, \operatorname{grad} h\rangle(X) d t, \quad f, h \in C^{\infty}(M)
$$

Thus $[f(X), h(X)]=\int(d f \otimes d h)(d X, d X)=\int \operatorname{trace}(d f \otimes d h)(X) d t$. By means of the uniqueness part of Theorem 1.3.3, we get

$$
\int b(d X, d X)=\int(\operatorname{trace} b)(X) d t
$$

for any bilinear form $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$.
B. (iii) $\Rightarrow$ (i): Part A applied to the given $\nabla$-martingale $X$ shows $b(d X, d X)=($ trace $b)(X) d t$ for bilinear forms $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$, in particular for $b=\nabla d f$,

$$
d(f(X)) \stackrel{\mathrm{m}}{=} \frac{1}{2} \nabla d f(d X, d X)=\frac{1}{2}(\Delta f)(X) d t
$$

thus $X$ is $\operatorname{BM}(M, g)$.
C. (i) $\Rightarrow$ (ii): Now let $X$ be $\operatorname{BM}(M, g)$ and $f \in C^{\infty}(M)$. According to $\nabla d f^{2}=$ $2(f \nabla d f+d f \otimes d f)$ we first note $\Delta\left(f^{2}\right)=2 f \Delta f+2|\operatorname{grad} f|^{2}$, thus

$$
d\left(f^{2}(X)\right) \stackrel{m}{=} \frac{1}{2}\left(\Delta f^{2}\right)(X) d t=(f \Delta f)(X) d t+|\operatorname{grad} f|^{2}(X) d t
$$

On the other hand, by means of Itô's formula,

$$
\begin{aligned}
d\left(f^{2}(X)\right) & =2 f(X) d(f(X))+d[f(X), f(X)] \\
& \stackrel{\text { m }}{=} f(X)(\Delta f)(X) d t+d[f(X), f(X)]
\end{aligned}
$$

Uniqueness in the Doob-Meyer decomposition implies

$$
[f(X), f(X)]=\int|\operatorname{grad} f|^{2}(X) d t
$$

Finally, once again by means of part A, the last formula gives

$$
\nabla d f(d X, d X)=(\operatorname{trace} \nabla d f)(X) d t=(\Delta f)(X) d t
$$

from where we conclude that $X$ is a $\nabla$-martingale.
Example 1.5.19. According to Lévy's characterization of flat Brownian motions, Brownian motions on ( $\mathbb{R}^{n}$, eucl) in the sense of Definition 1.5.17 coincide with the usual class of $\mathbb{R}^{n}$-valued Brownian motions.

THEOREM 1.5.20 ( $M$-valued Brownian motions as solutions of an SDE). Let ( $M, g$ ) be a Riemannian manifold and $\nabla$ the Levi-Civita connection on $M$. We consider the SDE

$$
\begin{equation*}
d X=A_{0}(X) d t+A(X) \circ d Z \tag{1.5.10}
\end{equation*}
$$

with $A_{0} \in \Gamma(T M), A \in \Gamma\left(\operatorname{Hom}\left(M \times \mathbb{R}^{r}, T M\right)\right)$, and $Z$ a BM on $\mathbb{R}^{r}$. Then maximal solutions to (1.5.10) are Brownian motions on $(M, g)$ if the two subsequent conditions are satisfied:
(i) $A_{0}=-\frac{1}{2} \sum_{i} \nabla_{A_{i}} A_{i}$ with $A_{i} \equiv A(\cdot) e_{i}$ for $i=1, \ldots, r$.
(ii) The map $A(x)^{*}: T_{x} M \rightarrow \mathbb{R}^{r}$ is an isometric embedding for every $x \in M$, i.e., $A(x) A(x)^{*}=\mathrm{id}_{T_{x} M}$ where $A(x)^{*}$ is the adjoint to $A(x) \in \operatorname{Hom}\left(\mathbb{R}^{r}, T_{x} M\right)$.

Proof. Let $X$ be a solution to Eq. (1.5.10) and assume that conditions (i) and (ii) are satisfied. According to Remark 1.4.35 condition (i) guarantees that $X$ is a $\nabla$-martingale. In addition, we have for $f \in C^{\infty}(M)$,

$$
d(f(X)) \stackrel{\mathrm{m}}{=} \frac{1}{2} \sum_{i=1}^{r}(\nabla d f)\left(A_{i}, A_{i}\right)(X) d t
$$

It is thus sufficient to verify that

$$
\sum_{i=1}^{r}(\nabla d f)\left(A_{i}, A_{i}\right)=\Delta f
$$

Letting $x \in M$ and $\left(a_{1}, \ldots, a_{n}\right)$ an orthonormal basis of $T_{x} M$, we obtain

$$
\begin{aligned}
(\Delta f)(x) & =\operatorname{trace}(\nabla d f)_{x}=\sum_{i=1}^{n}(\nabla d f)_{x}\left(a_{i}, a_{i}\right) \\
& =\sum_{i=1}^{n}(\nabla d f)_{x}\left(A(x) A(x)^{*} a_{i}, A(x) A(x)^{*} a_{i}\right)
\end{aligned}
$$

Completing $\left(A(x)^{*} a_{1}, \ldots, A(x)^{*} a_{n}\right)$ to an orthonormal basis $\left(\tilde{e}_{1}, \ldots, \tilde{e}_{r}\right)$ of $\mathbb{R}^{r}$, taking into account that $\left(\operatorname{im} A(x)^{*}\right)^{\perp}=\operatorname{ker} A(x)$ and denoting by $\left(e_{1}, \ldots, e_{r}\right)$ the standard basis of $\mathbb{R}^{r}$, we obtain

$$
\begin{aligned}
(\Delta f)(x) & =\sum_{i=1}^{r}(\nabla d f)_{x}\left(A(x) \tilde{e}_{i}, A(x) \tilde{e}_{i}\right) \\
& =\sum_{i=1}^{r}(\nabla d f)_{x}\left(A(x) e_{i}, A(x) e_{i}\right)=\sum_{i=1}^{r}(\nabla d f)_{x}\left(A_{i}(x), A_{i}(x)\right)
\end{aligned}
$$

which completes the proof.
REMARK 1.5.21. Conditions (i) and (ii) of Theorem 1.5.20 can always be satisfied for $r$ sufficiently large. For instance, let $M \hookrightarrow \mathbb{R}^{r}$ be a Whitney embedding. Then $T_{x} M$ can be seen as a subspace $\mathbb{R}^{r}$ for each $x \in M$. Defining $A \in \Gamma\left(\operatorname{Hom}\left(M \times \mathbb{R}^{r}, T M\right)\right)$ fiberwise as orthogonal projection $A(x): \mathbb{R}^{r} \rightarrow T_{x} M$ onto $T_{x} M$ and setting $A_{0}=-\frac{1}{2} \sum_{i} \nabla_{A_{i}} A_{i}$, then every solution to the $\operatorname{SDE}$ (1.5.10) (with a given initial condition) is a Brownian motion on $(M, g)$. The drawback of this construction is that to a given Riemannian manifold $(M, g)$ there is no canonical choice of the coefficients $A_{0}$ and $A$; there is however a canonical SDE on the orthonormal frame bundle $\mathrm{O}(T M)$ over $M$ such that its solutions project to Brownian motions on $(M, g)$. We develop this construction in the next Section.

We conclude this Section with a specification of Theorem 1.5.20 in the case of submanifolds of $\mathbb{R}^{n}$.

THEOREM 1.5.22 (Brownian motions on submanifolds of $\mathbb{R}^{n}$ ). Let $M \hookrightarrow \mathbb{R}^{n}$ be a submanifold of $\mathbb{R}^{n}$ endowed with the induced Riemannian metric $g=\iota^{*}$ eucl. Consider the SDE

$$
\begin{equation*}
d X=A(X) \circ d Z \tag{1.5.11}
\end{equation*}
$$

where $Z$ is a Brownian motion on $\mathbb{R}^{n}$ and

$$
A \in \Gamma\left(\operatorname{Hom}\left(M \times \mathbb{R}^{n}, T M\right)\right), \quad(x, v) \mapsto A(x) v
$$

such that $A(x): \mathbb{R}^{n} \rightarrow T_{x} M$ is the orthogonal projection onto $T_{x} M$. Then every solution to (1.5.11) gives a Brownian motion on $(M, g)$.

Proof. For each $x \in M$, the map $d \iota_{x}: T_{x} M \rightarrow \mathbb{R}^{n}$ is an isometric embedding and we consider $T_{x} M$ as a linear subspace of $\mathbb{R}^{n}$. Note that $A(x)^{*}=d \iota_{x}$ and $A(x) A(x)^{*}=$ $\mathrm{id}_{T_{x} M}$. In terms of the vector fields $A_{i} \equiv A(\cdot) e_{i} \in \Gamma(T M), i=1, \ldots, n$, by Theorem 1.5.20, it is sufficient to verify that $\sum_{i=1}^{n} \nabla_{A_{i}} A_{i}=0$ where $\nabla$ denotes the Levi-Civita connection on $M$. For $e \in \mathbb{R}^{n}$ let $A^{e}:=A(\cdot) e \in \Gamma(T M)$. We show that

$$
\nabla_{v} A^{e}=0 \quad \text { for all } e \in \operatorname{im} A(x)^{*}=(\operatorname{ker} A(x))^{\perp}, v \in T_{x} M
$$

To this end, let $Q(x):=A(x)^{*} A(x): \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$. We note that $Q(x)^{2}=Q(x)$ and $Q(x) e=e$ for $e \in \operatorname{im} A(x)^{*}$. Thus, using $Q Q=Q$, we have

$$
(d Q)_{x} e=(d Q)_{x} Q(x) e+Q(x)(d Q)_{x} e=(d Q)_{x} e+Q(x)(d Q)_{x} e
$$

from where we conclude that $Q(x)(d Q)_{x} e=0$ for all $e \in \operatorname{im} A(x)^{*}$. In explicit terms this shows that

$$
0=A(x)^{*} A(x) d\left(A(\cdot)^{*} A(\cdot) e\right)_{x}=A(x)^{*} A(x) d\left(\bar{A}^{e}\right)_{x}
$$

where $\bar{A}^{e}=\iota_{*} A^{e}$. Since by Eq. (1.5.4), $\nabla A^{e}=A d\left(\bar{A}^{e}\right)$, we finally obtain

$$
\left(\nabla A^{e}\right)_{v}=\nabla_{v} A^{e}=0, \quad v \in T_{x} M
$$

which completes the proof.
REMARK 1.5.23. On a differentiable manifold $M$ consider an SDE of the type

$$
\begin{equation*}
d X=A_{0}(X) d t+A(X) \circ d Z \tag{1.5.12}
\end{equation*}
$$

with $A_{0} \in \Gamma(T M), A \in \Gamma\left(\operatorname{Hom}\left(M \times \mathbb{R}^{r}, T M\right)\right)$, and $Z$ a BM on $\mathbb{R}^{r}$. We assume that the $\operatorname{SDE}$ (1.5.12) is elliptic in the sense that

$$
A(x): \mathbb{R}^{r} \rightarrow T_{x} M \quad \text { is surjective for each } x \in M
$$

Then there exists a Riemannian metric $g$ on $M$ such that $A(x)^{*}: T_{x} M \rightarrow \mathbb{R}^{r}$ is an isometric embedding for each $x \in M$. Indeed, let $Q \in \Gamma\left(T^{*} M \otimes \mathbb{R}^{r}\right)$ be a right-inverse to $A$, i.e. for each $x \in M$,

$$
Q(x): T_{x} M \rightarrow \mathbb{R}^{r} \text { is linear and } A(x) Q(x)=\mathrm{id}_{T_{x} M}
$$

We define

$$
g(u, v):=\langle Q(x) u, Q(x) v\rangle_{\mathbb{R}^{r}}, \quad u, v \in T_{x} M
$$

It's easy to see that $Q(x)=A(x)^{*}$. Now let $\nabla$ be the Levi-Civita connection on $(M, g)$. As in the proof of Theorem 1.5.20 we have

$$
\sum_{i=1}^{r}(\nabla d f)\left(A_{i}, A_{i}\right)=\Delta f
$$

where $A_{i}=A(\cdot) e_{i}$ for $i=1, \ldots, r$. Since $A_{i}^{2} f=(\nabla d f)\left(A_{i}, A_{i}\right)+\nabla_{A_{i}} A_{i}$, we observe that solutions to (1.5.12) define $L$-diffusions for

$$
L=\frac{1}{2} \sum_{i=1}^{r} A_{i}^{2}+A_{0}=\frac{1}{2} \Delta+\left(A_{0}+\sum_{i=1}^{r} \nabla_{A_{i}} A_{i}\right)=\frac{1}{2} \Delta+V
$$

where $V:=A_{0}+\sum_{i=1}^{r} \nabla_{A_{i}} A_{i} \in \Gamma(T M)$. In other words, with respect to an appropriately chosen Riemannian metric $g$ on $M$, maximal solutions to (1.5.12) are Brownian motions on $(M, g)$ with drift $V$.

### 1.6. Parallel Transport and Stochastically Moving Frames

In the last Section we have shown that Brownian motion on a Riemannian manifold $M$ can be constructed as solution to an appropriate stochastic differential equation on $M$ (driven by a standard Euclidean Wiener process). These constructions are however not canonical which is due to the fact that in general, unless the tangent bundle of $M$ is trivial (i.e. for M parallelizable), the Laplace-Beltrami operator does not have a natural representation in Hörmander form as a sum of squares of vector fields.

The fundamental observation that Brownian motions on Riemannian manifolds can be horizontally lifted via a Riemannian connection to semimartingales on the orthonormal frame bundle $\mathrm{O}(T M) \rightarrow M$ over $M$ and satisfy there globally defined canonical stochastic differential equations (SDEs) goes back to the pioneering work of Malliavin, Eells and Elworthy. Conversely, solving SDEs on the frame bundle and projecting the solutions down to the manifold $M$ allows canonical constructions of diffusion processes on $M$ (see [7, 8, 9, 10, 31, 32, 33, 42]).

Intuitively the procedure of constructing $M$-valued processes $X$ from continuous $\mathbb{R}^{n}$ valued semimartingales $Z$ corresponds to a "rolling without slipping" of the manifold $M$ along the trajectories of $Z$ in $\mathbb{R}^{n}$. It allows to construct to each continuous semimartingale $Z$ in $T_{x} M \equiv \mathbb{R}^{n}$ a stochastic development $X$ on $M$, together with a notion of parallel transport along the paths of $X$ on $M$. Brownian motion $X$ on $M$ starting at $x \in M$ can be thought as the trace printed on $M$ by the paths of an Euclidean Brownian motion $Z$ in $T_{x} M \cong \mathbb{R}^{n}$ when " $M$ is rolled along the trajectories of the flat process". The obvious difficulty that paths of Brownian motion are non-differentiable almost surely requires to work with stochastic differential equations instead of pathwise ordinary differential equations.

We begin the discussion with necessary prerequisites on principal bundles and connection forms. Apart from the already studied vector bundles (with a finite-dimensional vector space $V \cong \mathbb{R}^{n}$ as typical fiber) a further type of fiber bundles is needed, that is principal bundles with a Lie group $G$ as typical fiber. The most important examples will be the frame bundle $\mathrm{L}(T M)$ with $G=\mathrm{GL}(n ; \mathbb{R})$ and, in the Riemannian case, the orthonormal frame bundle $\mathrm{O}(T M)$ with the orthogonal group $G=\mathrm{O}(n)$. Vector bundles and principal bundles belong both to the common category of fiber bundles with structure group.

Notation 1.6.1. Let $G$ be a Lie group and $F$ a manifold. A left action of $G$ on $F$ (" $G$ operates on $F$ from the left") is a differentiable map

$$
G \times F \rightarrow F, \quad(g, v) \rightarrow g v=: L_{g} v
$$

with the properties:
(a) $e v=v$ for $v \in F$ where $e$ is the neutral element in $G$,
(b) $g_{2}\left(g_{1} v\right)=\left(g_{2} g_{1}\right) v$ for all $g_{1}, g_{2} \in G$ and $v \in F$.

A left action of $G$ on $F$ is hence a group homomorphism $G \rightarrow \operatorname{Diff}(F)$ with the property that the operation $G \times F \rightarrow F$ is differentiable. A left action of $G$ on $F$ is called effective if $G \rightarrow \operatorname{Diff}(F), g \mapsto L_{g}$ is injective, and free if $g v=v$ for some $v \in F$ implies $g=e$.

If $F=V$ is a finite dimensional real vector space, then left actions of $G$ on $V$ given by a differentiable group homomorphism $G \rightarrow \operatorname{Aut}(V)$ are called linear, respectively representations of $G$, if they are in addition effective.

These concepts carry over correspondingly to right actions of $G$ (" $G$ operates on $F$ from the right"). Note that if $F \times G \rightarrow F,(v, g) \mapsto v g$ is a right action of $G$ on $F$, then $G \times F \rightarrow F,(g, v) \mapsto v g^{-1}$ defines a left action of $G$ on $F$.

DEFINITION 1.6.2 (Fibre bundle with structure group). Let $\pi: E \rightarrow M$ be a fiber bundle with typical fiber $F$ and $G$ a Lie group with an effective left action of $G$ on $F$. A bundle atlas for $\pi: E \rightarrow M$ is called $G$-atlas if all transition functions are given by differentiable maps taking values in $G \subset \operatorname{Diff}(F)$. The change of two charts $\left(\varphi_{i}, U_{i}\right)$, $\left(\varphi_{j}, U_{j}\right)$ is thus given by a differentiable map $\phi_{i j}: U_{i} \cap U_{j} \rightarrow G$ such that

$$
\begin{gathered}
\text { E/(U, } \left.\cap U_{j}\right) \\
\left(U_{i} \cap U_{j}\right) \times F \cdots\left(U_{i} \cap U_{j}\right) \times F \\
(x, v) \longmapsto\left(x, \phi_{i j}(x) v\right)
\end{gathered}
$$

The bundle $\pi: E \rightarrow M$ equipped with a $G$-Atlas is called fiber bundle with typical fiber $F$ and structure group $G$.

As usual, differentiable right inverses of $\pi: E \rightarrow M$ are called sections of $E$. Global resp. local sections of $E$ are denoted by $\Gamma(E)$ resp. $\Gamma(E / U)$.

REMARK 1.6.3. The $m$-dimensional vector bundles over a manifold $M$ are the fiber bundles over $M$ with typical fiber $\mathbb{R}^{m}$ and $\operatorname{GL}(m ; \mathbb{R})$ as structure group; see Definition 1.0.22. In particular, the tangent bundle $T M \rightarrow M$ of a differentiable $n$-dimensional manifold is a fiber bundle with typical fiber $\mathbb{R}^{n}$ and structure group $\operatorname{GL}(n ; \mathbb{R})$ where the transition functions of the canonical $\operatorname{GL}(n ; \mathbb{R})$-atlas take the form $x \mapsto J_{x}\left(k \circ h^{-1}\right)$.

DEFINITION 1.6.4 (Principal bundle). Let $G$ be a Lie group. A principal $G$-bundle is a fiber bundle $\pi: P \rightarrow M$ with typical fiber $G$ and structure group $G$ which operates on $G$ (effectively) from the left by the group multiplication.

Let $\pi: P \rightarrow M$ be a principal $G$-bundle. Then $G$ operates in a natural way from the right on $P$. At first, $G$ operates on $G$ itself on the right by group multiplication, and via charts this action extends to a differentiable free right action on $P$ : an element $u \in P / U$ reads in a chart $P / U \rightarrow U \times G$ as $(x, g)$, so that $u a$ for $a \in G$ corresponds to ( $x, g a$ ),

$$
\begin{aligned}
P / U \longrightarrow U \times G, \quad u & (x, g) \\
u a & \longleftrightarrow(x, g a)
\end{aligned}
$$

and since this assignment is independent of charts, it gives a well-defined global right action of $G$ on $P$. The orbits of this action are the fibers of $P$. Bundle charts $(\varphi, U)$ from a $G$-atlas

are automatically $G$-compatible (or equivariant) in the sense of

$$
\varphi(u a)=\left(\pi(u),\left(\operatorname{pr}_{G} \circ \varphi\right)(u a)\right), \quad u \in \pi^{-1} U, a \in G
$$

DEFINITION 1.6.5 (Reduction of the structure group). If the $G$-Atlas of a fiber bundle $\pi: E \rightarrow M$ with structure group $G$ contains a $G^{\prime}$-Atlas for $E$ where $G^{\prime} \subset G$ is a Lie subgroup, then $\pi: E \rightarrow M$ can also be considered as fiber bundle with structure group $G^{\prime}$. This procedure is called reduction of the structure group $G$ to $G^{\prime}$.

REMARK 1.6.6. Reductions of the structure group appear naturally in case of an additional structure on the fibers of $E$ corresponding in suitable charts to a $G^{\prime}$-invariant structure on the typical fiber. For example, if a vector bundle $\pi: E \rightarrow M$ of rank $m$ carries
a Riemannian metric, then the fiberwise isometric linear bundle charts provide an $\mathrm{O}(m)$ atlas and hence a reduction of the structure group $\mathrm{GL}(m ; \mathbb{R})$ to $\mathrm{O}(n)$. Such bundle charts can be constructed from a $\operatorname{GL}(m ; \mathbb{R})$-atlas for $\pi: E \rightarrow M$ via Gram-Schmidt orthogonalization. In this sense the tangent bundle of an $n$-dimensional Riemannian manifold becomes a vector bundle with structure group $\mathrm{O}(n)$.

REMARK 1.6.7. (a) (Associated fiber bundles) Let $\pi: P \rightarrow M$ be a principal $G$ bundle and $F$ a manifold with an effective left action of $G$ on $F$. We may consider the right action of $G$ (diagonal action) on $P \times F$ given by

$$
(P \times F) \times G \longrightarrow P \times F, \quad(p, v) g:=\left(p g, g^{-1} v\right)
$$

The projection $P \times F \rightarrow P \rightarrow M$ is invariant under the diagonal action, and $E \equiv$ $P \times_{G} F:=(P \times F) / G \rightarrow M$ defines a fiber bundle with typical fiber $F$ and structure group $G$. Each chart $(\varphi, U)$ for $P$ gives a chart for $E$ via

$$
\begin{aligned}
E / U=(P / U) \times_{G} F & \longrightarrow(U \times G) \times_{G} F \xrightarrow{\longrightarrow} U \times F, \\
{[u, v] } & \longmapsto[\varphi(u), v]
\end{aligned}
$$

where the second bijection is given through

$$
\begin{aligned}
(U \times G) \times_{G} F \longrightarrow U \times F, & {[(x, g), v] \longmapsto(x, g v) } \\
& {[(x, 1), w] \longleftrightarrow(x, w) . }
\end{aligned}
$$

The charts for $E$ then have the same transition functions as the $G$-atlas for $P$. We call $P \times_{G} F \rightarrow M$ a fiber bundle associated to $P$.
(b) (Associated principal bundles) Let $\pi: E \rightarrow M$ be a fiber bundle with typical fiber $F$ and structure group $G$. For $x \in M$ let

$$
P_{x} \equiv \operatorname{Iso}_{G}\left(F ; E_{x}\right):=\left\{\left(\varphi \mid E_{x}\right)^{-1} \circ L_{g}: F \rightarrow E_{x} \mid g \in G\right\}
$$

denote the entity of maps $F \rightarrow E_{x}$ induced with respect to one (and then every) bundle chart $(\varphi, U)$ of $E$ at $x$ by group elements in $G$. Then $P \equiv \operatorname{Iso}_{G}(F ; E):=\bigcup_{x \in M} P_{x}$ is naturally a principal $G$-bundle. Bundle charts for $P$ are obtained by assigning to each chart $(\varphi, U)$ for $E$ the bijection

$$
\begin{array}{r}
P / U \equiv \pi^{-1} U \longleftarrow U \times G \\
\left(\varphi \mid E_{x}\right)^{-1} \circ L_{g} \longleftarrow(x, g) .
\end{array}
$$

We call $P \rightarrow M$ the principal $G$-bundle associated to $E$.
REMARK 1.6.8. The two procedures described in Remark 1.6.7 are functorial constructions inverting each other. More precisely, for a principal $G$-bundle $P \rightarrow M$ and a fiber bundle $E \rightarrow M$ with typical fiber $F$ and structure group $G$, functorial bundle isomorphisms are given as follows:

$$
\begin{aligned}
P \longrightarrow \operatorname{Iso}_{G}\left(F ; P \times_{G} F\right), & u & \longmapsto(v \mapsto[u, v]) \\
E \longleftarrow \operatorname{Iso}_{G}(F ; E) \times_{G} F, & u(v) & \longleftrightarrow[u, v] .
\end{aligned}
$$

Remark 1.6.7 allows to construct from a principal $G$-bundle $P$ and a representation $G \rightarrow \operatorname{Aut}(V)$ of $G$ the vector bundle associated to $P$ with fiber $V$, respectively to pass from a vector bundle $\pi: E \rightarrow M$ of rank $m$ to the principal $\operatorname{GL}(m ; \mathbb{R})$-bundle associated to $E$ (or from a Riemannian vector bundle of rank $m$ to the associated principal $\mathrm{O}(m)$ bundle).

Example 1.6.9. Let $M$ be a differentiable manifold of dimension $n$ and $T M \rightarrow M$ its tangent bundle considered as fiber bundle with typical fiber $\mathbb{R}^{n}$ and structure group $\mathrm{GL}(n ; \mathbb{R})$. The associated principal $\mathrm{GL}(n ; \mathbb{R})$-bundle

$$
\mathrm{L}(T M):=\operatorname{Iso}_{\mathrm{GL}(n ; \mathbb{R})}\left(\mathbb{R}^{n} ; T M\right) \rightarrow M
$$

is called frame bundle over $M$. If $T M$ carries a Riemannian metric $g$ then $T M$ is a vector bundle with structure group $\mathrm{O}(n)$. The associated principal $\mathrm{O}(n)$-bundle

$$
\mathrm{O}(T M):=\operatorname{Iso}_{\mathrm{O}(n)}\left(\mathbb{R}^{n} ; T M\right) \rightarrow M
$$

is called orthonormal frame bundle over $M$.

1. The frame bundle $P=\mathrm{L}(T M)$ over $M$ is the principal $\mathrm{GL}(n ; \mathbb{R})$-bundle associated to the vector bundle $T M$. By construction $u \in P_{x}$ is a linear isomorphism $u: \mathbb{R}^{n} \rightarrow T_{x} M$ which may be identified with the $\mathbb{R}$-basis

$$
\left(u_{1}, \ldots, u_{n}\right):=\left(u e_{1}, \ldots, u e_{n}\right)
$$

for $T_{x} M$ where $e_{i}$ denotes the $i$ th standard coordinate vector of $\mathbb{R}^{n}$. The general linear group $G=\mathrm{GL}(n ; \mathbb{R})$ operates on $\mathrm{L}(T M)$ from the right via

$$
u g: \mathbb{R}^{n} \xrightarrow{g} \mathbb{R}^{n} \xrightarrow{u} T_{x} M .
$$

Thus $u g \in \mathrm{~L}(T M)$ for $g=\left(g_{i j}\right) \in G$ is given by $(u g)_{j}=\sum_{i} g_{i j} u_{i}$. Bundle charts for $\mathrm{L}(T M)$ are obtained from charts $(h, U)$ for $M$. Indeed, $\left(\frac{\partial}{\partial h^{1}}, \ldots, \frac{\partial}{\partial h^{n}}\right)$ defines a local section $\mathrm{L}(T M)$ over $U$ so that each $u \in \mathrm{~L}(T M)$ with $\pi(u)=x \in U$ writes as $u_{j}=\sum_{i} a_{i j}(u)\left(\frac{\partial}{\partial h^{i}}\right)_{x}$ where $a(u):=\left(a_{i j}(u)\right) \in \mathrm{GL}(n ; \mathbb{R})$. Then

$$
\pi^{-1} U \xrightarrow{\sim} U \times \operatorname{GL}(n ; \mathbb{R}), \quad u \longmapsto(\pi(u), a(u))
$$

is a bundle chart for $P=\mathrm{L}(T M)$.
2. The orthonormal frame bundle $\mathrm{O}(T M)$ over a Riemannian manifold $(M, g)$ is the principal $\mathrm{O}(n)$-bundle associated to $T M$ (now $T M$ considered as vector bundle with structure group $\mathrm{O}(n)$ ). To each chart $(h, U)$ for $M$, Gram-Schmidt orthogonalization of $\left(\frac{\partial}{\partial h^{1}}, \ldots, \frac{\partial}{\partial h^{n}}\right)$ gives a local section of $\mathrm{O}(T M)$ over $U$. Bundle charts for $\mathrm{O}(T M)$ are constructed as above.

Bundle charts (local trivializations) allow to identify neighbouring fibers of principal bundles, but not in a canonical way. As for vector bundles, to relate fibers intrinsically to each others, a connection as additional structure is required. We will see that for principal $\operatorname{GL}(n ; \mathbb{R})$-bundles connections correspond canonically to linear connections in associated vector bundles.

DEFINITION 1.6.10 (Connection in a principal bundle). Let $\pi: P \rightarrow M$ be a principal bundle over $M$ with structure group $G$. A $G$-connection in $P$ is a differentiable $G$-invariant splitting $h$ of the following exact sequence of vector bundles over $P$ :

where $d \pi \circ h=\mathrm{id}$. This splitting induces a decomposition of $T P$ :

$$
T P=V \oplus H:=\operatorname{ker} d \pi \oplus h\left(\pi^{*} T M\right)
$$

The $G$-invariance of the splitting means that $H_{u g}=\left(d R_{g}\right) H_{u}$ for each $u \in P$, where $R_{g} u:=u g$ denotes the right action of $g \in G$. For $u \in P$, we call $H_{u}$ the horizontal space at $u$ and $V_{u}=\left\{v \in T_{u} P:(d \pi) v=0\right\}$ the vertical space at $u$. The bundle isomorphism

$$
\begin{equation*}
h: \pi^{*} T M \xrightarrow{\sim} H \longleftrightarrow T P \tag{1.6.1}
\end{equation*}
$$

is called horizontal lift of the $G$-connection; fiberwise it reads as $h_{u}: T_{\pi(u)} M \xrightarrow{\sim} H_{u}$.
REMARK 1.6.11. For each $u \in P$ the vertical space $V_{u}$ is canonically given and $G$ invariant. However there is no canonical choice of a complement $H_{u}$ : A $G$-connection in $P$ corresponds exactly to a $G$-invariant choice of a horizontal space $H_{u}$ for each $u \in P$. By means of the $G$-connection in $P$ each vector field $X \in \Gamma(T P)$ decomposes in a horizontal and a vertical part:

$$
X=\operatorname{hor} X+\operatorname{vert} X
$$

DEFINITION 1.6.12 (Standard-vertical vector field). Let $\pi: P \rightarrow M$ be a principal $G$-bundle over $M$. Each $u \in P$ defines an embedding

$$
I_{u}: G \longleftrightarrow P, \quad g \mapsto u g
$$

Its differential at the unit element $e \in G$,

$$
\begin{equation*}
\iota_{u} \equiv\left(d I_{u}\right)_{e}: T_{e} G \rightarrow T_{u} P, \quad A \longmapsto \hat{A}(u) \tag{1.6.2}
\end{equation*}
$$

gives an identification $\kappa_{u}: \mathfrak{g} \xrightarrow{\sim} V_{u}$ of the Lie algebra $\mathfrak{g}=T_{e} G$ of $G$ with the vertical fiber $V_{u}$ at $u$. The vertical vector field $\hat{A} \in \Gamma(T P)$ on $P$ defined by (1.6.2) is called standard-vertical vector field to $A \in \mathfrak{g}$.

DEFINITION 1.6.13 (Connection form). Let $\pi: P \rightarrow M$ be a principal $G$-bundle over $M$ equipped with a $G$-connection. The $\mathfrak{g}$-valued one-form $\omega \in \Gamma\left(T^{*} P \otimes \mathfrak{g}\right)$ on $P$,

$$
\begin{equation*}
\omega_{u}\left(X_{u}\right):=\kappa_{u}^{-1}(\operatorname{vert} X)_{u}, \quad X \in \Gamma(T P) \tag{1.6.3}
\end{equation*}
$$

is called connection form of the $G$-connection.
By definition, the connection form $\omega$ of a $G$-connection is horizontal, i.e., $\omega(X)=0$ for a vector field $X$ on $P$ if and only if $X$ is horizontal.

REMARK 1.6.14. Let $\pi: P \rightarrow M$ be a principal $G$-bundle over $M$ and let $\omega \in$ $\Gamma\left(T^{*} P \otimes \mathfrak{g}\right)$ be the connection form of a $G$-connection in $P$. Then:
(i) $\omega(\hat{A})=A$ for $A \in \mathfrak{g}$;
(ii) $\omega$ is equivariant, i.e.,

$$
R_{g}^{*} \omega=\operatorname{Ad}\left(g^{-1}\right) \omega, \quad g \in G
$$

where $\operatorname{Ad}$ is the adjoint representation of $G$ in $\mathfrak{g}$. Recall that $\operatorname{Ad}(g): \mathfrak{g} \rightarrow \mathfrak{g}$ for $g \in G$ is defined as the differential of the automorphism $I(g): G \rightarrow G, h \longmapsto R_{g^{-1}} L_{g} h:=g h g^{-1}$ at $e \in G$. Identifying $\mathfrak{g}$ and left-invariant vector fields on $G$, and taking into account that $\left(d L_{g}\right) A \equiv A$, i.e., $\left(d L_{g}\right)_{h} A_{h} \equiv A_{g h}$, we get $\operatorname{Ad}(g)=d R_{g^{-1}}$.

Proof. Item (i) is a direct consequence of Definition 1.6.12 and 1.6.13. We want to verify item (ii) which reads as

$$
\begin{equation*}
\omega\left(\left(d R_{g}\right) X\right)=\operatorname{Ad}\left(g^{-1}\right) \omega(X), \quad X \in T_{u} P, g \in G \tag{1.6.4}
\end{equation*}
$$

It is obviously sufficient to consider the following two cases:
(1) $X$ is horizontal, i.e. $X \in H_{u}$. Since $\left(d R_{g}\right) H_{u}=H_{u g}$, however also $\left(d R_{g}\right) X$ is horizontal, so that both sides of (1.6.4) vanish.
(2) $X$ is vertical, i.e. $X \in V_{u}$, and hence $X=\hat{A}(u)$ for some $A \in \mathfrak{g}$. But then $\left(d R_{g}\right)_{u} \hat{A}(u)=\left(\operatorname{Ad}\left(g^{-1}\right) A\right)^{\wedge}(u g)$ with $\left(\operatorname{Ad}\left(g^{-1}\right) A\right)^{\wedge} \in \Gamma(T P)$ the standard-vertical vector field to $\operatorname{Ad}\left(g^{-1}\right) A \in \mathfrak{g}$, and one obtains

$$
\begin{aligned}
\left(\left(R_{g}\right)^{*} \omega\right)_{u}(X) & =\omega_{u g}\left(\left(d R_{g}\right)_{u} X\right)=\omega_{u g}\left(\left(d R_{g}\right)_{u} \hat{A}(u)\right) \\
& =\omega_{u g}\left(\left(\operatorname{Ad}\left(g^{-1}\right) A\right)^{\wedge}(u \cdot g)\right)=\operatorname{Ad}\left(g^{-1}\right) A \\
& =\operatorname{Ad}\left(g^{-1}\right) \omega_{u}(\hat{A}(u))=\operatorname{Ad}\left(g^{-1}\right) \omega_{u}(X)
\end{aligned}
$$

which shows the claim.
REMARKS 1.6.15. (i) A $G$-connection in $P$ is uniquely determined by its connection form $\omega$ : the map $\omega_{u}: T_{u} P \rightarrow \mathfrak{g}$ is linear for each $u \in P$ with $\operatorname{ker} \omega_{u}=H_{u}$. Conversely, every equivariant differential form

$$
\omega \in \Gamma\left(T^{*} P \otimes \mathfrak{g}\right) \quad \text { with } \omega(\hat{A})=A \text { for } A \in \mathfrak{g}
$$

defines a $G$-connection in $P$ whose connection form is given by $\omega$.
(ii) $G$-invariant splittings of the exact sequence (1.6.10), and hence connections in principal bundles, can thus be described in different ways, for instance,

- as horizontal lifts $h, d \pi \circ h=\mathrm{id}$, with the property that $\left(d R_{g}\right) h=h$ for each $g \in G$;
- as $\mathfrak{g}$-valued differential forms $\omega \in \Gamma\left(T^{*} P \otimes \mathfrak{g}\right), \bar{\omega} \circ \iota=$ id, where $\bar{\omega}_{u}=\left(u, \omega_{u}\right)$ for $u \in P$, satisfying the condition $R_{g}^{*} \omega=\operatorname{Ad}\left(g^{-1}\right) \omega$ for $g \in G$;
- as horizontal subbundles $H, V \oplus H=T P$, such that $\left(d R_{g}\right) H_{u}=H_{u g}$ for each $u \in P$.


Fiberwise diagram (1.6.5) reads as


As already mentioned, there is a one-to-one correspondence between connections in principal bundles and linear connections in associated vector bundles. More precisely we have the following situation.

REMARK 1.6.16. (a) Let $\pi: P \rightarrow M$ be a principal $G$-bundle, $V$ a real vector space, $G \rightarrow \operatorname{Aut}(V)$ a representation of $G$ and $E=P \times_{G} V$ the associated vector bundle with fiber $V$. Each $G$-connection in $P$ induces a linear connection in $E$ as follows. Denote by

$$
P \times V \longrightarrow E=(P \times V) / G, \quad(u, \xi) \longmapsto(u, \xi) G
$$

the canonical projection and consider elements $\xi \in V$ as

$$
\xi: P \rightarrow E, \quad u \longmapsto u \xi:=(u, \xi) G .
$$

By definition we have $(u g) \xi=u(g \xi)$. We fix $e \in E$ and choose an arbitrary $u \in P_{\pi_{E}(e)}$. Then there exists exactly one $\xi \in V$ such that $u \xi=e$, and one obtains:


The differential $(d \xi)_{u}$ of $\xi: P \rightarrow E$ at $u$ gives a map $T_{u} P \rightarrow T_{e} E$ such that


By assumption we have a decomposition $T_{u} P=V_{u}(P) \oplus H_{u}(P)$ which induces a decomposition $T_{e} E=V_{e}(E) \oplus H_{e}(E)$ with $V_{e}(E)=\operatorname{ker}\left(d \pi_{E}\right)_{e}$ and $H_{e}(E)$ the image of $H_{u}(P)$ under $(d \xi)_{u}$. One verifies that $H_{e}(E)$ is well-defined, i.e., independent of the choice of $u$. Indeed taking $u g$ instead of $u$ leads to $(u \cdot g)\left(g^{-1} \xi\right)=e$, but by assumption $\left(d R_{g}\right) H_{u}(P)=H_{u g}(P)$ and consequently

$$
d\left(g^{-1} \xi\right)_{u g} H_{u g}(P)=d\left(\left(g^{-1} \xi\right) \circ R_{g}\right)_{u} H_{u}(P)=(d \xi)_{u} H_{u}(P)
$$

Hence $H(E)$ defines a subbundle of $T E$ which determines a linear connection in $E$.
(b) Conversely let $E$ be a vector bundle over $M$ with structure group $G$ and $\operatorname{Iso}_{G}(V ; E)$ the associated principal $G$-bundle over $M$; without restrictions we may assume $V=\mathbb{R}^{m}$. Then each linear connection in $E$ induces a $G$-connection in $P$ as follows. According to Remark 1.4.18 (i) the horizontal space at $e \in E$ (with $\pi_{E}(e)=x$ ) induced by the connection in $E$ writes as

$$
H_{e}(E)=\left\{(d X)_{x} v: v \in T_{x} M, X \in \Gamma(E) \text { with } X(x)=e \text { and } \nabla_{v} X=0\right\}
$$

Note that each section $\hat{X} \in \Gamma(P)$ is of the form $\hat{X}=\left(X_{1}, \ldots, X_{n}\right)$ with $X_{i}:=\hat{X} e_{i} \in$ $\Gamma(E)$ and $e_{i}$ the $i$-th standard coordinate vector of $\mathbb{R}^{n}$. For $u \in P$ with $\pi(u)=x$ let now

$$
\begin{aligned}
H_{u}(P):=\left\{(d \hat{X})_{x} v: v\right. & \in T_{x} M, \hat{X} \in \Gamma(P / U) \text { with } \hat{X}(x)=u \\
& \text { and } \left.\nabla_{v} \hat{X}:=\left(\nabla_{v} X_{1}, \ldots, \nabla_{v} X_{n}\right)=0\right\}
\end{aligned}
$$

This determines horizontal subbundle $H(P)$ of $T P$ which satisfies $\left(d R_{g}\right) H_{u}=H_{u g}$ for $u \in P$. Hence it defines a $G$-connection in $P$.

One verifies that (a) and (b) are inverse to each other when passing from frame bundles to associated vector bundles, resp., from vector bundles to the associated frame bundles (see Remark 1.6.7). In particular, we have one-to-one correspondences:

$$
\begin{aligned}
\text { linear connections in } T M \longleftrightarrow \mathrm{GL}(n ; \mathbb{R}) \text {-connections in } \mathrm{L}(T M) \text {; } \\
\text { Riemannian connections in } T M \longleftrightarrow \mathrm{O}(n) \text {-connections in } \mathrm{O}(T M)
\end{aligned}
$$

In the sequel we call $\operatorname{GL}(n ; \mathbb{R})$-connections in the frame bundle $\mathrm{L}(T M)$ briefly linear connections on $M$, and $\mathrm{O}(n)$-connections in the orthonormal frame bundle $\mathrm{O}(T M)$ Riemannian connections on $M$.

THEOREM 1.6.17 (Horizontal lifts in principal bundles). Let $\pi: P \rightarrow M$ be a principal $G$-bundle over $M$ equipped with a $G$-connection. Furthermore, let $x: I \rightarrow M$, $t \mapsto x(t)$, be a differentiable curve and $t_{0} \in I$. Then, to each $u_{0} \in P$ with $\pi\left(u_{0}\right)=x\left(t_{0}\right)$,
there exists exactly one horizontal curve $u: I \rightarrow P$ with $u\left(t_{0}\right)=u_{0}$ which is above $t \mapsto x(t)$, i.e., such that $(\pi \circ u)(t)=x(t)$ and $\dot{u}(t) \in H_{u(t)}$ for each $t \in I$.

Proof. It is obviously sufficient to verify existence and uniqueness of the horizontal lift locally about $t_{0}$. By means of the $G$-invariance, along with $t \mapsto u(t)$ and $g \in G$ also $t \mapsto u(t) g$ is a horizontal curve above $t \mapsto x(t)$. Hence pieces of local horizontal lifts can be patched together to obtain the horizontal lift defined on all of $I$.

Let $\Phi: P \times G \rightarrow P,(u, g) \mapsto u g$, denote the right action of $G$ on $P$. For fixed $g \in G$ then $R_{g} \equiv \Phi(\cdot, g): P \rightarrow P$ is right multiplication by $g$, and for $u \in P$ we have the already considered embedding $I_{u} \equiv \Phi(u, \cdot): G \rightarrow P$. Recall that we used its differential $\iota_{u} \equiv\left(d I_{u}\right)_{e}: \mathfrak{g} \rightarrow T_{u} P$ at $e \in G$ to identify $\mathfrak{g}$ and the vertical fiber $V_{u}$, in particular $\omega \circ \iota_{u}=\operatorname{id}_{\mathfrak{g}}$. Let now $t \mapsto x(t)=: x_{t}$ be the curve in $M$ where without restrictions we may assume that $x(\cdot)$ takes values in the domain $U$ of a bundle chart $(\varphi, U)$. By means of $\varphi: \pi^{-1}(U) \xrightarrow{\sim} U \times G$ we first procure some differentiable curve $t \mapsto v_{t}$ in $P$ which lies above $t \mapsto x_{t}$, for instance $v_{t}:=\varphi^{-1}\left(x_{t}, e\right)$ with $e$ the unit element in $G$. Next we search a curve $t \mapsto \tilde{g}_{t}$ in $G$ such that $t \mapsto u_{t}:=v_{t} \tilde{g}_{t}$ becomes horizontal, i.e. $\left.\omega(\dot{u}) \equiv 0\right)$. Writing $v_{t}=u_{t} g_{t}$ with $g_{t}:=\tilde{g}_{t}^{-1}$, we get $\dot{v}_{t}=\left(R_{g_{t}}\right)_{*} \dot{u}_{t}+\left(I_{u_{t}}\right)_{*} \dot{g}_{t}$. By means of the equivariance of $\omega$ (Remark 1.6.14) along with the relation $I_{u_{t}}=I_{u_{t} g_{t}} \circ L_{g_{t}^{-1}}$, we then have

$$
\begin{equation*}
\omega\left(\dot{v}_{t}\right)=\operatorname{Ad}\left(g_{t}^{-1}\right) \omega\left(\dot{u}_{t}\right)+\left(L_{g_{t}^{-1}}\right)_{*} \dot{g}_{t} \tag{1.6.7}
\end{equation*}
$$

To make $t \mapsto u_{t}$ horizontal the condition $\dot{g}_{t}=\left(L_{g_{t}}\right)_{*} \omega\left(\dot{v}_{t}\right)$ is required. This is a differential equation for $t \mapsto g_{t}$ which we write as

$$
\begin{equation*}
\dot{g}_{t}=A\left(t, g_{t}\right) \tag{1.6.8}
\end{equation*}
$$

where $A(t, \cdot) \in \Gamma(T G)$ defined as $A(t, h)=\left(L_{h}\right)_{*} \omega\left(\dot{v}_{t}\right)$ is just the left-invariant vector field on $G$ associated to $\omega\left(\dot{v}_{t}\right) \in \mathfrak{g}$. The proof is completed by the fact that Eq. (1.6.8) has a unique local solution to each given initial condition.

REMARK 1.6.18. Note that uniqueness in Theorem 1.6.17 comes from the unique solvability of differential equation (1.6.8) for a given initial condition. Uniqueness can also be seen directly: If $t \mapsto u_{t}$ and $t \mapsto v_{t}$ are two horizontal lifts, we may write $v_{t}=u_{t} g_{t}$. However, since $\omega(\dot{v})=\omega(\dot{u}) \equiv 0$ we conclude from (1.6.7) that $g_{t} \equiv$ constant. Thus if $u_{t_{0}}=v_{t_{0}}$ for one $t_{0}$, then necessarily $g_{t} \equiv e$.

REMARK 1.6.19. If in the situation of Theorem 1.6 .17 there is a representation $G \rightarrow$ $\operatorname{Aut}(V)$ of $G$ and $E=P \times_{G} V$ the vector bundle associated to $P$, then $P \cong \operatorname{Iso}_{G}(V ; E)$, and Theorem 1.6.17 be reduced to the already established existence and uniqueness of horizontal lifts in vector bundles. Indeed, by Remark 1.6.16 (a) one obtains a linear connection in $E$, and by Theorem 1.4 .19 to $\xi \in V$ there exists exactly one horizontal curve $I \rightarrow E$, $t \mapsto e_{\xi}(t)$ over $t \mapsto x(t)$ such that $e_{\xi}\left(t_{0}\right)=u_{0} \xi \equiv\left(u_{0}, \xi\right) \cdot G$. Then the curve

$$
t \mapsto u(t) \in \operatorname{Iso}_{G}\left(V ; E_{x(t)}\right), \quad u(t):=\left(\xi \mapsto e_{\xi}(t)\right)
$$

in $P$ is horizontal and has the wanted properties.
Corollary 1.6.20. Let $P$ be a principal $G$-bundle over a manifold $M$.
(1) Every $G$-connection in $P$ defines canonically a parallel transport in $P$ along differentiable curves $t \mapsto x(t)$ in $M$, namely for $t_{0}, t_{1} \in I$ as

$$
/ / t_{0}, t_{1}: P_{x\left(t_{0}\right)} \xrightarrow{\sim} P_{x\left(t_{1}\right)}, \quad u_{0} \longmapsto u\left(t_{1}\right)
$$

where $t \mapsto u(t)$ is the according to Theorem 1.6.17 uniquely determined horizontal lift of $t \mapsto x(t)$ to $P$ such that $u\left(t_{0}\right)=u_{0}$.
(2) If $E$ a vector bundle associated to $P$ with fiber $V$, then this parallel transport induces a parallel transport in $E$, namely as

$$
/ / t_{0}, t_{1}: E_{x\left(t_{0}\right)} \xrightarrow{\sim} E_{x\left(t_{1}\right)}, \quad e_{0} \longmapsto u\left(t_{1}\right) \xi
$$

where as above $t \mapsto u(t)$ with $u\left(t_{0}\right)=u_{0}$ is horizontal lift of $t \mapsto x(t)$ to $P$, and $\xi \in V$ is chosen such that $u\left(t_{0}\right) \xi \equiv u_{0} \xi=e_{0}$.

For the remainder of this section we restrict ourselves to the principal bundle $P=$ $\mathrm{L}(T M)$ over a differentiable manifold $M$ with structure group $G=\mathrm{GL}(n ; \mathbb{R})$, respectively, $P=\mathrm{O}(T M)$ over a Riemannian manifold $M$ with $G=\mathrm{O}(n)$. The corresponding Lie algebras are then the matrix algebras

$$
\mathfrak{g}=\mathbf{M}(n \times n ; \mathbb{R}), \quad \text { resp. }, \mathfrak{g}=\{A \in \mathbf{M}(n \times n ; \mathbb{R}): A \text { skew symmetric }\}
$$

Fixing a $G$-connection in $P$, we have the $\mathfrak{g}$-valued connection form (see Definition 1.6.13)

$$
\begin{equation*}
\omega \in \Gamma\left(T^{*} P \otimes \mathfrak{g}\right), \quad \omega_{u}\left(X_{u}\right)=\kappa_{u}^{-1}(\operatorname{vert} X)_{u}, u \in P \text { and } X \in \Gamma(T P) \tag{1.6.9}
\end{equation*}
$$

In addition to the connection form $\omega$ we have the canonical one-form of the principal bundle $\pi: P \rightarrow M$,

$$
\begin{equation*}
\vartheta \in \Gamma\left(T^{*} P \otimes \mathbb{R}^{n}\right), \quad \vartheta_{u}\left(X_{u}\right):=u^{-1}\left(d \pi X_{u}\right), u \in P \text { and } X \in \Gamma(T P) \tag{1.6.10}
\end{equation*}
$$

where as usual $u \in P$ is read as linear isomorphism, resp. isometry, $u: \mathbb{R}^{n} \xrightarrow{\sim} T_{\pi(u)} M$. Note that contrary to the connection form the canonical one-form $\vartheta$ does not depend on the chosen $G$-connection.

THEOREM 1.6.21. The frame bundles $P=\mathrm{L}(T M)$ ( $M$ manifold), resp. $P=\mathrm{O}(T M)$ ( $M$ Riemannian manifold), considered as manifolds, are parallelizable, i.e., the tangent bundles $T \mathrm{~L}(T M) \rightarrow \mathrm{L}(T M)$ and $T \mathrm{O}(T M) \rightarrow \mathrm{O}(T M)$ are trivial.

Proof. Indeed a $G$-connection in $P$ decomposes $T P=V \oplus H$. A canonical trivialization for $T P$ is given as follows: the vertical subbundle $V$ is trivialized by the standardvertical vector fields $\hat{A}$ to $A$, where $A$ runs through a basis of $\mathfrak{g}$; the horizontal subbundle $H$ is trivialized by the standard-horizontal vector fields $L_{1}, \ldots, L_{n}$ in $\Gamma(T P)$ defined by

$$
L_{i}(u):=h_{u}\left(u e_{i}\right)
$$

For every $u \in P$,

$$
\left(\hat{A}(u), L_{i}(u): A \in \text { basis for } \mathfrak{g}, i=1, \ldots, n\right)
$$

is a basis of $T_{u} P=V_{u} \oplus H_{u}$. This is obvious from the isomorphisms $\mathfrak{g} \xrightarrow{\sim} V_{u}, A \mapsto \hat{A}(u)$ and $h_{u}: T_{\pi(u)} M \xrightarrow{\sim} H_{u}$.

REMARK 1.6.22. The standard-vertical, respectively standard-horizontal vector fields are determined by the relations

$$
\begin{aligned}
& \vartheta(\hat{A})=0 \text { and } \vartheta\left(L_{i}\right)=e_{i} \\
& \omega(\hat{A})=A \text { and } \omega\left(L_{i}\right)=0
\end{aligned}
$$

The canonical second order partial differential operator

$$
\Delta^{\mathrm{hor}}:=\sum_{i=1}^{n} L_{i}^{2}
$$

is called horizontal Laplacian on $\mathrm{L}(T M)$, resp. $\mathrm{O}(T M)$.

NOTATION 1.6.23. Let $\pi: P \rightarrow M$ be a principal $G$-bundle over $M$ equipped with a $G$-connection. For a vector field $X \in \Gamma(T M)$ we denote by

$$
\bar{X} \in \Gamma(T P), \quad \bar{X}_{u}=h_{u}\left(X_{\pi(u)}\right), u \in P
$$

the corresponding horizontal lift to $P$.
Lemma 1.6.24. Let $M$ be a differentiable manifold with a linear connection in $T M$. If $X, Y \in \Gamma(T M)$ and $\alpha \in \Gamma\left(T^{*} M\right)$, then

$$
\begin{aligned}
& \left(\nabla_{X} Y\right)_{x}=\lim _{\varepsilon \downarrow 0} \frac{/_{0, \varepsilon}^{-1} Y_{\gamma(\varepsilon)}-Y_{\gamma(0)}}{\varepsilon} \\
& \left(\nabla_{X} \alpha\right)_{x}=\lim _{\varepsilon \downarrow 0} \frac{/ /_{0, \varepsilon}^{-1} \alpha_{\gamma(\varepsilon)}-\alpha_{\gamma(0)}}{\varepsilon}
\end{aligned}
$$

where $/ /_{0, \varepsilon}: T_{\gamma(0)} M \rightarrow T_{\gamma(\varepsilon)} M$ is the parallel transport along a curve $\gamma$ on $M$ with the properties that $\gamma(0)=x$ and $\dot{\gamma}(0)=X_{x}$.
(Note that for $Y \in \Gamma(T M)$ by definition $\left(/ / /_{0, \varepsilon}^{-1} \alpha_{\gamma(\varepsilon)}\right)\left(Y_{\gamma(0)}\right)=\alpha_{\gamma(\varepsilon)}\left(/ / 0_{0, \varepsilon} Y_{\gamma(0)}\right)$.
Proof. The first formula has been shown in Remark 1.4.14, to verify the second one first note that

$$
\begin{aligned}
\frac{/ /_{0, \varepsilon}^{-1} \alpha_{\gamma(\varepsilon)}-\alpha_{\gamma(0)}}{\varepsilon}\left(Y_{\gamma(0)}\right) & =\frac{\alpha_{\gamma(\varepsilon)}\left(/ /_{0, \varepsilon} Y_{\gamma(0)}\right)-\alpha_{\gamma(0)} Y_{\gamma(0)}}{\varepsilon} \\
& =\frac{\alpha_{\gamma(\varepsilon)}\left(/ /_{0, \varepsilon} Y_{\gamma(0)}-Y_{\gamma(\varepsilon)}\right)}{\varepsilon}+\frac{\alpha_{\gamma(\varepsilon)} Y_{\gamma(\varepsilon)}-\alpha_{\gamma(0)} Y_{\gamma(0)}}{\varepsilon}
\end{aligned}
$$

Taking the limit as $\varepsilon \downarrow 0$, the right-hand side converges to

$$
\begin{aligned}
& \lim _{\varepsilon \downarrow 0} \alpha_{\gamma(0)} \frac{/ /_{0, \varepsilon}^{-1}\left(/ /_{0, \varepsilon} Y_{\gamma(0)}-Y_{\gamma(\varepsilon)}\right)}{\varepsilon}+\lim _{\varepsilon \downarrow 0} \frac{\alpha_{\gamma(\varepsilon)} Y_{\gamma(\varepsilon)}-\alpha_{\gamma(0)} Y_{\gamma(0)}}{\varepsilon} \\
& \quad=\alpha_{\gamma(0)}\left(\lim _{\varepsilon \downarrow 0} \frac{Y_{\gamma(0)}-/ /_{\varepsilon, 0} Y_{\gamma(\varepsilon)}}{\varepsilon}\right)+\lim _{\varepsilon \downarrow 0} \frac{\alpha_{\gamma(\varepsilon)} Y_{\gamma(\varepsilon)}-\alpha_{\gamma(0)} Y_{\gamma(0)}}{\varepsilon} \\
& \quad=-\alpha_{x}\left(\nabla_{X_{x}} Y\right)+X_{x}(\alpha Y)=\left(\nabla_{X} \alpha\right)_{x}\left(Y_{x}\right)
\end{aligned}
$$

which gives the claim.
NOTATION 1.6.25. Let $M$ be a differentiable manifold and $P=\mathrm{L}(T M)$ be the frame bundle over $M$, respectively, $M$ a Riemannian manifold and $P=\mathrm{O}(T M)$ the orthonormal frame bundle over $M$. It is convenient to write vector fields $Y \in \Gamma(T M)$ and differential forms $\alpha \in \Gamma\left(T^{*} M\right)$ as equivariant functions on the frame bundle $P$,

$$
\begin{array}{ll}
f_{Y}: P \rightarrow \mathbb{R}^{n}, & f_{Y}(u):=u^{-1} Y_{\pi(u)} \\
F_{\alpha}: P \rightarrow \mathbb{R}^{n}, & F_{\alpha}^{i}(u):=\alpha_{\pi(u)}\left(u e_{i}\right), \quad i=1, \ldots, n
\end{array}
$$

Equivariance means that for $g \in G=\mathrm{GL}(n ; \mathbb{R})$, respectively $\mathrm{O}(n)$,

$$
\begin{aligned}
f_{Y}(u g) & =g^{-1} f_{Y}(u) \\
F_{\alpha}(u g) & =g^{*} F_{\alpha}(u)
\end{aligned}
$$

where $g^{-1}$ and $g^{*}$ are the inverse, resp. dual linear map to $g$.

THEOREM 1.6.26. Let $M$ be a differentiable manifold and $P=\mathrm{L}(T M)$ be the frame bundle over $M$ endowed with a $\mathrm{GL}(n ; \mathbb{R})$-connection, respectively, $M$ a Riemannian manifold and $P=\mathrm{O}(T M)$ the associated orthonormal frame bundle over $M$ endowed with a $\mathrm{O}(n)$-connection. Then, for vector fields $X, Y \in \Gamma(T M)$, the covariant derivative $\nabla_{X} Y \in \Gamma(T M)$ with respect to the induced linear connection in $T M$ is given by

$$
\begin{equation*}
\left(\nabla_{X} Y\right)_{x}=u\left(\bar{X}_{u} \vartheta(\bar{Y})\right) \quad \text { for } u \in P \text { with } \pi(u)=x \tag{1.6.11}
\end{equation*}
$$

Writing $Y \in \Gamma(T M)$ as equivariant function $f_{Y}$ on $P$, this formula reads as

$$
\left(\nabla_{X} Y\right)_{x}=u\left(\bar{X}_{u} f_{Y}\right)
$$

or equivalently:

$$
\begin{equation*}
\left\langle u^{-1}\left(\nabla_{X} Y\right)_{x}, e_{i}\right\rangle=\bar{X}_{u}\left(f_{Y}^{i}\right), \quad i=1, \ldots, n \tag{1.6.12}
\end{equation*}
$$

Proof. We choose a curve $\gamma$ on $M$ such that $\gamma(0)=x$ and $\dot{\gamma}(0)=X_{x}$. Let $t \mapsto u(t)$ be a horizontal lift of $t \mapsto \gamma(t)$ to $P$. Then, by Corollary 1.6.20, the parallel transport $/ /_{0, \varepsilon}: T_{\gamma(0)} M \rightarrow T_{\gamma(\varepsilon)} M$ along $\gamma$ is given by $/ /_{0, \varepsilon}=u(\varepsilon) u(0)^{-1}$. By Lemma 1.6.24 we get

$$
\begin{aligned}
\left(\nabla_{X} Y\right)_{\pi(u)} & =\lim _{\varepsilon \downarrow 0} \frac{u(0) u(\varepsilon)^{-1} Y_{\pi(u(\varepsilon))}-Y_{\pi(u(0))}}{\varepsilon} \\
& =\lim _{\varepsilon \downarrow 0} \frac{u\left(u(\varepsilon)^{-1} Y_{\pi(u(\varepsilon))}-u(0)^{-1} Y_{\pi(u(0))}\right)}{\varepsilon} \\
& =u\left(\lim _{\varepsilon \downarrow 0} \frac{f(u(\varepsilon))-f(u(0))}{\varepsilon}\right)=u\left(\bar{X}_{u}(f)\right) .
\end{aligned}
$$

In the last equality we used $\dot{u}(t)=h_{u(t)} \dot{\gamma}(t)$ which implies $\dot{u}(0)=h_{u}\left(X_{x}\right)=\bar{X}_{u}$.
We can give formulas analogous to Eq. (1.6.12) also for the covariant derivative of differential forms can be described. We note the result for later reference.

THEOREM 1.6.27. Let $M$ be a differentiable manifold and $P=\mathrm{L}(T M)$ be the frame bundle over $M$ endowed with a $\mathrm{GL}(n ; \mathbb{R})$-connection, respectively, $M$ a Riemannian manifold and $P=\mathrm{O}(T M)$ the associated orthonormal frame bundle over $M$ endowed with a $\mathrm{O}(n)$-connection. Furthermore, let $\alpha \in \Gamma\left(T^{*} M\right)$ be a differential form on $M$, according to Notation 1.6.25, read as equivariant function on $P$. Then, for $X \in \Gamma(T M)$ with horizontal lift $\bar{X} \in \Gamma(T P)$ and $u \in P$, the following formula holds:

$$
\begin{equation*}
\left(\nabla_{X} \alpha\right)_{\pi(u)}\left(u e_{i}\right)=\bar{X}_{u} F^{i}, \quad i=1, \ldots, d \tag{1.6.13}
\end{equation*}
$$

Proof. We may proceed as in the proof of Theorem 1.6.26. Let $u \in P$ and $t \mapsto \gamma(t)$ be a curve on $M$ such that $\gamma(0)=\pi(u)$ and $\dot{\gamma}(0)=X_{\pi(u)}$. Furthermore, let $t \mapsto u(t) \in P$ be the horizontal lift of $\gamma$ with $u(0)=u$. By Lemma 1.6.24 we obtain:

$$
\begin{aligned}
&\left(\nabla_{X} \alpha\right)_{\pi(u)}\left(u e_{i}\right)=\lim _{\varepsilon \downarrow 0} \frac{\left(/ /_{0, \varepsilon}^{-1} \alpha_{\gamma(\varepsilon)}\right)\left(u e_{i}\right)-\alpha_{\gamma(0)}\left(u e_{i}\right)}{\varepsilon} \\
&\left.=\lim _{\varepsilon \downarrow 0} \frac{\alpha_{\gamma(\varepsilon)}(/ / 0, \varepsilon}{}\left(u e_{i}\right)\right)-\alpha_{\gamma(0)}\left(u e_{i}\right) \\
& \varepsilon \\
&=\lim _{\varepsilon \downarrow 0} \frac{\alpha_{\pi \circ u(\varepsilon)}\left(u(\varepsilon) e_{i}\right)-\alpha_{\pi \circ u(0)}\left(u(0) e_{i}\right)}{\varepsilon} \\
&=\lim _{\varepsilon \downarrow 0} \frac{F^{i}(u(\varepsilon))-F^{i}(u(0))}{\varepsilon}=\bar{X}_{u} F^{i} .
\end{aligned}
$$

DEFINITION 1.6.28 (Horizontal lift of an $M$-valued semimartingale). Let $P$ a principal $G$-bundle over a differentiable manifold $M$ and $\omega \in \Gamma\left(T^{*} P \otimes \mathfrak{g}\right)$ be the connection form of a $G$-connection in $P$. For a $P$-valued semimartingale $U$, the process $\int_{U} \omega$ takes values by definition in the Lie algebra $\mathfrak{g}$ and is defined with respect to a basis of $\mathfrak{g}$ as

$$
\int_{U} \omega \equiv\left(\int_{U} \omega^{1}, \ldots, \int_{U} \omega^{r}\right), \quad \omega=\left(\omega^{1}, \ldots, \omega^{r}\right) .
$$

(a) The process $U$ is called horizontal if $\int_{U} \omega=0$ a.s.
(b) For an $M$-valued semimartingale $X$, a semimartingale $U$ taking values in $P$ is called horizontal lift of $X$ if $\pi \circ U=X$ a.s. and if in addition $U$ is horizontal.
Definition 1.6 .28 generalizes the classical notion of horizontal lift for differentiable curves in $M$ (see Theorem 1.6.17): a curve $t \mapsto u(t)$ over $t \mapsto x(t)$ is called horizontal if $\pi \circ u=x$ and $\omega(\dot{u})=0$ (see Example 1.3.11). Existence of horizontal lifts for semimartingales will be proved in Theorem 1.6.35 below.

For the remainder of this Section we deal with the following situation: Let $M$ be an $n$-dimensional manifold equipped with a torsion-free linear connection, respectively, a Riemannian manifold with the Levi-Civita connection. Above $M$ we consider the frame bundle $P=\mathrm{L}(T M)$ with structure group $\mathrm{GL}(n ; \mathbb{R})$, respectively, the orthonormal frame bundle $P=\mathrm{O}(T M)$ with structure group $\mathrm{O}(n)$, each with the induced $G$-connection on $P$. In addition to the connection form $\omega \in \Gamma\left(T^{*} P \otimes \mathfrak{g}\right)$ we have the canonical one-form $\vartheta \in \Gamma\left(T^{*} P \otimes \mathbb{R}^{n}\right)$, see (1.6.10). The induced decomposition $T P=V \oplus H$ is then given by $V_{u}=\operatorname{ker} \vartheta_{u}$ and $H_{u}=\operatorname{ker} \omega_{u}$ for $u \in P$.

DEFINITION 1.6.29 (Anti-development of an $M$-valued semimartingale). Let $X$ be an $M$-valued semimartingale and $U$ a horizontal lift of $X$ taking values in $P=\mathrm{L}(T M)$, resp. $\mathrm{O}(T M)$. The $\mathbb{R}^{n}$-valued semimartingale

$$
Z=\int_{U} \vartheta \equiv \int \vartheta(\circ d U)
$$

is called anti-development of $X$ into $\mathbb{R}^{n}$ (with respect to the initial frame $U_{0}$ ). In terms of the standard basis of $\mathbb{R}^{n}$ we have $Z \equiv\left(Z^{1}, \ldots, Z^{n}\right)$ where $Z^{i}=\int_{U} \vartheta^{i}$. We call

$$
\mathscr{A}(X)=U_{0} \int_{U} \vartheta \equiv \int \vartheta(\circ d U)
$$

anti-development of $X$ into $T_{X_{0}} M$, or briefly anti-development of $X$. Note that $\mathscr{A}(X)$ is independent of the choice of $U_{0}$.

THEOREM 1.6.30. Let $X$ be an $M$-valued semimartingale, $U$ a horizontal lift of $X$ to $P=\mathrm{L}(T M)$ resp. $\mathrm{O}(T M)$, and $Z$ an anti-development of $X$ into $\mathbb{R}^{n}$. The following statements hold:
(i) $\int_{U} \sigma=\sum_{i=1}^{n} \int \sigma(U) L_{i}(U) \circ d Z^{i}$ for each differential form $\sigma \in \Gamma\left(T^{*} P\right)$;
(ii) $\int_{X} \alpha=\sum_{i=1}^{n} \int \alpha(X) U e_{i} \circ d Z^{i}$ for each differential form $\alpha \in \Gamma\left(T^{*} M\right)$.

In particular, $d(f(U))=\sum_{i=1}^{n}\left(L_{i} f\right)(U) \circ d Z^{i}$ for each function $f \in C^{\infty}(P)$, in short-terms

$$
\begin{equation*}
d U=\sum_{i=1}^{n} L_{i}(U) \circ d Z^{i} \tag{1.6.14}
\end{equation*}
$$

as well as $d(f(X))=\sum_{i=1}^{n}\left(U e_{i}\right)(f) \circ d Z^{i}$ for each function $f \in C^{\infty}(M)$, or in short-terms

$$
\begin{equation*}
d X=U \circ d Z \tag{1.6.15}
\end{equation*}
$$

Proof. The additional claims follow from (i) and (ii) with $\sigma=d f$ for $f \in C^{\infty}(P)$, resp. $\alpha=d f$ for $f \in C^{\infty}(M)$.

To (i): According to Theorem 1.3.9 it is sufficient that the right-hand side of (i) has the defining properties of $\int_{U} \sigma$. For $f \in C^{\infty}(P)$ we have to show that

$$
d(f(U))=\sum_{i}(d f)(U) L_{i}(U) \circ d Z^{i} \equiv \sum_{i}\left(L_{i} f\right)(U) \circ d Z^{i}
$$

which is equivalent to

$$
\begin{equation*}
f(U)-f\left(U_{0}\right)=\int_{U} \sigma \quad \text { where } \sigma \in \Gamma\left(T^{*} P\right), \sigma_{u}:=\sum_{i}\left(L_{i} f\right)(u) \vartheta_{u}^{i} \tag{1.6.16}
\end{equation*}
$$

However observe that $\sum_{i}\left(L_{i} f\right)(u) \vartheta_{u}^{i}=(d f)_{u} \circ \operatorname{pr}_{H_{u}}$, indeed for $A \in T_{u} P$ we have

$$
\begin{aligned}
\sum_{i}\left(L_{i} f\right)(u) \vartheta_{u}^{i}(A) & =\sum_{i}(d f)_{u} L_{i}(u) \vartheta_{u}^{i}(A) \\
& =\sum_{i}(d f)_{u} h_{u}\left(u e_{i}\right)\left(u^{-1}(d \pi)_{u} A\right)^{i} \\
& =(d f)_{u} h_{u}\left(u u^{-1}(d \pi)_{u} A\right) \\
& =(d f)_{u} h_{u}\left((d \pi)_{u} A\right) \\
& =\left((d f)_{u} \circ \operatorname{pr}_{H_{u}}\right)(A) .
\end{aligned}
$$

On the other side, we have $\left(d f \circ \operatorname{pr}_{V}\right)_{u}=(d f)_{u} \kappa_{u} \omega_{u}=d\left(f \circ I_{u}\right)_{e} \omega_{u}$. But $U$ is horizontal and hence $\int_{U} d f \circ \operatorname{pr}_{V}=0$ which shows that

$$
f(U)-f\left(U_{0}\right)=\int_{U} d f=\int_{U} d f \circ \operatorname{pr}_{H}+\int_{U} d f \circ \operatorname{pr}_{V}=\int_{U} d f \circ \operatorname{pr}_{H}=\int_{U} \sigma
$$

The second defining property of the Stratonovich integral is obvious.
To (ii): It is sufficient to show that

$$
d(f(X))=\sum_{i}(d f)(X) U e_{i} \circ d Z^{i} \equiv \sum_{i}\left(U e_{i}\right)(f) \circ d Z^{i}
$$

holds for each function $f \in C^{\infty}(M)$. With part (i) using that $(d \pi)_{u} L_{i}(u)=u e_{i}$, we obtain

$$
\begin{aligned}
d((f \circ \pi)(U)) & =\sum_{i} d(f \circ \pi)(U) L_{i}(U) \circ d Z^{i} \\
& =\sum_{i}(d f)(\pi(U))(d \pi)(U) L_{i}(U) \circ d Z^{i} \\
& =\sum_{i}(d f)(X) U e_{i} \circ d Z^{i},
\end{aligned}
$$

which shows the claim.
Theorem 1.6.31. Let $X$ be an $M$-valued semimartingale, $U$ a horizontal lift of $X$ to $P=\mathrm{L}(T M)$ resp. $\mathrm{O}(T M)$, and $Z$ an anti-development of $X$ into $\mathbb{R}^{n}$. Then
(i) $\int a(d U, d U)=\sum_{i, j=1}^{n} \int a(U)\left(L_{i}(U), L_{j}(U)\right) d\left[Z^{i}, Z^{j}\right]$ for $a \in \Gamma\left(T^{*} P \otimes T^{*} P\right)$;
(ii) $\int b(d X, d X)=\sum_{i, j=1}^{n} \int b(X)\left(U e_{i}, U e_{j}\right) d\left[Z^{i}, Z^{j}\right]$ for $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$.

Proof. It is again sufficient to consider the special case $a=d \varphi_{1} \otimes d \varphi_{2}$ where $\varphi_{1}, \varphi_{2} \in C^{\infty}(P)$, resp. $b=d f_{1} \otimes d f_{2}$ where $f_{1}, f_{2} \in C^{\infty}(M)$. The statements then follow directly from the properties of the quadratic variation or from Theorem 1.6.30 with formula (1.3.10).

Proposition 1.6.32 (Left-invariant SDE on a Lie group). Let Ge a Lie group and $\mathfrak{g}$ the corresponding Lie algebra. We identify

$$
\mathfrak{g} \xrightarrow{\sim}\{\text { left-invariant vector fields on } G\}, \quad A \longmapsto A(\cdot),
$$

where $A(g)=\left(L_{g}\right)_{*} A(e) \equiv\left(L_{g}\right)_{*} A$ and $\left(L_{g}\right)_{*}: \mathfrak{g} \xrightarrow{\sim} T_{g} G$ is the differential of the left multiplication $L_{g}$. Let $A_{1}, \ldots, A_{r} \in \mathfrak{g}$ and $A_{1}(\cdot), \ldots, A_{r}(\cdot) \in \Gamma(T G)$ the corresponding left-invariant vector fields. Let $\gamma$ be a continuous $\mathbb{R}^{r}$-valued semimartingale. Then each maximal solution of the Stratonovich SDE

$$
\begin{equation*}
d g=\sum_{i=1}^{r} A_{i}(g) \circ d \gamma^{i} \tag{1.6.17}
\end{equation*}
$$

has infinite lifetime. If $\left(g_{t}\right)_{t \geq 0}$ is a solution to SDE (1.6.17), then $\tilde{g}_{t}:=g_{t}^{-1}$ satisfies the SDE

$$
\begin{equation*}
d \tilde{g}=-\sum_{i=1}^{r}\left(\operatorname{Ad}\left(\tilde{g}^{-1}\right) A_{i}\right)(\tilde{g}) \circ d \gamma^{i}, \quad \tilde{g}_{0}=g_{0}^{-1} \tag{1.6.18}
\end{equation*}
$$

Proof. (a) Note that $\operatorname{SDE}$ (1.6.18) is equivalent to

$$
\begin{equation*}
d \tilde{g}=-\sum_{i=1}^{r}\left(R_{\tilde{g}}\right)_{*} A_{i}(e) \circ d \gamma^{i} \tag{1.6.19}
\end{equation*}
$$

Let now $\left(\tilde{g}_{t}\right)_{t \geq 0}$ be a semimartingale satisfying (1.6.19) and $\left(g_{t}\right)_{t \geq 0}$ be a solution to (1.6.17). Then we have:

$$
\begin{equation*}
d(f(g \tilde{g}))=f_{*}\left(L_{g}\right)_{*} \circ d \tilde{g}+f_{*}\left(R_{\tilde{g}}\right)_{*} \circ d g=0, \quad f \in C^{\infty}(G) \tag{1.6.20}
\end{equation*}
$$

Indeed letting $Q: G \times G \rightarrow G,(g, \tilde{g}) \mapsto g \tilde{g}=L_{g} \tilde{g}=R_{\tilde{g}} g$ denote multiplication on $G$, by Remark 1.2.24, to verify the first equality in (1.6.20), it is sufficient to show that

$$
(f \circ Q)_{*}: T_{(g, \tilde{g})}(G \times G) \cong T_{g} G \times T_{\tilde{g}} G \rightarrow \mathbb{R}
$$

satisfies the formula:

$$
\begin{equation*}
(f \circ Q)_{*}(v, w)=f_{*}\left(L_{g}\right)_{*} w+f_{*}\left(R_{\tilde{g}}\right)_{*} v \tag{1.6.21}
\end{equation*}
$$

This is however easy to see by curve transport. Let $v$ be represented by the curve $\alpha$ : $\alpha(0)=g, \dot{\alpha}(0)=v$, and analogously $w$ by $\beta: \beta(0)=\tilde{g}, \dot{\beta}(0)=w$, then $(f \circ Q)_{*}(v, w)$ is represented by the $t \mapsto f(\alpha(t) \beta(t))$ at 0 . For this we have

$$
\frac{d}{d t} f(\alpha \beta)=f_{*}\left(L_{\alpha}\right)_{*} \dot{\beta}+f_{*}\left(R_{\beta}\right)_{*} \dot{\alpha}
$$

and hence

$$
\left.\frac{d}{d t}\right|_{t=0} f(\alpha \beta)=f_{*}\left(L_{g}\right)_{*} w+f_{*}\left(R_{\tilde{g}}\right)_{*} v
$$

The second equality in (1.6.20) is then immediate from (1.6.17) and (1.6.19). From (1.6.20) we then conclude that $\left(\tilde{g}_{t}\right)_{t \geq 0} \equiv\left(g_{t}^{-1}\right)_{t \geq 0}$ modulo indistinguishability.
(b) Note that if $\left(g_{t}\right)_{t \geq 0}$ solves SDE (1.6.17) with initial condition $g_{0}=e$ and if $\xi_{0}$ is an $\mathscr{F}_{0}$-measurable $G$-valued random variable, then $\left(g_{t}^{\prime}\right)_{t \geq 0}$ where $g_{t}^{\prime}:=\xi_{0} g_{t}$ is the solution with initial condition $g_{0}^{\prime}=\xi_{0}$.
(c) It remains to verify that the maximal solution to

$$
\begin{equation*}
d g=\sum_{i=1}^{r} A_{i}(g) \circ d \gamma^{i}, \quad g_{0}=e \tag{1.6.22}
\end{equation*}
$$

has infinite lifetime. To this end, we fix a relatively compact open coordinate neighbourhood $V$ of the unit element $e$ in $G$ and construct inductively an increasing sequence $\left(\tau_{n}\right)_{n \geq 0}$ of stopping times:

$$
\tau_{0}=0, \text { and } \tau_{n+1}=\inf \left\{t \geq \tau_{n}: g_{t}^{n} \notin V\right\} \wedge(n+1), \quad n \geq 0
$$

where $g^{n}$ denotes the solution to (1.6.17) on $\left[\tau_{n}, \tau_{n+1}\right]$ satisfying $g_{\tau_{n}}^{n}=e$. A global solution $g \equiv\left(g_{t}\right)_{t \geq 0}$ to (1.6.22) is then inductively put together by $g \mid\left[\tau_{n}, \tau_{n+1}\right]:=g_{\tau_{n}} g^{n}$. It remains to show that it has infinite lifetime which means that $\mathbb{P}\left\{\sup \tau_{n}<\infty\right\}=0$. Let $\gamma=\mu+\beta$ be the Doob-Meyer decomposition of $\gamma$. Possibly after a time transformation, we may assume without restrictions that $[\mu, \mu]_{t}+\sum_{i} \int_{0}^{t}\left|d \beta^{i}\right| \leq$ const $\times t$. We want to show that $\mathbb{P}\left\{\sup \tau_{n}<N\right\}=0$ for each $N \in \mathbb{N}$. To this end, we first note that for any $f \in C^{\infty}(G)$,

$$
\begin{equation*}
\int_{\tau_{n}}^{\tau_{n+1}} \sum_{i=1}^{r}\left(A_{i} f\right)\left(g^{n}\right) \circ d \gamma^{i}=\int_{\tau_{n}}^{\tau_{n+1}} \sum_{i=1}^{r} d f A_{i}\left(g^{n}\right) \circ d \gamma^{i}=f\left(g_{\tau_{n+1}}^{n}\right)-f(e) \tag{1.6.23}
\end{equation*}
$$

On the other hand, since the functions $A_{i}(f) \in C^{\infty}(G)$ are bounded on $\bar{V}$ and since $\tau_{n+1} \wedge N-\tau_{n} \wedge N \rightarrow 0$ for $n \rightarrow \infty$, we get that a.s.

$$
\begin{equation*}
\int_{0}^{N} 1_{] \tau_{n}, \tau_{n+1}\right]} \sum_{i=1}^{r}\left(A_{i} f\right)\left(g^{n}\right) \circ d \gamma^{i} \rightarrow 0, \quad n \rightarrow \infty \tag{1.6.24}
\end{equation*}
$$

Since the left-hand sides of (1.6.23) and (1.6.24) agree on $\left\{\sup \tau_{n}<N\right\}$, we conclude that $\mathbb{P}\left\{\sup \tau_{n}<N\right\}=0$.

DEFINITION 1.6.33 (Canonical one-form of a Lie group). Let $G$ be a Lie group with Lie algebra $\mathfrak{g}$. The one-form $\theta \in \Gamma\left(T^{*} G \otimes \mathfrak{g}\right)$ taking values in $\mathfrak{g}$ and defined by $\theta_{g}\left(A_{g}\right):=$ $\left(L_{g}\right)_{*}^{-1} A_{g}$ is called the canonical one-form on $G$.

REMARK 1.6.34. Let $g=\left(g_{t}\right)_{t \geq 0}$ be a continuous semimartingale taking values in a Lie group $G$, then $\gamma:=\int_{g} \theta$ defines a $\mathfrak{g}$-valued semimartingale which writes as $\gamma=$ $\sum_{i=1}^{r} \gamma^{i} A_{i}$ after fixing a basis $A_{1}, \ldots, A_{r}$ for $\mathfrak{g}$. Note that the semimartingale $\left(g_{t}\right)$ can be recovered from $\gamma$ as solution to the SDE

$$
\begin{equation*}
d X=\sum_{i=1}^{r} A_{i}(X) \circ d \gamma^{i}, \quad X_{0}=g_{0} \tag{1.6.25}
\end{equation*}
$$

This shows that each continuous semimartingale $\left(g_{t}\right)$ on $G$ is solution to an SDE of the form (1.6.25) driven by $\mathfrak{g}$-valued semimartingale $\left(\gamma_{t}\right)$. In particular, according to Proposition 1.6.32, the inverse process $\left(g_{t}^{-1}\right)_{t \geq 0}$ satisfies the SDE

$$
d \tilde{X}=-\sum_{i=1}^{r}\left(R_{\tilde{X}}\right)_{*} A_{i} \circ d \gamma^{i}, \quad \tilde{X}_{0}=g_{0}^{-1}
$$

THEOREM 1.6.35 (Horizontal lifts of $M$-valued semimartingales). Let $P$ be principal $G$-bundle over a differentiable manifold $M$ endowed with a $G$-connection. Furthermore, let $x_{0}$ be an $M$-valued random variable and $u_{0}$ a $P$-valued random variable above $x_{0}$, i.e. $\pi\left(u_{0}\right)=x_{0}$ a.s. Then, to each $M$-valued semimartingale $X$ with $X_{0}=x_{0}$, there exists $a$ unique horizontal lift $U$ of $X$ onto $P$ such that $U_{0}=u_{0}$ a.s.

Proof. We follow the proof in the deterministic case of differentiable curves (Theorem 1.6.17). Without restrictions we may assume that $X$ has infinite lifetime. Choosing a countable covering $\left(V_{k}\right)_{k \geq 0}$ of $M$ by bundle chart domains, Lemma 1.3.1 allows inductively by means of the bundle charts $\varphi: \pi^{-1}\left(V_{k}\right) \xrightarrow{\sim} V_{k} \times G$ to lift $X$ first in some way to $P$, that is to find a $P$-valued semimartingale $\tilde{U}$ such that $\pi(\tilde{U})=X$ and $\tilde{U}_{0}=u_{0}$. The problem is now reduced to determine a $G$-valued semimartingale $\left(\tilde{g}_{t}\right)_{t \geq 0}$ in such a way that $U:=\tilde{U} \tilde{g}$ satisfies the wanted properties. First the connection form $\omega \in \Gamma\left(T^{*} P \otimes \mathfrak{g}\right)$ provides a $\mathfrak{g}$-valued semimartingale $\gamma:=\int_{\tilde{U}} \omega \equiv \int \omega(\circ d \tilde{U})$ which we write as $\gamma=\sum_{i=1}^{r} A_{i} \gamma^{i}$ with respect to a fixed basis basis $A_{1}, \ldots, A_{r}$ of $\mathfrak{g}$. With it we define $\left(g_{t}\right)_{t \geq 0}$ as the maximal solution of the SDE

$$
d g=\sum_{i=1}^{r} A_{i}(g) \circ d \gamma^{i}, \quad g_{0}=e
$$

for which Remark 1.6.32 guarantees that $\left(g_{t}\right)$ has infinite lifetime. Letting $\tilde{g}_{t}:=g_{t}^{-1}$ we want to verify that $U_{t}=\tilde{U}_{t} \tilde{g}_{t}$ is horizontal. According to Remark 1.6.32, the inverted process $\left(\tilde{g}_{t}\right)_{t \geq 0} \equiv\left(g_{t}^{-1}\right)_{t \geq 0}$ solves the SDE

$$
d \tilde{g}=-\sum_{i=1}^{r}\left(\operatorname{Ad}\left(\tilde{g}^{-1}\right) A_{i}\right)(\tilde{g}) \circ d \gamma^{i}, \quad \tilde{g}_{0}=e .
$$

Letting again $\Phi: P \times G \rightarrow P,(u, g) \mapsto u \cdot g$, furthermore $R_{g} \equiv \Phi(\cdot, g)$ and $I_{u} \equiv$ $\Phi(u, \cdot)$, then

$$
\left(\Phi^{*} \omega\right)_{(u, g)}=\left(R_{g}^{*} \omega\right)_{u}+\left(I_{u}^{*} \omega\right)_{g}=\left(R_{g}^{*} \omega\right)_{u}+\theta_{g}
$$

where $\theta \in \Gamma\left(T^{*} G \otimes \mathfrak{g}\right)$ is the canonical one-form on $G$ given in Definition 1.6.33. By means of the pullback formula (1.3.9) for Stratonovich integrals of differential forms along semimartingales and Remark 1.6.14 (ii) one then obtains

$$
\begin{aligned}
\int_{U} \omega=\int_{\Phi(\tilde{U}, \tilde{g})} \omega & =\int_{(\tilde{U}, \tilde{g})} \Phi^{*} \omega=\int_{(\tilde{U}, \tilde{g})}\left(R^{*} \omega+\theta\right) \\
& =\int\left(R_{\tilde{g}}^{*} \omega\right)(\circ d \tilde{U})+\int \theta(\circ d \tilde{g}) \\
& =\int \operatorname{Ad}\left(\tilde{g}^{-1}\right) \omega(\circ d \tilde{U})+\int\left(L_{\tilde{g}^{-1}}\right)_{*}(\circ d \tilde{g})=0
\end{aligned}
$$

since $\omega(\circ d \tilde{U})=\sum A_{i} \circ d \gamma^{i}$ and $d \tilde{g}=-\sum \operatorname{Ad}\left(\tilde{g}^{-1}\right) A_{i}(\tilde{g}) \circ d \gamma^{i}=-\sum\left(R_{\tilde{g}}\right)_{*} A_{i} \circ d \gamma^{i}$. This shows that $U$ is indeed a horizontal process.

Uniqueness of $U$ is immediate, since given two lifts $U$ and $\tilde{U}$ with the wanted properties, then $U=\tilde{U} g$ where $g \equiv\left(g_{t}\right)_{t \geq 0}$ is a $G$-valued semimartingale with $g_{0}=e$, almost surely. By the calculation above we obtain

$$
\omega(\circ d U)=\operatorname{Ad}\left(g^{-1}\right) \omega(\circ d \tilde{U})+\theta(\circ d g)
$$

But $U$ and $\tilde{U}$ are horizontal by assumption, hence $\theta(\circ d g) \equiv\left(L_{g^{-1}}\right)_{*}(\circ d g)=0$ which implies $d g=0$ and thus $g_{t} \equiv g_{0}=e$, almost surely.

The proof of Theorem 1.6 .35 provides a structural statement for semimartingales in $P$ which we state in the case of frame bundles in explicit form.

Corollary 1.6.36. Let $M$ be a differentiable manifold and $P=\mathrm{L}(T M)$ with a $G$-connection where $G=\operatorname{GL}(n ; \mathbb{R})$, respectively, let $M$ be a Riemannian manifold and $P=\mathrm{O}(T M)$ with a $G$-connection where $G=\mathrm{O}(n)$. Assume that $\tilde{U}$ is an arbitrary
semimartingale taking values in $P$. Denote its starting value by $\tilde{U}_{0}=u_{0}$. Integration of the connection form $\omega \in \Gamma\left(T^{*} P \otimes \mathfrak{g}\right)$ and the canonical one-form $\vartheta \in \Gamma\left(T^{*} P \otimes \mathbb{R}^{n}\right)$ along $\tilde{U}$ gives the semimartingales $\gamma=\int_{\tilde{U}} \omega$ with values in $\mathfrak{g}$, respectively $Z=\int_{\tilde{U}} \vartheta$ with values in $\mathbb{R}^{n}$. Fixing a basis $\left(A_{1}, \ldots, A_{r}\right)$ for $\mathfrak{g}$ and writing $\gamma=\gamma^{1} A_{1}+\ldots+\gamma^{r} A_{r}$, we define semimartingales $g_{t}$ taking values in $G$ and $U_{t}$ taking values in $P$ as solutions to the following SDEs:

$$
\begin{array}{ll}
d g=\sum_{i=1}^{r} A_{i}(g) \circ d \gamma^{i}, & g_{0}=e, \quad \text { resp. } \\
d U=\sum_{i=1}^{n} L_{i}(U) \circ d Z^{i}, & U_{0}=u_{0},
\end{array}
$$

where we read $A_{1}, \ldots, A_{r}$ as left-invariant vector fields on $G$ and where $L_{1}, \ldots, L_{d}$ denote the standard-horizontal vector fields on $P$. Then, by definition, $U$ is horizontal and $\tilde{U}=$ $U g$ holds, modulo indistinguishability.

Proof. Along with $U$ also $\tilde{U} g^{-1}$ is a horizontal lift of $\pi(\tilde{U})$; since both coincide for $t=0$ they must be equal.

REMARK 1.6.37. There is an alternative proof of Theorem 1.6 .35 (see [41]) which uses the fact that according to Theorem 1.2.23, each semimartingale $X$ on $M$ can be realized as solution to a Stratonovich SDE of the form

$$
\begin{equation*}
d X=\sum_{i=1}^{\ell} A_{i}(X) \circ d Z^{i}, \quad X_{0}=x_{0} \tag{1.6.26}
\end{equation*}
$$

where $Z$ is an $\mathbb{R}^{\ell}$-valued semimartingale for some $\ell$. Let $\bar{A}_{i} \in \Gamma(T P)$ be the horizontal lift of $A_{i} \in \Gamma(T M)$, i.e. $\bar{A}_{i}(u)=h_{u}\left(A_{i}(\pi u)\right)$ for $u \in P$, and consider the "horizontally lifted SDE" on $P$ :

$$
\begin{equation*}
d U=\sum_{i=1}^{\ell} \bar{A}_{i}(U) \circ d Z^{i}, \quad U_{0}=u_{0} \tag{1.6.27}
\end{equation*}
$$

It is clear that solutions to (1.6.27) are canonical candidates for the wanted horizontal lift. Indeed, we have $d(\pi(U))=\sum_{i}(d \pi)_{U} \bar{A}_{i}(U) \circ d Z^{i} \equiv \sum_{i} A_{i}(\pi(U)) \circ d Z^{i}$ with $\pi\left(U_{0}\right)=x_{0}$, and hence $\pi(U)=X$ by uniqueness of solutions to (1.6.26). On the other hand, we have $\int_{U} \omega=\sum_{i} \int \omega(U) \bar{A}_{i}(U) \circ d Z^{i}=0$. It thus remains to verify that $U$ and $X$ have identical lifetimes which is however not immediately clear from the construction.

We want to summarize the theory developed so far. Let $M$ be a differentiable manifold equipped with a torsion-free connection, or a Riemannian manifold with the LeviCivita connection. Over $M$ we then have the frame bundle $P=\mathrm{L}(T M)$ with the induced $\mathrm{GL}(n ; \mathbb{R})$-connection, respectively the orthonormal frame bundle $P=\mathrm{O}(T M)$ with the induced $\mathrm{O}(n)$-connection.

REMARK 1.6.38. Let $u_{0}$ be a $P$-valued random variable and $x_{0}=\pi\left(u_{0}\right)$. If $X$ is a semimartingale on $M$ with starting value $X_{0}=x_{0}$, then by Theorem 1.6.35 there is a unique horizontal lift $U$ of $X$ such that $U_{0}=u_{0}$. By Definition 1.6.29 the antidevelopment $Z$ of $X$ into $\mathbb{R}^{n}$ (with initial frame $u_{0}$ ) is given as $Z=\int_{U} \vartheta$. Modulo choice of initial conditions $X_{0}=x, U_{0}=u$, each of the three processes $X, U, Z$ determines the two others. Indeed, we have:
(a) $Z$ determines $U$ as solution to the SDE

$$
d U=\sum_{i=1}^{n} L_{i}(U) \circ d Z^{i}, \quad U_{0}=u
$$

(b) $U$ determines $X$ via

$$
X=\pi(U)
$$

(c) $X$ determines $Z$ as

$$
Z=\int_{U} \vartheta
$$

where $U$ is the unique horizontal lift of $X$ to $P$ with $U_{0}=u$.
Typically, one starts with $Z$ on $\mathbb{R}^{n}$ (without restrictions $Z_{0}=0$ ) to determine $X$ on $M$. We call $X$ the stochastic development of $Z$. Stochastic development provides at the same time the horizontal lift $U$ to $P$ with $U_{0}=u_{0}$. The frame $U$ moves then along $X$ by parallel transport. The process $Z$ is recovered via $Z=\int_{U} \vartheta$.

REMARKS 1.6.39. (1) The described procedure depends in an obvious way on the choice $u_{0}$ above $x_{0}$. Choosing instead of $u_{0}$ another $\mathscr{F}_{0}$-measurable $P$-valued random variable $\tilde{u}_{0}$ such that $\pi \circ \tilde{u}_{0}=x_{0}$ a.s. leads to $\tilde{u}_{0}=u_{0} g_{0}$ for an $\mathscr{F}_{0}$-measurable random variable $g_{0}$ taking values in the Lie group $G$ of invertible, respectively orthogonal $n \times n$ matrices, so that $U$ changes to $\tilde{U}=U g_{0}$. Since $R_{g}^{*} \vartheta=g^{-1} \vartheta$ for $g \in G$, the antidevelopment $Z$ transforms to

$$
\begin{equation*}
\tilde{Z}=\int_{\tilde{U}} \vartheta=\int_{U} R_{g_{0}}^{*} \vartheta=\int_{U} g_{0}^{-1} \vartheta=g_{0}^{-1} Z \tag{1.6.28}
\end{equation*}
$$

(2) Writing

$$
d U=\sum_{i=1}^{n} L_{i}(U) \circ d Z^{i}=\sum_{i=1}^{n} h_{U}\left(U e_{i}\right) \circ d Z^{i} \quad \text { and } \quad d X=\sum_{i=1}^{n} U e_{i} \circ d Z^{i}
$$

we arrive at the intrinsic formulas

$$
\begin{equation*}
d U=h_{U}(\circ d X) \quad \text { and } \quad d X=U \circ d Z \tag{1.6.29}
\end{equation*}
$$

(3) Fixing $u \in P$, read as isomorphism (isometry) $u: \mathbb{R}^{n} \xrightarrow{\sim} T_{x} M$ where $x=\pi(u)$, we may identify $Z$ with the $T_{x} M$-valued semimartingale $\tilde{Z}=u Z$. Stochastic development then provides a one-to-one correspondence between continuous semimartingales $\tilde{Z}$ in the tangent space $T_{x} M$ with $\tilde{Z}_{0}=0$ and semimartingales $X$ on the manifold $M$ with $X_{0}=x$, where $\tilde{Z} \longmapsto X=\pi(U)$ and $U$ defined as solution to the SDE

$$
d U=\sum_{i=1}^{n} L_{i}(U) u^{-1} \circ d \tilde{Z}^{i}, \quad U_{0}=u
$$

We want to give a geometric illustration of stochastic development. For instance, let $P=\mathrm{O}(T M)$ the orthonormal frame bundle over a Riemannian manifold $M$. We fix $u \in \mathrm{O}(T M)$ as isometry $u: \mathbb{R}^{n} \xrightarrow{\sim} T_{\pi(u)} M$ and let $x:=\pi(u)$.


Figure 1.6.1. Stochastic development

One should think of $X$ as the trace which the paths of $Z$ print on the manifold $M$, under the identification $U: \mathbb{R}^{n} \xrightarrow{\sim} T_{X} M$, when $M$ is "rolled" along $t \mapsto Z_{t}$ (rolling without slipping). In the probabilistic case however this interpretation requires further explication as in general the trajectories of $Z$ are not differentiable and thus a pathwise procedure does not make immediate sense. Let us thus first have a look at the deterministic case of a differentiable curve $Z: t \mapsto z(t)$. We will show that in this case "stochastic development" reduces to the classical Cartan development of the curve $t \mapsto z(t)$.

EXAMPLE 1.6.40 (Cartan development). The Cartan development of an $\mathbb{R}^{n}$-valued curve $t \mapsto z(t)$ is the construction of curves $x: t \mapsto x(t) \in M$ and $u: t \mapsto u(t) \in P$ (where $P=\mathrm{L}(T M)$, resp. $P=\mathrm{O}(T M)$ in the Riemannian case) such that $u(\cdot)$ lies above $x(\cdot)$ and such that
(i) $\dot{x}=u \dot{z}$;
(ii) $u$ is parallel along $x$.

Condition (i) can be rewritten as

$$
d x(t)=u(t) d z(t)
$$

and " $u$ is parallel along $x$ " is understood in the sense that $\nabla_{D} u \equiv\left(\nabla_{D} u^{1}, \ldots, \nabla_{D} u^{n}\right)=0$ where $D=\partial / \partial t$. Condition (ii) means then that $u(\cdot)$ is a horizontal curve; thus $\dot{u} \in H_{u} \equiv$
$h_{u}\left(T_{\pi(u)} M\right)$, and since $\dot{x}=(\pi \circ u)^{\cdot}=\pi_{*} \dot{u}=u \dot{z}$ we obtain $\dot{u}=h_{u}(\dot{x})=h_{u}(u \dot{z})$ by using (i). On the other hand, since $h_{u}(u \dot{z})=\sum_{i} h_{u}\left(u e_{i}\right) \dot{z}^{i}=\sum_{i} L_{i}(u) \dot{z}^{i}$, conditions (i) and (ii) are seen to be equivalent to

$$
\left\{\begin{array}{l}
d u=\sum_{i=1}^{n} L_{i}(u) d z^{i}  \tag{1.6.30}\\
x(\cdot)=(\pi \circ u)(\cdot)
\end{array}\right.
$$

REMARK 1.6.41. (a) Note that Eq. (1.6.30) is the equation introduced above for the procedure of "rolling without slipping" in the special case of a deterministic driving process $z(t)$. In this case stochastic development reduces to classical Cartan development. In the general case of a non-trivial semimartingale $Z$ the ordinary differential equation (1.6.30) for Cartan development needs to be rewritten as a Stratonovich type SDE.
(b) The term $L_{i}\left(u_{0}\right)$ can be interpreted as infinitesimal direction of the parallel transport of $u_{0} \in P$ along a curve in $M$ with initial velocity $u_{0} e_{i}$ at $\pi\left(u_{0}\right)$, i.e.,

$$
\left.\frac{d}{d \varepsilon}\right|_{\varepsilon=0} / / 0, \varepsilon u_{0}=L_{i}\left(u_{0}\right)
$$

As already explained, we are mainly interested in the case of frame bundles over $M$. We distinguished so far the two cases of the frame bundle $\mathrm{L}(T M)$ and the orthonormal frame bundle $P=\mathrm{O}(T M)$ if $M$ carries in addition a Riemannian metric. We want to check quickly that the two points of view are compatible for the procedure of stochastic development.

REmARK 1.6.42. Let $M$ be a Riemannian manifold equipped with the Levi-Civita connection. The inclusion $(\mathrm{O}(T M), \mathrm{O}(n)) \stackrel{j}{\longrightarrow}(\mathrm{~L}(T M), \mathrm{GL}(n ; \mathbb{R}))$ defines a homomorphism of principal bundles with $\mathfrak{g}=\mathfrak{o}(d) \stackrel{j_{0}}{\longrightarrow} \overline{\mathfrak{g}}:=\mathfrak{g l}(d ; \mathbb{R})$ the inclusion of the corresponding Lie algebras. This gives the following situation:


Let $X$ be an $M$-valued semimartingale and $u_{0}$ an $\mathscr{F}_{0}$-measurable $\mathrm{O}(T M)$-valued random variable. In addition, let $U$ be the horizontal lift of $X$ to $\mathrm{O}(T M)$ and $\bar{U}$ the horizontal lift of $X$ to $\mathrm{L}(T M)$ such that $U_{0}=\bar{U}_{0}=u_{0}$ a.s. Let $Z=\int_{U} \vartheta$ and $\bar{Z}=\int_{\bar{U}} \bar{\vartheta}$. Then, modulo indistinguishability, $Z=\bar{Z}$ and $j(U)=\bar{U}$ hold.

Proof. It is straightforward to see that $j_{*} H=\bar{H}, j^{*} \bar{\omega}=j_{0} \omega$ and $j^{*} \bar{\vartheta}=\vartheta$ where $\omega, \vartheta$ (respectively $\bar{\omega}, \bar{\vartheta}$ ) denote the connection form and canonical one-form on $P=\mathrm{O}(T M)$, respectively on $P=\mathrm{L}(T M)$. This gives

$$
\int_{j(U)} \bar{\omega}=\int_{U} j^{*} \bar{\omega}=j_{0} \int_{U} \omega=0
$$

which by uniqueness of the horizontal lift implies $j(U)=\bar{U}$. On the other hand we have

$$
\int_{j(U)} \bar{\vartheta}=\int_{U} j^{*} \bar{\vartheta}=\int_{U} \vartheta
$$

which shows $\bar{Z}=Z$.

DEFINITION 1.6.43 (Parallel transport along a semimartingale). Let $M$ be a differentiable manifold equipped with a torsion-free connection, or a Riemannian manifold equipped with the Levi-Civita connection. Let $X$ be a semimartingale on $M$ and $U$ an arbitrary horizontal lift of $X$ to $\mathrm{L}(T M)$ resp. $\mathrm{O}(T M)$. For $0 \leq s \leq t$ let $/ /_{s, t}:=U_{t} \circ U_{s}^{-1}$ be given by


The isomorphisms (resp. isometries in the Riemannian case)

$$
/_{0, t}: T_{X_{0}} M \rightarrow T_{X_{t}} M
$$

are called stochastic parallel transport along $X$.
REMARK 1.6.44. The parallel transports $/ /_{0, t}$ extend canonically from the tangent bundle $T M$ to tensors of type $(p, q)$, i.e., to $T M^{\otimes p} \otimes\left(T^{*} M\right)^{\otimes q}$, and then to

$$
\bigoplus_{p, q \geq 0}^{N} T M^{\otimes p} \otimes\left(T^{*} M\right)^{\otimes q}, \quad N \in \mathbb{N} .
$$

Note that, for $\alpha \in \Gamma\left(T^{*} M\right)$ and $A \in \Gamma(T M)$, by definition

$$
\left(/ /_{0, t} \alpha_{X_{0}}\right)\left(A_{X_{t}}\right)=\alpha_{X_{0}}\left(/ / t, 0 A_{X_{t}}\right)
$$

THEOREM 1.6.45 (Geometric Itô formula). Let $M$ be a differentiable manifold endowed with a linear connection $\nabla$ (without restriction $\nabla$ torsion-free). Let $X$ be an $M$-valued semimartingale, $U$ a horizontal lift of $X$ to $L(T M)$ and $Z=\int_{U} \vartheta$ the corresponding anti-development of $X$ into $\mathbb{R}^{n}$. For each $f \in C^{\infty}(M)$ the following formula holds:

$$
\begin{equation*}
d(f(X))=\sum_{i=1}^{n}(d f)(X)\left(U e_{i}\right) d Z^{i}+\frac{1}{2} \sum_{i, j=1}^{n}(\nabla d f)(X)\left(U e_{i}, U e_{j}\right) d\left[Z^{i}, Z^{j}\right] \tag{1.6.31}
\end{equation*}
$$

or in abbreviated form (see Theorem 1.6.31),

$$
\begin{equation*}
d(f(X))=(d f)(U d Z)+\frac{1}{2} \nabla d f(d X, d X) \tag{1.6.32}
\end{equation*}
$$

Proof. From $d U=\sum_{i} L_{i}(U) \circ d Z^{i}$ we first see that

$$
\begin{aligned}
& d(f(X))=d((f \circ \pi)(U))=\sum_{i} L_{i}(f \circ \pi)(U) \circ d Z^{i} \\
& \quad=\sum_{i} L_{i}(f \circ \pi)(U) d Z^{i}+\frac{1}{2} \sum_{i, j} L_{i} L_{j}(f \circ \pi)(U) d\left[Z^{i}, Z^{j}\right]
\end{aligned}
$$

where $L_{i}(f \circ \pi)(u)=d(f \circ \pi)_{u} L_{i}(u)=(d f)_{\pi(u)}(d \pi)_{u} h_{u}\left(u e_{i}\right)=(d f)_{\pi(u)}\left(u e_{i}\right)$. Hence we have $L_{i}(f \circ \pi)(u)=F^{i}(u)$ where $F \equiv F_{d f}: \mathrm{L}(T M) \rightarrow \mathbb{R}^{n}$ is the equivariant function $F^{i}(u)=(d f)_{\pi(u)}\left(u e_{i}\right)$ associated to $d f$ (see Notation 1.6.25). Denoting $\overline{u e_{i}}:=h_{u}\left(u e_{i}\right)$, then by means of Eq. (1.6.13),

$$
L_{i} L_{j}(f \circ \pi)(u)=\left(L_{i} F^{j}\right)(u)=\overline{u e_{i}} F^{j}=\left(\nabla_{u e_{i}} d f\right)_{\pi(u)}\left(u e_{j}\right)=\nabla d f\left(u e_{i}, u e_{j}\right)
$$

from where formula (1.6.31) results.

REMARK 1.6.46. Let $M$ be a Riemannian manifold with its Levi-Civita connection. Denoting by $\Delta^{\text {hor }}=\sum_{i} L_{i}^{2}$ the horizontal Laplacian on $\mathrm{O}(T M)$ and by $\Delta$ the LaplaceBeltrami operator on $M$, then for each $f \in C^{\infty}(M)$ the following relation holds:

$$
\Delta^{\mathrm{hor}}(f \circ \pi)=(\Delta f) \circ \pi
$$

Proof. Indeed, for $u \in \mathrm{O}(T M)$, we have

$$
\sum_{i} L_{i}^{2}(f \circ \pi)(u)=\sum_{i} \nabla d f\left(u e_{i}, u e_{i}\right)=(\operatorname{trace} \nabla d f) \pi(u)=(\Delta f) \circ \pi(u)
$$

Notation 1.6.47. In terms of the Itô integral of the one-form df along $X$, defined as

$$
\begin{equation*}
(\nabla) \int_{X} d f:=\int d f(U d Z) \tag{1.6.33}
\end{equation*}
$$

Eq. (1.6.32) writes as

$$
\begin{equation*}
\int_{X} d f=(\nabla) \int_{X} d f+\frac{1}{2} \int \nabla d f(d X, d X) . \tag{1.6.34}
\end{equation*}
$$

Note that (1.6.33) extends naturally to differential forms $\alpha \in \Gamma\left(T^{*} M\right)$ as

$$
(\nabla) \int_{X} \alpha:=\int \alpha(U d Z)
$$

Stochastic development of $\mathbb{R}^{n}$-valued semimartingales (along with the anti-development of $M$-valued semimartingales into $\mathbb{R}^{n}$ as inverse operation) allows to construct to each class of $\mathbb{R}^{n}$-valued semimartingales a corresponding class of $M$-valued semimartingales. We want to verify next that under the procedure of stochastic development local martingales on $\mathbb{R}^{n}$ correspond to $\nabla$-martingales on $M$, as well as on Riemannian manifolds $\mathrm{BM}\left(\mathbb{R}^{n}\right)$ and $\mathrm{BM}(M, g)$ correspond to each other via stochastic development.

THEOREM 1.6.48. Let $M$ be a differentiable manifold equipped with a torsion-free linear connection $\nabla$. Let $X$ be an $M$-valued semimartingale and $U_{0}$ an $\mathrm{L}(T M)$-valued $\mathscr{F}_{0}$-measurable random variable such that $\pi\left(U_{0}\right)=X_{0}$ a.s. Furthermore let $Z=\int_{U} \vartheta$ be the anti-development of $X$ into $\mathbb{R}^{n}$ with respect to the initial frame $U_{0}$. Then
(i) $X$ is a $\nabla$-martingale on $M$ if and only if $Z$ is a local martingale on $\mathbb{R}^{n}$.
(ii) If $\nabla$ is the Levi-Civita connection to some Riemannian metric $g$ on $M$ and if $U_{0}$ takes its values in $\mathrm{O}(T M)$, then $X$ is a Brownian motion on $(M, g)$ if and only if $Z$ is a Brownian motion on $\mathbb{R}^{n}$ (more precisely, a Brownian motion on $\mathbb{R}^{n}$ stopped at the lifetime $\zeta$ of $X$ ).

Proof. (i) According to Definition 1.4.32 $X$ is a $\nabla$-martingale, if

$$
d(f(X))-\frac{1}{2}(\nabla d f)(d X, d X) \stackrel{m}{=} 0
$$

for functions $f \in C^{\infty}(M)$. By means of the Geometric Itô formula 1.6.45 this means that

$$
\sum_{i}(d f)(X)\left(U e_{i}\right) d Z^{i} \stackrel{\mathrm{~m}}{=} 0
$$

for any $f \in C^{\infty}(M)$ which is easily seen (with the help of Lemma 1.3.1) to be equivalent to the condition that $Z$ is a local martingale.
(ii) According to Definition 1.5.17, the semimartingale $X$ is a Brownian motion on $(M, g)$ if

$$
d(f(X))-\frac{1}{2}(\Delta f)(X) d t \stackrel{\underline{m}}{=} 0
$$

for all $f \in C^{\infty}(M)$. By formula (1.6.31), clearly if $Z$ is a Brownian motion $\mathbb{R}^{n}$, then $X$ will be Brownian motion on $(M, g)$. Conversely, if $X$ is Brownian motion on $(M, g)$ then by Lévy's characterization of $M$-valued Brownian motions (Theorem 1.5.18) $X$ is a $\nabla$-martingale, and thus $Z$ a local martingale by part (i). On the other hand, we have $Z^{i}=\int_{U} \vartheta^{i}$ where $\vartheta_{u}^{i}=\left\langle d \pi(\cdot), u e_{i}\right\rangle=\pi^{*}\left\langle\cdot, u e_{i}\right\rangle$. We may calculate the quadratic variation of $Z$ using Remark 1.3.14 as follows:

$$
\begin{aligned}
d\left[Z^{i}, Z^{j}\right] & =d\left[\int_{U} \vartheta^{i}, \int_{U} \vartheta^{j}\right]=\left(\vartheta^{i} \otimes \vartheta^{j}\right)(d U, d U) \\
& =\pi^{*}\left(\left\langle\cdot, U e_{i}\right\rangle \otimes\left\langle\cdot, U e_{j}\right\rangle\right)(d U, d U) \\
& =\left(\left\langle\cdot, U e_{i}\right\rangle \otimes\left\langle\cdot, U e_{j}\right\rangle\right)(d X, d X) \\
& =\operatorname{trace}\left(\left\langle\cdot, U e_{i}\right\rangle \otimes\left\langle\cdot, U e_{j}\right\rangle\right)(X) d t=\delta_{i j} d t .
\end{aligned}
$$

By means of Lévy's characterization for Brownian motions on $\mathbb{R}^{n}$ we see that $Z$ is a Brownian motion.

REMARK 1.6.49. 1) Theorem 1.6.48 provides a canonical way to construct Brownian motions on Riemannian manifolds. One obtains Brownian motions on $(M, g)$ with starting point $x \in M$ as stochastic development of a Euclidean Brownian motion $B$ on $\mathbb{R}^{n}$. To this end we choose $u \in \mathrm{O}(T M)$ such that $\pi(u)=x$ and solve the SDE

$$
\begin{equation*}
d U=\sum_{i=1}^{n} L_{i}(U) \circ d B^{i}, \quad U_{0}=u \tag{1.6.35}
\end{equation*}
$$

According to Theorem 1.6.48, then $X=\pi(U)$ is a Brownian motion on $(M, g)$ starting from $X_{0}=x$. Note that choosing a different initial frame $u \in \pi^{-1}\{x\}$ in (1.6.35) only changes the underlying Euclidean Brownian motion, in particular, the law of $X$ will be independent of these choices. Indeed, for any $g \in \mathrm{O}(T M)$, along with $B$ also $g B$ is a $\operatorname{BM}\left(\mathbb{R}^{n}\right)$, and hence $X$ constructed by means of $u g$ and $B$ coincides with $X$ constructed by means of $u$ and $g B$.
2) More generally we have the following observation: For an arbitrary $\mathscr{F}_{0}$-measurable $\mathrm{O}(n)$-valued random variable $g_{0}$ along with $B$ also $g_{0} B$ is an $\mathbb{R}^{n}$-valued Brownian motion. Hence if $U$ is the solution to $d U=\sum_{i} L_{i}(U) \circ d\left(g_{0} B\right)^{i} \tilde{\tilde{U}}^{\text {with }} U_{0}=u_{0}$, then $\tilde{U}:=U g_{0}$ solves the $\operatorname{SDE} d \tilde{U}=\sum_{i} L_{i}(\tilde{U}) \circ d B^{i}$ with initial value $\tilde{U}_{0}=u_{0} g_{0}$. Indeed, as a consequence of $\left(R_{g}\right)_{*} h_{u}=h_{u g}$ for $g \in \mathrm{O}(n)$, we have

$$
\begin{aligned}
d\left(U g_{0}\right) & =d\left(R_{g_{0}} U\right)=\left(d R_{g_{0}}\right)_{U} \circ d U=\sum_{i}\left(R_{g_{0}}\right)_{*} L_{i}(U) \circ d\left(g_{0} B\right)^{i} \\
& =\left(R_{g_{0}}\right)_{*} h_{U}\left(U \circ d\left(g_{0} B\right)\right)=h_{U g_{0}}\left(U g_{0} \circ d B\right)=\sum_{i} L_{i}\left(U g_{0}\right) \circ d B^{i}
\end{aligned}
$$

see also the argumentation related to formula (1.6.28).
REMARK 1.6.50. Let $X$ be an $M$-valued semimartingale with starting point $x \in M$. The anti-development $Z$ of $X$ into $\mathbb{R}^{n}$ (see Definition 1.6.29) requires the choice of a frame $u$ above $x$,

$$
Z=\int_{U} \vartheta, \quad U_{0}=u
$$

Considering the anti-development of $X$ into $T_{x} M$, i.e.

$$
\mathscr{A}(X)=U_{0} \int_{U} \vartheta
$$

makes the notion intrinsic. Note that $d(\mathscr{A}(X))=U_{0} U_{t}^{-1} \circ d X$. Our formulas then read as

$$
d(\mathscr{A}(X))=\|_{0, t}^{-1} \circ d X, \quad \text { respectively } d X=\|_{0, t} \circ d(\mathscr{A}(X))
$$

In the same way we have

$$
\begin{equation*}
d U=h_{U}\left(/ /_{0, t} \circ d \mathscr{A}(X)\right) \equiv h_{U}(\circ d X) \tag{1.6.36}
\end{equation*}
$$

The intrinsic version of the Geometric Itô formula (Theorem 1.6.45) takes the form

$$
\begin{equation*}
d(f(X))=(d f)(/ / 0, t \quad d(\mathscr{A}(X)))+\frac{1}{2} \nabla d f(d X, d X) \tag{1.6.37}
\end{equation*}
$$

or in integrated form

$$
\begin{equation*}
\int_{X} d f=(\nabla) \int_{X} d f+\frac{1}{2} \int \nabla d f(d X, d X) \tag{1.6.38}
\end{equation*}
$$

where now

$$
\begin{equation*}
\left.(\nabla) \int_{X} d f=\int(d f)(/ / 0, t) d(\mathscr{A}(X))\right) \tag{1.6.39}
\end{equation*}
$$

see Eq. (1.6.33) for the definition of the Itô integral of $d f$ along $X$.
We want to come back briefly to the deterministic case of development of differentiable curves by pointing out that via development and anti-development geodesics on $M$ correspond to straight lines passing through the origin in $\mathbb{R}^{n}$.

REMARK 1.6.51. Let $M$ be a differentiable manifold, $\nabla$ a linear connection on $M$ and $u_{0} \in \mathrm{~L}(T M)$ fixed. To each curve $t \mapsto z(t)$ in $\mathbb{R}^{n}$ with $z(0)=0$ we consider its development $t \mapsto \gamma(t)$ on $M$ with $\dot{\gamma}(0)=u_{0} \dot{z}(0)$. (Or conversely: to a curve $t \mapsto \gamma(t)$ in $M$ with $\gamma(0)=\pi\left(u_{0}\right)$ we consider its "anti-development" $z(\cdot)=\int_{u} \vartheta$ where $t \mapsto u(t)$ is the horizontal lift of $\gamma$ to $\mathrm{L}(T M)$ with initial value $\left.u(0)=u_{0}\right)$. Then $t \mapsto \gamma(t)$ is a geodesic on $M$ if and only if $z(t)=\dot{z}(0) t$ for each $t$.

Proof. Suppose first that $t \mapsto \gamma(t)$ is a geodesic on $M$. Then both $\dot{\gamma}$ and $u(\cdot) \dot{z}(0)$ are parallel sections along $\gamma$ satisfying $\dot{\gamma}(0)=u(0) \dot{z}(0)$. By Theorem 1.4.11 hence $\dot{\gamma}(s)=$ $u(s) \dot{z}(0)$, and we have

$$
\begin{aligned}
\left(\int_{u} \vartheta\right)(t) & =\int_{0}^{t} \vartheta(\dot{u}(s)) d s=\int_{0}^{t} u(s)^{-1} \pi_{*} \dot{u}(s) d s \\
& =\int_{0}^{t} u(s)^{-1} \dot{\gamma}(s) d s=\int_{0}^{t} \dot{z}(0) d s=\dot{z}(0) t
\end{aligned}
$$

Conversely, if $z(t)=\dot{z}(0) t$ then $\dot{u}(t)=h_{u(t)}(u(t) \dot{z}(t))=h_{u(t)}(u(t) \dot{z}(0))$ and hence $\dot{\gamma}(t)=(\pi \circ u)^{\cdot}(t)=\pi_{*} h_{u(t)}(u(t) \dot{z}(0))=u(t) \dot{z}(0) \equiv / /_{0, t} \dot{\gamma}(0)$. This shows that $\dot{\gamma}$ is parallel along $\gamma$.

Definition 1.6.52. Let $M$ be a differentiable manifold, $\nabla$ a torsion-free linear connection on $M$ and $x \in M$ a point in $M$. Furthermore let $X$ be an $M$-valued semimartingale with $X_{0}=x$ and $U$ be a horizontal lift of $X$ to $\mathrm{L}(T M)$ such that $U_{0}=u_{0} \in \pi^{-1}\{x\}$. The semimartingale $X$ is called one-dimensional if there exists a real-valued semimartingale $Z^{1}$ and a vector $a \in \mathbb{R}^{n}$ such that the anti-development $Z=\int_{U} \vartheta$ of $X$ into $\mathbb{R}^{n}$ takes the form $Z=Z^{1} a$. In addition, $X$ is called one-dimensional martingale, respectively onedimensional Brownian motion, if $Z^{1}$ is even a real local martingale, respectively $\mathrm{BM}(\mathbb{R})$.

The properties above obviously do not depend on the choice of $u \in \pi^{-1}\{x\}$.

THEOREM 1.6.53 (One-dimensional semimartingales move along geodesics). Let $M$ be a differentiable manifold, $\nabla$ a torsion-free linear connection on $M$ and $X$ a semimartingale taking values in $M$ with $X_{0}=x \in M$. Then:
(i) $X$ is a one-dimensional semimartingale if and only if there exist a geodesic $\gamma: I \rightarrow$ $M$ (defined on some open interval $I \subset \mathbb{R}$ ) and a real semimartingale $X^{\prime}$ taking values in $I$ such that $X_{0}^{\prime}=$ const and $X=\gamma\left(X^{\prime}\right)$.
(ii) $X$ is a one-dimensional martingale (one-dimensional Brownian motion) if and only if
$X=\gamma\left(X^{\prime}\right)$ as in (i) and $X^{\prime}$ is in addition a continuous local martingale (Brownian motion).

Proof. Let $X$ be an $M$-valued semimartingale with $X_{0}=x$ and $U$ a horizontal lift of $X$ to $\mathrm{L}(T M)$ with $U_{0}=u_{0}$ for some $u_{0} \in \pi^{-1}\{x\}$. Furthermore let $Z=\int_{U} \vartheta$ be the anti-development of $X$ in $\mathbb{R}^{n}$.
(1) First assume $Z=Z^{\prime} a$ (where $Z_{0}^{\prime}=0$ ). Then $U$ satisfies the SDE

$$
\begin{equation*}
d U=\sum_{i=1}^{n} L_{i}(U) \circ d Z^{i}=L_{a}(U) \circ d Z^{\prime}, \quad U_{0}=u_{0} \tag{1.6.40}
\end{equation*}
$$

where the horizontal vector field $L_{a}$ on $\mathrm{L}(T M)$ is given by $L_{a}(u)=h_{u}(u a)$. Let $t \mapsto u(t)$ be the maximal flow curve to $L_{a}$ with initial value $u(0)=u_{0}$, i.e., $\dot{u}(t)=L_{a}(u(t))$ with $u(0)=u_{0}$. Then the projection $\gamma:=\pi(u)$ defines a geodesic on $M$ : indeed $\dot{\gamma}=(d \pi)_{u} \dot{u}=$ $(d \pi)_{u} h_{u}(u a)=u a$ shows that $\dot{\gamma}$ is parallel along $\gamma$. On the other hand, we have

$$
d\left(u\left(Z^{\prime}\right)\right)=\dot{u}\left(Z^{\prime}\right) \circ d Z^{\prime}=L_{a}\left(u\left(Z^{\prime}\right)\right) \circ d Z^{\prime}, \quad\left(u\left(Z^{\prime}\right)\right)_{0}=u(0)=u_{0}
$$

so that by uniqueness of solutions to Eq. (1.6.40) we get $U=u\left(Z^{\prime}\right)$ modulo indistinguishability. This implies $X=(\pi \circ u)\left(Z^{\prime}\right)=\gamma\left(Z^{\prime}\right)$. With $X^{\prime}:=Z^{\prime}$ we get the claim.
(2) Conversely, suppose that $X=\gamma\left(X^{\prime}\right)$ for some geodesic $\gamma$ and a real semimartingale $X^{\prime}$ where by assumption $X_{0}^{\prime}=$ const. Without restrictions we may assume $X_{0}^{\prime}=0$. Letting $t \mapsto u(t)$ be the horizontal lift of $\gamma$ to $\mathrm{L}(T M)$ with $u(0)=u_{0}$, we get by $\tilde{U}:=u\left(X^{\prime}\right)$ a semimartingale on $\mathrm{L}(T M)$ which projects to $X$ and satisfies $\int_{\tilde{U}} \omega=0$ by the pullback formula (1.3.9) since trivially $\omega(\dot{u}) \equiv 0$. Hence $\tilde{U}$ is a horizontal lift of $X$ with $U_{0}=\tilde{U}_{0}$ a.s. and thus by uniqueness $U=\tilde{U}$ modulo indistinguishability. On the other hand, $\dot{\gamma}$ is parallel along $\gamma$ and hence $\dot{\gamma}(\cdot)=u(\cdot) a$ for some $a \in \mathbb{R}^{n}$ from where we get $\dot{\gamma}\left(X^{\prime}\right)=U a$. The last equality implies

$$
d X=\dot{\gamma}\left(X^{\prime}\right) \circ d X^{\prime}=U a \circ d X^{\prime}=U \circ d\left(X^{\prime} a\right)=\sum_{i}\left(U e_{i}\right) \circ d\left(X^{\prime} a\right)^{i}
$$

and hence $d U=\sum_{i} L_{i}(U) \circ d\left(X^{\prime} a\right)^{i}$ from where $Z=X^{\prime} a$ follows. Hence $X$ is a one-dimensional semimartingale.
(3) Part (ii) of the Theorem is obvious since according to (1) and (2) we may choose $X^{\prime}=Z^{1}$.

### 1.7. Morphisms of Martingales and Brownian Motions

In Section 1.6 we have seen in great generality how to construct martingales and Brownian motions on manifolds. In this Section we are going to give functional characterizations of martingales and Brownian motions, in terms of their behaviour under transformations by maps between manifolds. It will turn out that only very specific maps, so-called harmonic morphisms, map Brownian motions to Brownian motions. Harmonic morphisms in higher dimensions are difficult to find. If however it is only required that Brownian motions are
transformed to martingales, then there is the larger class of harmonic maps. Conversely, harmonic maps are completely characterized by this property. This point of view leads to the general goal in this Section of studying maps between manifolds by analyzing how they change the stochastic behaviour of certain classes of manifold-valued stochastic processes.

Before introducing the necessary vocabulary, we want to briefly summarize how linear connections in vector bundles canonically induce connections in new vector bundles obtained by vector bundle operations.

REMARK 1.7.1. Let $\pi: E \rightarrow M$ be a vector bundle over a differentiable manifold $M$ and

$$
\Gamma(T M) \times \Gamma(E) \rightarrow \Gamma(E), \quad(V, A) \mapsto \nabla_{V} A
$$

a linear connection in $E$. Then, according to Leibniz rule, $\nabla$ extends to a linear linear connection on

$$
\bigoplus_{r, s=0}^{N}(\underbrace{E \otimes \ldots \otimes E}_{r} \otimes \underbrace{E^{*} \otimes \ldots \otimes E^{*}}_{s})
$$

More specifically, if $\nabla^{E}, \nabla^{F}$ are linear connections in $E$, respectively $F$ (both vector bundles over $M$ ), then we have

- the direct sum of the connections $\nabla$ in $E \oplus F$ defined by

$$
\nabla_{V}(A \oplus B)=\nabla_{V}^{E} A \oplus \nabla_{V}^{F} B, \quad V \in \Gamma(T M), A \in \Gamma(E), B \in \Gamma(F)
$$

- the product connection $\nabla$ in $E \otimes F$ defined by
$\nabla_{V}(A \otimes B)=\left(\nabla_{V}^{E} A\right) \otimes B+A \otimes\left(\nabla_{V}^{F} B\right), \quad V \in \Gamma(T M), A \in \Gamma(E), B \in \Gamma(F) ;$
- the dual connection $\nabla^{E^{*}}$ in $E^{*}$ (see Definition 1.4.28) defined by

$$
\left(\nabla_{V}^{E^{*}} \alpha\right)(A)=V(\alpha A)-\alpha\left(\nabla_{V} A\right), \quad V \in \Gamma(T M), \alpha \in \Gamma\left(E^{*}\right), A \in \Gamma(E)
$$

- the pullback connection $\nabla=\nabla^{f^{*} E}$ in $f^{*} E$ (see Definition 1.4.7) for a differentiable map $f: M \rightarrow N$ and $E$ a vector bundle over $N$, determined by

$$
\nabla_{v} f^{*} A=\nabla_{f_{*} v}^{E} A, \quad v \in T M, A \in \Gamma(E)
$$

where $f^{*} A=A \circ f \in \Gamma\left(f^{*} E\right)$.
It is easy to see that pullback of connections is compatible with the other operations for vector bundles.

Lemma 1.7.2. Let $E, F$ be vector bundle over a differentiable manifold $M$ and $\nabla^{E}, \nabla^{F}$ linear connections in $E$ respectively $F$. Furthermore let $\phi \in \Gamma\left(E^{*} \otimes F\right)$, i.e., a homomorphism of vector bundles $\phi: E \rightarrow F$ over $M$. For $B \in \Gamma(E)$, let $\phi B \in \Gamma(F)$ where $(\phi B)_{x}:=\phi_{x} B_{x}$. Then:

$$
\nabla_{A}^{F}(\phi B)=\left(\nabla_{A}^{E^{*} \otimes F} \phi\right) B+\phi \nabla_{A}^{E} B, \quad A \in \Gamma(T M)
$$

Proof. By linearity we can restrict ourselves to the case $\phi=e \otimes \varphi$ where $e \in \Gamma\left(E^{*}\right)$ and $\varphi \in \Gamma(F)$, but then

$$
\begin{aligned}
\nabla_{A}^{F}((e \otimes \varphi) B) & =\nabla_{A}^{F}((e B) \varphi)=A(e B) \varphi+(e B) \nabla_{A}^{F} \varphi \\
& =\left(\left(\nabla_{A}^{E^{*}} e\right) B+e \nabla_{A}^{E} B\right) \varphi+(e B) \nabla_{A}^{F} \varphi \\
& =\left(\left(\left(\nabla_{A}^{E^{*}} e\right) B\right) \varphi+(e B) \nabla_{A}^{F} \varphi\right)+e\left(\nabla_{A}^{E} B\right) \varphi \\
& =\left(\nabla_{A}^{E^{*} \otimes F}(e \otimes \varphi)\right) B+(e \otimes \varphi)\left(\nabla_{A}^{E} B\right)
\end{aligned}
$$

DEFINITION 1.7.3 (Affine and convex mappings). Let $M$ and $N$ be differentiable manifolds, endowed with a torsion-free linear connection in $T M$, respectively $T N$, and $f: M \rightarrow N$ be a differentiable map. For each $x \in M$ the differential $(d f)_{x}: T_{x} M \rightarrow$ $T_{f(x)} N$ of $f$ at $x$ is a linear map. The covariant derivative of the section

$$
d f \in \Gamma\left(T^{*} M \otimes f^{*} T N\right)
$$

gives the Hessian or second fundamental form of $f$,

$$
\nabla d f \in \Gamma\left(T^{*} M \otimes T^{*} M \otimes f^{*} T N\right), \quad(\nabla d f)(A, B)=\left(\nabla_{A} d f\right) B \in \Gamma\left(f^{*} T N\right)
$$

For each $x \in M$ this gives a bilinear form $(\nabla d f)_{x}: T_{x} M \times T_{x} M \rightarrow T_{f(x)} N$. The map $f$ is called affine or totally geodesic if $\nabla d f \equiv 0$.

In the special case $N=\mathbb{R}$, the map $f$ is called convex at $x$ if $(\nabla d f)_{x} \geq 0$ (i.e., positively semidefinite), and strictly convex at $x$ if $(\nabla d f)_{x}>0$ (i.e., positively definite). Finally, $f$ is called convex, respectively strictly convex, if $f$ is convex, respectively strictly convex at each $x \in M$.

REMARK 1.7.4. Let $\pi: \mathrm{L}(T M) \rightarrow M$ be the frame bundle over a differentiable manifold $M$. Then $L_{i} L_{j}(f \circ \pi)(u)=\nabla d f\left(u e_{i}, u e_{j}\right)$ (see the proof of Theorem 1.6.45), and hence $f$ is convex at $x$ if and only if $\left(L_{i} L_{j}(f \circ \pi)\right)_{1 \leq i, j \leq n}$ is positively semidefinite along the fiber $\pi^{-1}\{x\}$.

DEFINITION 1.7.5 (Energy density, tension field, harmonic map). Let ( $M, g$ ) and $(N, h)$ be Riemannian manifolds, endowed with the Levi-Civita connection. To a differentiable map $f: M \rightarrow N$ we have the two fundamental forms, namely
(i) the first fundamental form of $f$ defined as pullback $f^{*} h$ of the metric $h$ under $f$, i.e., $f^{*} h \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$ where $\left(f^{*} h\right)_{x}(u, v)=h_{f(x)}\left(f_{*} u, f_{*} v\right)$ for $u, v \in T_{x} M$;
(ii) the second fundamental form $\nabla d f \in \Gamma\left(T^{*} M \otimes T^{*} M \otimes f^{*} T N\right)$ of $f$ defined as covariant derivative of $d f \in \Gamma\left(T^{*} M \otimes f^{*} T N\right)$.
Taking trace with respect to the given metrics gives

$$
\begin{aligned}
& \text { trace } f^{*} h=|d f|^{2} \in C^{\infty}(M) \quad \text { (the energy density of } f \text { ), } \\
& \text { trace } \nabla d f=\tau(f) \in \Gamma\left(f^{*} T N\right) \quad \text { (the tension field of } f \text { ). }
\end{aligned}
$$

Mappings $f \in C^{\infty}(M, N)$ with vanishing tension field $\tau(f)=0$ are called harmonic. In the special case $N=\mathbb{R}$, the map $f$ is called subharmonic if $\tau(f)=\Delta f \geq 0$.

Lemma 1.7.6. Let $M$ and $N$ be differentiable manifolds, endowed with a torsionfree linear connection in $T M$, respectively $T N$, and $f: M \rightarrow N$ a differentiable map. For $B \in \Gamma(T M)$ let df $B \in \Gamma\left(f^{*} T N\right)$ be defined by $(d f B)_{x}=(d f)_{x} B_{x} \in T_{f(x)} N$. Then for $A, B \in \Gamma(T M)$ :

$$
\begin{equation*}
\nabla_{A}^{f^{*} T N}(d f B)=\left(\nabla_{A} d f\right) B+d f \nabla_{A} B \tag{1.7.1}
\end{equation*}
$$

or equivalently: $(\nabla d f)(A, B)=\nabla_{A}^{f^{*} T N}(d f B)-d f \nabla_{A} B$.
Proof. The claim is a direct consequence of Lemma 1.7.2 with $E=T M, F=$ $f^{*} T N$ and $\phi=d f$.

COROLLARY 1.7.7. In the situation of Lemma 1.7 .6 the bilinear form $\nabla d f$ is symmetric, i.e.,

$$
\nabla d f(A, B)=\nabla d f(B, A) \quad \text { for } A, B \in \Gamma(T M)
$$

Proof. Since the connection on $N$ is torsion-free, we get from the first of Cartan's structural equations (see Theorem 1.4.27) the relation

$$
\nabla_{A}^{f^{*} T N}(d f B)-\nabla_{B}^{f^{*} T N}(d f A)=d f[A, B]
$$

Since also the connection on $M$ is torsion-free, i.e., $\nabla_{A} B-\nabla_{B} A=[A, B]$, the claim follows from Eq. (1.7.1):

$$
\begin{aligned}
\nabla d f(A, B)-\nabla d f(B, A) & =\nabla_{A}^{f^{*} T N}(d f B)-\nabla_{B}^{f^{*} T N}(d f A)-d f\left(\nabla_{A} B-\nabla_{B} A\right) \\
& =d f[A, B]-d f[A, B]=0
\end{aligned}
$$

THEOREM 1.7.8 (Composition formula). Let $M \xrightarrow{f} N \xrightarrow{\varphi} N^{\prime}$ be smooth maps between differentiable manifolds, each manifold endowed with a torsion-free connection. For the Hessian of $\varphi \circ f$ it holds:

$$
\begin{equation*}
\nabla d(\varphi \circ f)=\varphi_{*} \nabla d f+f^{*} \nabla d \varphi \tag{1.7.2}
\end{equation*}
$$

In the case of Riemannian manifolds this gives

$$
\begin{equation*}
\tau(\varphi \circ f)=\varphi_{*} \tau(f)+\operatorname{trace}\left(f^{*} \nabla d \varphi\right) \tag{1.7.3}
\end{equation*}
$$

Proof. For the verification of Eq. (1.7.2) we use Lemma 1.7.2 with $E=f^{*} T N$, $F=(\varphi \circ f)^{*} T N^{\prime} \equiv f^{*}\left(\varphi^{*} T N^{\prime}\right)$ and $\phi=f^{*} d \varphi$ (then $\phi_{x}=(d \varphi)_{f(x)}$ for $x \in M$; see Example 1.0.30). For vector fields $A, B \in \Gamma(T M)$ this gives the formula

$$
\begin{aligned}
\nabla_{A}^{F}(d(\varphi \circ f) B) & =\nabla_{A}^{F}\left(\left(f^{*} d \varphi\right) d f B\right) \\
& =\left(\nabla_{A}^{E^{*} \otimes F}\left(f^{*} d \varphi\right)\right)(d f B)+\left(f^{*} d \varphi\right) \nabla_{A}^{E}(d f B) \\
& =\left(f^{*} \nabla d \varphi\right)(A, B)+\left(f^{*} d \varphi\right) \nabla_{A}^{E}(d f B)
\end{aligned}
$$

where the last equality comes from the definition of the pullback connection on $E^{*} \otimes F \cong$ $f^{*}\left(T^{*} N \otimes \varphi^{*} T N^{\prime}\right)$. Altogether this gives

$$
\begin{aligned}
\nabla d(\varphi \circ f)(A, B) & =\left(\nabla_{A} d(\varphi \circ f)\right) B=\nabla_{A}(d(\varphi \circ f) B)-d(\varphi \circ f) \nabla_{A} B \\
& =\left(f^{*} \nabla d \varphi\right)(A, B)+\left(f^{*} d \varphi\right) \nabla_{A}^{E}(d f B)-\left(f^{*} d \varphi\right) d f \nabla_{A} B \\
& =\left(f^{*} \nabla d \varphi\right)(A, B)+\left(f^{*} d \varphi\right) \nabla d f(A, B) .
\end{aligned}
$$

Eq. (1.7.3) follows from Eq. (1.7.2) by taking trace.
REMARK 1.7.9. Theorem 1.7.8 shows in particular that also the composition $\varphi \circ f$ is affine if $f$ and $\varphi$ are affine. In case of $f$ harmonic and $\varphi$ affine, also $\varphi \circ f$ is harmonic. However, in general, the composition of harmonic maps is not again harmonic.

Corollary 1.7.10. Let $M$ be a manifold endowed with a torsion-free linear connection and $f: M \rightarrow \mathbb{R}$ be a differentiable function. The following characterizations hold:
(i) $f$ is affine if and only if the composition $f \circ \gamma$ is affine, i.e., $(f \circ \gamma)^{\prime \prime} \equiv 0$ for any geodesic $\gamma: I \rightarrow M(I \subset \mathbb{R}$ interval $)$.
(ii) $f$ is convex (resp. strictly convex) if and only if for each geodesic curve $\gamma: I \rightarrow M$ the composition $f \circ \gamma$ is convex (resp. strictly convex), i.e., $(f \circ \gamma)^{\prime \prime} \geq 0$ (resp. $>0$ ).
Proof. First observe that for a smooth curve $\gamma$ on $M$ by Eq. (1.7.2)

$$
(f \circ \gamma)^{\prime \prime}=f_{*} \nabla d \gamma+\gamma^{*} \nabla d f
$$

Since $(\nabla d \gamma)(D, D)=\nabla_{D} \dot{\gamma}$, a curve $\gamma$ is a geodesic if and only if $\gamma$ is affine. On the other hand, $\left(\gamma^{*} \nabla d f\right)(D, D)=(\nabla d f)(\dot{\gamma}, \dot{\gamma})$ so that for geodesic curves $\gamma: I \rightarrow M$ the equation $(f \circ \gamma)^{\prime \prime}(t)=(\nabla d f)(\dot{\gamma}(t), \dot{\gamma}(t))$ holds from where all claims are immediate.

Corollary 1.7.11. Let $M, N$ be differentiable manifolds endowed with torsion-free linear connections. A differentiable map $f: M \rightarrow N$ is affine if and only if $f$ transfers geodesics on $M$ to geodesics on $N$.

Proof. As already noted, for curves $\gamma$ on $M$, "affine" has the same meaning "geodesic" so that the claim follows from $\nabla d(f \circ \gamma)=f_{*} \nabla d \gamma+\gamma^{*} \nabla d f$.

We now return to random motions on manifolds with the goal to investigate maps between manifolds under the aspect of how they transform classes of processes such as Brownian motions or $\nabla$-martingales. To motivate this procedure we consider first the example of Brownian motions on $(M, g)$.

Let $(M, g)$ and $(N, h)$ be Riemannian manifolds, each endowed with the Levi-Civita connection, and $f: M \rightarrow N$ a differentiable map. Let $X$ be a $\operatorname{BM}(M, g)$ starting in $x \in M$ and fix $u \in \mathrm{O}(T M)$ over $x$. There is a unique horizontal lift $U$ of $X$ to $\mathrm{O}(T M)$ such that $U_{0}=u$. The lifted Brownian motion $U$ is a flow process to $\frac{1}{2} \Delta^{\text {hor }}$ where

$$
\Delta^{\mathrm{hor}} \equiv \sum_{i} L_{i}^{2}
$$

is the horizontal Laplacian on $\mathrm{O}(T M)$, and $U$ is called horizontal Brownian motion on the orthonormal frame bundle $\mathrm{O}(T M)$. On the other hand, $X$ comes by stochastic development from the Euclidean Brownian motion $Z=\int_{U} \vartheta$ in $\mathbb{R}^{n}$. Recall that the antidevelopment $\mathscr{A}(X)=u \int_{U} \vartheta$ of $X$ takes values in $T_{x} M$ and is independent of the choice of $U_{0}=u$.


The process $\tilde{X}:=f(X)$ on the target manifold $N$ is in general no longer a Brownian motion. By definition, it is however a semimartingale on $N$ (with $f(x)$ as starting point). We may take a horizontal lift $\tilde{U}$ of $\tilde{X}$ to $\mathrm{O}(T N)$ where $\tilde{U}_{0}=\tilde{u}$ for some $\tilde{u} \in \mathrm{O}(T N)$ above $f(x)$. In addition, we have the anti-development $\mathscr{A}(\tilde{X})$ of $\tilde{X}$ which by definition is a semimartingale taking values in $T_{f(x)} N$.


Figure 1.7.1. Anti-development of the target process
The idea is now to use the Doob-Meyer decomposition $d \tilde{Z}=d \tilde{Z}^{\text {Mart }}+d \tilde{Z}^{\text {drift }}$ of $\tilde{Z}$ to gain information about $f$. In particular, we shall see that the energy density $|d f|^{2}$ and the tension field $\tau(f)$ of $f$ can be recovered from the knowledge of $\tilde{Z}$, respectively $\mathscr{A}(\tilde{X})$. Before treating the case of a Brownian motion $X$ we want first consider the general situation.

THEOREM 1.7.12. Let $M$ and $N$ be differentiable manifolds, each endowed with a torsion-free linear connection, and $f: M \rightarrow N$ be a differentiable map. Furthermore, let $X$ be a semimartingale on $M$ and $\mathscr{A}(X)$ its anti-development to $T_{X_{0}} M$; correspondingly let $\mathscr{A}(\tilde{X})$ be the anti-development of $\tilde{X}:=f(X)$ taking values in $T_{f\left(X_{0}\right)} M$. Finally, let $U$ be a horizontal lift of $X$ to $\mathrm{L}(T M)$, respectively $\tilde{U}$ a horizontal lift of $\tilde{X}$ to $\mathrm{L}(T N)$. Then it holds

$$
\begin{equation*}
d \mathscr{A}(\tilde{X})=\widetilde{/}_{0, \bullet}^{-1}(d f)_{X} / /_{0, \bullet} d \mathscr{A}(X)+\frac{1}{2} \widetilde{/}_{0, \bullet}^{-1} \nabla d f(d X, d X) \tag{1.7.4}
\end{equation*}
$$

where $/ /_{0, t}=U_{t} \circ U_{0}^{-1}$ denotes parallel transport along $X$, respectively $\widetilde{/}_{0, t}=\tilde{U}_{t} \circ \tilde{U}_{0}^{-1}$ along $\tilde{X}$. Here $d f \equiv f_{*}$ is the tangent map to $f$, i.e., $d f_{x}: T_{x} M \rightarrow T_{f(x)} N$ for $x \in M$.

REMARK 1.7.13. In terms of the processes $Z=\int_{U} \vartheta$ in $\mathbb{R}^{\operatorname{dim} M}$, respectively $\tilde{Z}=$ $\int_{\tilde{U}} \vartheta$ in $\mathbb{R}^{\operatorname{dim} N}$, formula (1.7.4) writes as

$$
\begin{equation*}
d \tilde{Z}=\tilde{U}^{-1}(d f)_{X} U d Z+\frac{1}{2} \tilde{U}^{-1} \nabla d f(d X, d X) \tag{1.7.5}
\end{equation*}
$$

where $\tilde{U}^{-1}(d f)_{X} U d Z=\sum_{i} \tilde{U}^{-1}(d f)_{X} U e_{i} d Z^{i}$ and

$$
\tilde{U}^{-1} \nabla d f(d X, d X)=\sum_{i, j} \tilde{U}^{-1} \nabla d f\left(U e_{i}, U e_{j}\right) d Z^{i} d Z^{j}
$$

Proof of Theorem 1.7.12. Let $\varphi \in C^{\infty}(N)$. On one hand, we have by the geometric Itô formula (1.6.32)

$$
d(\varphi(\tilde{X}))=\varphi_{*} \widetilde{/}_{0, t} d \mathscr{A}(\tilde{X})+\frac{1}{2} \nabla^{N} d \varphi(d \tilde{X}, d \tilde{X})
$$

where $\nabla^{N} d \varphi(d \tilde{X}, d \tilde{X})=\left(f^{*} \nabla^{M} d \varphi\right)(d X, d X)$ by the pullback formula (Theorem 1.3.8). On the other hand, we can equally write

$$
d(\varphi(\tilde{X}))=d((\varphi \circ f)(X))=(\varphi \circ f)_{*} / /_{0, t} d \mathscr{A}(X)+\frac{1}{2} \nabla^{M} d(\varphi \circ f)(d X, d X)
$$

where $\nabla^{M} d(\varphi \circ f)(d X, d X)=\left(\varphi_{*} \nabla d f+f^{*} \nabla d \varphi\right)(d X, d X)$ according to the composition formula (Theorem 1.7.8). Comparing the two formulae shows that for each $\varphi \in C^{\infty}(N)$ it holds

$$
\varphi_{*} \widetilde{/ /}_{0, t} d \mathscr{A}(\tilde{X})=\varphi_{*} f_{*} / /_{0, t} d \mathscr{A}(X)+\frac{1}{2} \varphi_{*} \nabla d f(d X, d X)
$$

and thus

$$
\widetilde{/}_{0, t} d \mathscr{A}(\tilde{X})=f_{*} / /_{0, t} d \mathscr{A}(X)+\frac{1}{2} \nabla d f(d X, d X)
$$

which gives the claim.
COROLLARY 1.7.14. Let $(M, g)$ and $(N, h)$ be Riemannian manifolds, endowed with the Levi-Civita connection, and $f: M \rightarrow N$ a differentiable map. Let now $X$ be a Brownian motion on $(M, g)$ starting at $X_{0}=x \in M$. Then $\mathscr{A}(X)$ is a Brownian motion in $T_{x} M$, and for $\tilde{X}=f(X)$ on $N$ it holds

$$
\begin{equation*}
d \mathscr{A}(\tilde{X})=\tilde{/}_{0, \bullet}^{-1}(d f)_{X} / /_{0, \bullet} d \mathscr{A}(X)+\frac{1}{2} \tilde{/ /}_{0, \bullet}^{-1} \tau(f) d t . \tag{1.7.6}
\end{equation*}
$$

In addition, we have

$$
\begin{equation*}
h(d \tilde{X}, d \tilde{X})=|d f|^{2}(X) d t \tag{1.7.7}
\end{equation*}
$$

Proof. We now work with the orthonormal frame bundles $\mathrm{O}(T M)$, respectively $\mathrm{O}(T N)$. Let $U$ and $\tilde{U}$ be horizontal lifts of $X$ to $\mathrm{L}(T M)$, respectively of $\tilde{X}$ to $\mathrm{L}(T N)$. We shall show that

$$
h(d \tilde{X}, d \tilde{X})=d[\tilde{Z}, \tilde{Z}]=|d f|^{2}(X) d t
$$

where $\tilde{Z}=\int_{\tilde{U}} \vartheta$. Note that, by assumption, $Z=\int_{U} \vartheta$ is a Brownian motion on $\mathbb{R}^{n}$ where $n=\operatorname{dim} M$. Furthermore, we have $d X=\sum_{i} U e_{i} \circ d Z^{i}$ and $d \tilde{X}=\sum_{i}(d f)_{X} U e_{i} \circ d Z^{i}$. Hence we obtain

$$
\begin{aligned}
h(d \tilde{X}, d \tilde{X}) & =\left(f^{*} h\right)(d X, d X) \\
& =\sum_{i, j}\left(f^{*} h\right)_{X}\left(U e_{i}, U e_{j}\right) d Z^{i} d Z^{j} \\
& =\sum_{i}\left(f^{*} h\right)_{X}\left(U e_{i}, U e_{i}\right) d t \\
& =\sum_{i} h_{f(X)}\left((d f)_{X} U e_{i},(d f)_{X} U e_{i}\right) d t=|d f|^{2}(X) d t
\end{aligned}
$$

as well as $d[\tilde{Z}, \tilde{Z}]=\left\langle\tilde{U}^{-1} d \tilde{X}, \tilde{U}^{-1} d \tilde{X}\right\rangle=h(d \tilde{X}, d \tilde{X})$.
THEOREM 1.7.15 (Stochastic characterization of affine and harmonic maps). Let $M$ and $N$ be smooth manifolds and $f: M \rightarrow N$ be a differentiable map.
(i) Let $\nabla^{M}$ on $M$, respectively $\nabla^{N}$ on $N$, be torsion-free linear connections. Then $f$ is affine if and only if $f$ maps $\nabla^{M}$-martingales on $M$ to $\nabla^{N}$-martingales on $N$.
(ii) Let $(M, g)$ and $(N, h)$ be Riemannian manifolds and $\nabla^{M}$, respectively $\nabla^{N}$ the corresponding Levi-Civita connections. Then $f$ is harmonic if and only if $f$ maps $\operatorname{BM}(M, g)$ to the class $\operatorname{Mart}(N, h)$ of $\nabla^{N}$-martingales on $N$.

Proof. (i) By Theorem 1.6.48, $X$ is a $\nabla^{M}$-martingale on $M$ if and only if $\mathscr{A}(X)$ is a local martingale. Hence, by Theorem 1.7.12, $\tilde{X}=f(X)$ is a $\nabla^{N}$-martingale on $N$ (equivalently $\mathscr{A}(\tilde{X})$ a local martingale) for each $\nabla^{M}$-martingale $X$ if and only if $\nabla d f(d X, d X)=0$ for each $\nabla^{M}$-martingale $X$ which in turn is equivalent to $\nabla d f=0$.

Indeed, by stochastic development, $\nabla^{M}$-martingales $X$ on $M$ are of the form $d X=$ $\sum_{i} U e_{i} \circ d Z^{i}$ for some local martingale $Z$ in $\mathbb{R}^{n}(n=\operatorname{dim} M)$. Taking $Z=(B, 0, \ldots, 0)$ where $B$ a real Brownian motion thus gives $\nabla d f(d X, d X)=\nabla d f\left(U e_{1}, U e_{1},\right) d t=0$. For a constant starting point $X_{0}=x \in M$, the frame $U_{0}=u \in \mathrm{~L}(T M)$ with $\pi(u)=x$ can be chosen arbitrarily, so that necessarily $\nabla d f=0$ must hold.
(ii) According to formula (1.7.6), $f$ maps $\operatorname{BM}(M, g)$ to $\operatorname{Mart}(N, h)$ if and only if $\tau(f)(X)=0$ along each Brownian motion $X$ on $M$. Since the starting point of $X$ can be chosen arbitrarily, this however means $\tau(f)=0$.

The proof of Theorem 1.7 .15 shows that a map $f$ is already affine if it transfers onedimensional martingales on $M$ to one-dimensional martingales on $N$.

Corollary 1.7.16. Let $M, N$ be two differentiable manifolds, endowed with torsionfree linear connections, and $f: M \rightarrow N$ a differentiable map. The following items are equivalent:
(i) $f$ is affine;
(ii) $f$ maps one-dimensional martingales to one-dimensional martingales;
(iii) f maps one-dimensional Brownian motions to one-dimensional Brownian motions.

Proof. By Corollary 1.7.10, $f$ is affine if and only if $f$ maps geodesics to geodesics. On the other hand, by Theorem 1.6.53, one-dimensional martingales and one-dimensional Brownian motions move on geodesics. Affine maps $f$ hence transfer one-dimensional martingales (resp. one-dimensional Brownian motions) to one-dimensional martingales (resp. one-dimensional Brownian motions). Conversely, if $f$ is a differentiable map with this property, then for each geodesic curve $\gamma$ on $M$, the composition $f \circ \gamma$ maps continuous real local martingales (resp. real Brownian motions) to $\operatorname{Mart}(N, h)$; for each geodesic curve $\gamma$ on $M$, by Theorem 1.7.15, the composition $f \circ \gamma$ is thus affine ( $\equiv$ harmonic), and hence $f$ itself affine.

DEFINITION 1.7.17 (Horizontally conformal map, harmonic morphism). Let ( $M, g$ ) and $(N, h)$ be Riemannian manifolds and $f: M \rightarrow N$ a differentiable map. Then $f$ is said to be horizontally conformal, if
(a) at each point $x \in M$ at which $(d f)_{x} \neq 0$ the linear map $(d f)_{x}: T_{x} M \rightarrow T_{f(x)} N$ is surjective;
(b) there exists a function $\lambda: M \rightarrow \mathbb{R}_{+}$such that, for all $v, w \in\left(\operatorname{ker}(d f)_{x}\right)^{\perp}$,

$$
h_{f(x)}\left(f_{*} v, f_{*} w\right)=\lambda^{2}(x) g_{x}(v, w)
$$

The function $\lambda: M \rightarrow \mathbb{R}_{+}$is called dilatation of $f$ where $\lambda(x):=0$ if $(d f)_{x}=0$. A map $f: M \rightarrow N$ is called harmonic morphism (with dilatation $\lambda$ ) if $f$ is harmonic and horizontally conformal (with dilatation $\lambda$ ).

Lemma 1.7.18. Let $f:(M, g) \rightarrow(N, h)$ be a differentiable map between Riemannian manifolds and $\lambda: M \rightarrow \mathbb{R}_{+}$a function. The following items are equivalent:
(i) $f$ is horizontally conformal with dilatation $\lambda$;
(ii) $d f \circ(d f)^{\text {ad }}=\lambda^{2}$ id $\mid f^{*} T N$ where $(d f)^{\text {ad }}: f^{*} T N \rightarrow T M$ is the homomorphism of vector bundles fiberwise adjoint to $d f$;
(iii) $g(\operatorname{grad}(\varphi \circ f), \operatorname{grad}(\psi \circ f))=\lambda^{2} h(\operatorname{grad} \varphi, \operatorname{grad} \psi) \circ$ f for all $\varphi, \psi \in C^{\infty}(N)$.

Then necessarily $\lambda^{2}=|d f|_{\mathrm{op}}^{2}$ where $|d f|_{\mathrm{op}}$ is the operator norm of $d f$. Note that $k|d f|_{\mathrm{op}}^{2}=$ $|d f|^{2}$ where $k=\operatorname{dim} N$.

Proof. (i) $\Leftrightarrow$ (ii) The adjoint $(d f)_{x}^{\text {ad }}: T_{f(x)} N \rightarrow T_{x} M$ to $(d f)_{x}$ is determined by

$$
h_{f(x)}\left((d f)_{x} v, u\right)=g_{x}\left(v,(d f)_{x}^{\mathrm{ad}} u\right), \quad v \in T_{x} M, u \in T_{f(x)} N,
$$

and $f$ is hence horizontally conformal if and only if for all $x \in M$,

$$
(d f)_{x} \circ(d f)_{x}^{\mathrm{ad}}=\lambda^{2}(x) \operatorname{id} \mid T_{f(x)} N
$$

(ii) $\Leftrightarrow$ (iii) Since $g\left(A,(d f)^{\operatorname{ad}}\left(f^{*} \operatorname{grad} \varphi\right)\right)=A(\varphi \circ f)$ for $A \in \Gamma(T M)$, we have $(d f)^{\operatorname{ad}}\left(f^{*} \operatorname{grad} \varphi\right)=\operatorname{grad}(\varphi \circ f)$ for $\varphi \in C^{\infty}(N)$ from where the equivalence follows. The additional claim is obvious.

THEOREM 1.7.19 (Analytic characterization of harmonic morphisms). Let $f:(M, g) \rightarrow$ $(N, h)$ be a differentiable map between Riemannian manifolds and let $\lambda: M \rightarrow \mathbb{R}_{+}$be a function. The following conditions are equivalent:
(i) $f$ is a harmonic morphism (with dilatation $\lambda$ );
(ii) $\Delta_{M}(\varphi \circ f)=\lambda^{2}\left(\Delta_{N} \varphi \circ f\right)$ for $\varphi \in C^{\infty}(N)$.

Proof. (i) $\Rightarrow$ (ii) Since $f$ is harmonic, by composition formula (1.7.3) it holds

$$
\Delta_{M}(\varphi \circ f)=\varphi_{*} \tau(f)+\operatorname{trace}\left(f^{*} \nabla d \varphi\right)=\operatorname{trace}\left(f^{*} \nabla d \varphi\right)
$$

We have to show that $\Delta_{M}(\varphi \circ f)(x)=\lambda^{2}(x)\left(\Delta_{N} \varphi \circ f\right)(x)$ for $x \in M$. To this end, without restrictions, we may assume that $(d f)_{x} \neq 0$. If then $\left(a_{1}, \ldots, a_{\ell}\right)$ is an orthonormal basis of $\left(\operatorname{ker}(d f)_{x}\right)^{\perp}$, then by the horizontal conformality of $f$

$$
\left(\frac{1}{\lambda(x)}(d f)_{x} a_{i}: 1 \leq i \leq \ell\right)
$$

defines an orthonormal basis of $T_{f(x)} N$, and hence

$$
\Delta_{M}(\varphi \circ f)(x)=\sum_{i=1}^{\ell}(\nabla d \varphi)\left(f_{*} a_{i}, f_{*} a_{i}\right)=\lambda^{2}(x)\left(\Delta_{N} \varphi \circ f\right)(x)
$$

(ii) $\Rightarrow$ (i) For $\varphi, \psi \in C^{\infty}(N)$ we have on one hand

$$
\begin{aligned}
\Delta_{M}((\varphi \psi) \circ f)= & (\varphi \circ f) \Delta_{M}(\psi \circ f)+(\psi \circ f) \Delta_{M}(\varphi \circ f) \\
& +2 g(\operatorname{grad}(\varphi \circ f), \operatorname{grad}(\psi \circ f)),
\end{aligned}
$$

on the other hands it holds $\Delta_{N}(\varphi \psi)=\varphi \Delta_{N} \psi+\psi \Delta_{N} \varphi+2 h(\operatorname{grad} \varphi, \operatorname{grad} \psi)$. Composing the last equation with $f$ and multiplying by $\lambda^{2}$, then subtraction from the first equation gives

$$
g(\operatorname{grad}(\varphi \circ f), \operatorname{grad}(\psi \circ f))=\lambda^{2} h(\operatorname{grad} \varphi, \operatorname{grad} \psi) \circ f
$$

which shows that $f$ is horizontally conformal. It remains to verify $\tau(f)=0$. To this end, we conclude again as above from horizontal conformality of $f$ that for $\varphi \in C^{\infty}(N)$

$$
\operatorname{trace}\left(f^{*} \nabla d \varphi\right)=\lambda^{2}\left(\Delta_{N} \varphi \circ f\right)
$$

But since $\varphi_{*} \tau(f)+\operatorname{trace}\left(f^{*} \nabla d \varphi\right)=\Delta_{M}(\varphi \circ f)=\lambda^{2} \cdot\left(\Delta_{N} \varphi \circ f\right)$, we have $\varphi_{*} \tau(f)=0$ for any $\varphi \in C^{\infty}(N)$, and thus $\tau(f)=0$.

THEOREM 1.7.20. Let $f:(M, g) \rightarrow(N, h)$ be a differentiable map between Riemannian manifolds. The following conditions are equivalent:
(i) $f$ is a harmonic morphism (with dilatation $\lambda$ );
(ii) $f$ maps $\mathrm{BM}(M, g)$ to $\mathrm{BM}(N, h)$ modulo time change, more precisely: to each Brownian motion $X$ on $(M, g)$ there exists a Brownian motion $\tilde{X}$ on $(N, h)$ such that $f\left(X_{t}\right)=\tilde{X}_{T_{t}}$ a.s. where $T_{t}=\int_{0}^{t} \lambda^{2}\left(X_{s}\right) d s$.

REMARK 1.7.21. Note that the $N$-valued Brownian motion $\tilde{X}$ in (ii) is determined through the condition $f\left(X_{t}\right)=\tilde{X}_{T_{t}}$ only up to time $T_{\infty}$; it may however always be extended to maximal lifetime by "piecing on" an independent Brownian motion: the antidevelopment $\int_{\tilde{U}} \vartheta$ of $\tilde{X}$ gives first a stopped Brownian motion on $\mathbb{R}^{\operatorname{dim} N}$ which can be extended to all of $\mathbb{R}_{+}$. Stochastic development of this Brownian motion then gives the wanted prolongation of $\tilde{X}$. In this case the equality $f\left(X_{t}\right)=\tilde{X}_{T_{t}}$ then holds on an enlarged probability space.

Proof of Theorem 1.7.20. By Theorem 1.7 .15 (ii) the map $f$ is harmonic if and only if for each Brownian motion $X$ on $M$, the target process $f(X)$ is a $\nabla$-martingale on $N$ (with respect to the Levi-Civita connection $\nabla$ ) which according to Theorem 1.6.48 means that all anti-developments $\tilde{Z}$ of $f(X)$ in $\mathbb{R}^{\operatorname{dim} N}$ are local martingales. Since, modulo time change, $f(X)$ is a $\operatorname{BM}(N, h)$ if and only if $\tilde{Z}$ is a $\operatorname{BM}\left(\mathbb{R}^{\operatorname{dim} N}\right)$, it remains to show that if $f$ is in addition horizontally conformal then all anti-developments $\tilde{Z}$ of $f(X)$ are Brownian motions on $\mathbb{R}^{\operatorname{dim} N}$ modulo time change. Using the notations of Remark 1.7.13, we have

$$
\begin{aligned}
d \tilde{Z}^{k} d \tilde{Z}^{\ell} & =\sum_{i}\left\langle\tilde{U}^{-1}(d f)_{X} U e_{i}, e_{k}\right\rangle\left\langle\tilde{U}^{-1}(d f)_{X} U e_{i}, e_{\ell}\right\rangle d t \\
& =\sum_{i} g\left(U e_{i},(d f)^{\operatorname{ad}} \tilde{U} e_{k}\right) g\left(U e_{i},(d f)^{\mathrm{ad}} \tilde{U} e_{\ell}\right) d t \\
& =g\left((d f)^{\operatorname{ad}} \tilde{U} e_{k},(d f)^{\operatorname{ad}} \tilde{U} e_{\ell}\right) d t
\end{aligned}
$$

Hence it remains to observe that $g\left((d f)^{\text {ad }} \tilde{U} e_{k},(d f)^{\text {ad }} \tilde{U} e_{\ell}\right) d t=\lambda^{2}(X) \delta_{k \ell} d t$ holds for all $k, \ell$ and all horizontal lifts $\tilde{U}$ of semimartingales of the form $\tilde{X}=f(X)$ with $X$ in $\operatorname{BM}(M, g)$ if and only if $f$ is horizontally conformal with dilatation $\lambda$.

Theorem 1.7.15 (i) says in particular that the composition $\varphi(X)$ of an $M$-valued martingale $X$ with an affine function $\varphi \in C^{\infty}(M)$ gives a real local martingale. However, to use affine functions as "martingale testers" and to characterize $M$-valued martingales by this property usually fails due to the lacking richness of affine functions: in general, nonconstant real-valued affine functions may even not exist locally. A suitable substitute for affine functions are convex functions. There are typically also obstructions of topological and geometric nature for existence of globally defined non-trivial convex functions, but locally convex functions provide a rich class of functions.

Lemma 1.7.22. Let $M$ be a manifold endowed with a torsion-free linear connection. To each $x \in M$ there exists an open neighbourhood $U$ and a strictly convex function $\varphi \in C^{\infty}(U)$ with prescribed 2-jet, i.e., given $a \in \mathbb{R}, b \in T_{x}^{*} M$ and $C \in T_{x}^{*} M \otimes T_{x}^{*} M$ positive definite, there is an open neighbourhood $U$ of $x$ and a function $\varphi \in C^{\infty}(U)$ such that $\varphi(x)=a,(d \varphi)_{x}=b,(\nabla d \varphi)_{x}=C$ and $\nabla d \varphi>0$ on $U$.

Proof. We choose normal coordinates $h$ about $x$ as follows. The exponential map

$$
\exp _{x}:\left(T_{x} M, 0\right) \rightarrow(M, x), \quad v \mapsto \gamma_{v}(1)
$$

where $\gamma_{v}$ is the geodesic curve determined by $\gamma_{v}(0)=x$ and $\dot{\gamma}_{v}(0)=v$, is well-defined locally about 0 and has full rank at 0 , as can be seen from

$$
\left(d \exp _{x}\right)_{0} v=\left.\frac{d}{d t}\right|_{t=0} \exp _{x}(t v)=\left.\frac{d}{d t}\right|_{t=0} \gamma_{t v}(1)=\left.\frac{d}{d t}\right|_{t=0} \gamma_{v}(t)=v
$$

Fixing a linear isomorphism $\iota: \mathbb{R}^{n} \xrightarrow{\sim} T_{x} M$, then by the local inverse theorem, $\exp _{x} \circ \iota$ maps an open $\varepsilon$-ball $V_{\varepsilon}$ in $\mathbb{R}^{n}$ about 0 diffeomorphically onto an open neighbourhood of $x$
in $M$ and $h:=\left(\exp _{x} \circ \iota \mid V_{\varepsilon}\right)^{-1}$ defines a local chart at $x$. Now let $b=\sum_{i} b_{i}\left(d h^{i}\right)_{x}$ and $C=\sum_{i, j} C_{i j}\left(d h^{i} \otimes d h^{j}\right)_{x}$, and define

$$
\varphi=a+\sum_{i} b_{i} h^{i}+\sum_{i, j} C_{i j} h^{i} h^{j}
$$

then $\varphi \mid U$ has the wanted properties for some sufficiently small open neighbourhood $U$ of $x$. Indeed, letting $\partial_{i}=\frac{\partial}{\partial h^{i}}$ and $(\nabla d \varphi)_{i j}=(\nabla d \varphi)\left(\partial_{i}, \partial_{j}\right)$, we have

$$
\nabla d \varphi \mid U=\sum_{i, j}(\nabla d \varphi)_{i j} d h^{i} \otimes d h^{j}=\sum_{i, j}\left(\partial_{i} \partial_{j} \varphi-\sum_{k} \Gamma_{i j}^{k} \partial_{k} \varphi\right) d h^{i} \otimes d h^{j}
$$

Note that, by construction of the chart, $\Gamma_{i j}^{k}(x)=0$ which can be seen as follows: letting $\gamma_{v}(t)=\exp _{x}(t v)$ be again the geodesic curve defined locally about $t=0$ and determined by the properties $\gamma_{v}(0)=x$ and $\dot{\gamma}_{v}(0)=v$, we have for any $v \in T_{x} M$,

$$
0=\left.\frac{d^{2}}{d t^{2}}\right|_{t=0} \gamma_{v}^{k}(t)=\sum_{i, j} \Gamma_{i j}^{k}(x) v^{i} v^{j}
$$

which implies $\Gamma_{i j}^{k}(x)=0$ as $\nabla$ is torsion-free.
As a result of the richness of germs of convex functions guaranteed by Lemma 1.7.22, affine and harmonic maps can be characterized through their functional behaviour under pullback.

THEOREM 1.7.23 (Pullback properties of affine/harmonic maps).
(i) Let $f: M \rightarrow N$ be a differentiable map between manifolds equipped with torsionfree linear connections. The following items are equivalent:
(a) $f$ is affine;
(b) pullbacks $f^{*} \varphi$ of germs of convex functions on $N$ are convex, i.e., for each convex function $\varphi \in C^{\infty}(V)$ defined on an open subset $V$ of $N$, the composition $\varphi \circ f$ is convex on $f^{-1} V$.
(ii) Let $f:(M, g) \rightarrow(N, h)$ be a differentiable map between Riemannian manifolds. The following items are equivalent:
(a) $f$ is harmonic;
(b) pullbacks $f^{*} \varphi$ of germs of convex functions on $N$ are subharmonic, i.e., for each convex function $\varphi \in C^{\infty}(V)$ defined on an open subset $V$ of $N$, the composition $\varphi \circ f$ is subharmonic on $f^{-1} V$.

Proof. The implications (a) $\Rightarrow$ (b) are each direct consequences of the composition formulas (1.7.2) and (1.7.3)

$$
\nabla d(\varphi \circ f)=\varphi_{*} \nabla d f+f^{*} \nabla d \varphi \quad \text { resp. } \quad \Delta(\varphi \circ f)=\varphi_{*} \tau(f)+\operatorname{trace}\left(f^{*} \nabla d \varphi\right) .
$$

The implications (b) $\Rightarrow$ (a) rely on the richness of germs of convex functions as formulated in Lemma 1.7.22. For instance, as in part (i), whenever $\nabla d(\varphi \circ f) \geq 0$ holds if $\nabla d \varphi \geq 0$, then already $\nabla d f=0$ must be satisfied, otherwise there would exist $x \in M$ and $v \in$ $T_{x} M$ such that $w:=(\nabla d f)_{x}(v, v) \neq 0$ in $T_{f(x)} N$. By Lemma 1.7.22, there is then an open neighbourhood $V$ of $f(x)$ in $N$ and a convex function $\varphi \in C^{\infty}(V)$ such that $(d \varphi)_{f(x)} w<-\left|(d f)_{x} v\right|^{2}$ and $(\nabla d \varphi)_{f(x)}=h_{f(x)}$. But this would imply

$$
(\nabla d(\varphi \circ f))_{x}(v, v)=(d \varphi)_{f(x)} w+\left|(d f)_{x} v\right|^{2}<0
$$

in contradiction to $\nabla d(\varphi \circ f) \geq 0$. The implication (b) $\Rightarrow$ (a) in (ii) can be shown analogously.

In the stochastic context the richness of germs of convex functions allows a characterization of martingales, which has first been used by Darling [6] for the definition of $\nabla$-martingales.

THEOREM 1.7.24 (Darling's characterization of $\nabla$-martingales). Let $M$ be a differentiable manifold, $\nabla$ a torsion-free linear connection on $M$ and $X$ an $M$-valued semimartingale. Then $X$ is a $\nabla$-martingale if and only iffor each $\varphi \in C^{\infty}(M)$ and each open $V \subset M$ such that $\varphi \mid V$ is convex, the following holds true: If

$$
\varphi(X)=\varphi\left(X_{0}\right)+N+A
$$

is the Doob-Meyer decomposition of the real semimartingale $\varphi(X)$ and if $\sigma, \tau$ are stopping times such that $\sigma \leq \tau$ and $X \mid[\sigma, \tau[$ takes values in $V$, then the process $A$ is monotonically increasing on $[\sigma, \tau[$ a.s.

Proof. By the Geometric Itô formula (Theorem 1.6.45) and the notations there, we have for each $\varphi \in C^{\infty}(M)$ the formula

$$
d(\varphi(X))=\sum_{i}(d \varphi)(X)\left(U e_{i}\right) d Z^{i}+\frac{1}{2}(\nabla d \varphi)(d X, d X)
$$

Denoting by $Z=Z^{\text {Mart }}+Z^{\text {drift }}$ the Doob-Meyer decomposition of the $\mathbb{R}^{n}$-valued semimartingale $Z$, we obtain for the "drift part" $A$ of $\varphi(X)$ the representation

$$
\begin{equation*}
d A=\sum_{i}(d \varphi)(X)\left(U e_{i}\right) d\left(Z^{\mathrm{drift}}\right)^{i}+\frac{1}{2}(\nabla d \varphi)(d X, d X) \tag{1.7.8}
\end{equation*}
$$

According to Theorem 1.6 .48 (i), the process $X$ is a $\nabla$-martingale on $M$ if and only if $Z^{\text {drift }} \equiv 0$ modulo indistinguishability. Hence necessity of the given condition is obvious. Recall that $1_{[\sigma, \tau[ }(\nabla d \varphi)(d X, d X)$ is the differential of an increasing process. This is an immediate consequence of the definition of the $b$-quadratic variation, e.g. formula (1.3.3), since $X \mid\left[\sigma, \tau\left[\right.\right.$ takes values in $V$ and $(\nabla d \varphi)_{x}$ is positive semidefinite for $x \in V$.

Conversely, suppose now that for each $\varphi \in C^{\infty}(M)$ and each open subset $V \subset$ $M$ the following condition holds: If $\varphi \mid V$ is convex and $X \mid[\sigma, \tau[$ takes values in $V$ then $A \mid\left[\sigma, \tau\left[\right.\right.$ is almost surely monotonically increasing. We have to show that $Z^{\text {drift }} \equiv 0$ under this condition. By means of Lemma 1.3.1 the claim can be localized in space, and without restriction we may assume that $X$ takes its values in a fixed relatively compact open subset $V$ whose closure $\bar{V}$ lies completely in the domain of a chart $h$ for $M$. We fix a positive definite section $g$ of $T^{*} M \otimes T^{*} M$ over $V$, for instance, $g=\sum_{i} d h^{i} \otimes d h^{i}$, and are going to show that for each $f \in C^{\infty}(M)$ and $\varepsilon>0$, the process

$$
\begin{equation*}
\int \sum_{i}(d f)_{X}\left(U e_{i}\right) d\left(Z^{\mathrm{drift}}\right)^{i}+\frac{1}{2} \varepsilon \int g(d X, d X) \tag{1.7.9}
\end{equation*}
$$

is almost surely isotone. This then gives immediately the claim, since with $\varepsilon \downarrow 0$ in (1.7.9) one obtains that

$$
\int(d f)_{X}\left(U d Z^{\text {drift }}\right) \equiv \int \sum_{i}(d f)_{X}\left(U e_{i}\right) d\left(Z^{\text {drift }}\right)^{i}
$$

is almost surely isotone. Passing from $f$ to $-f$ thus shows that $\int(d f)(X)\left(U d Z^{\text {drift }}\right)$ is almost surely constant, and since this holds for all $f \in C^{\infty}(M)$, we conclude $Z^{\text {drift }} \equiv 0$ modulo indistinguishability.

In the sequel let $f \in C^{\infty}(M)$ and $g$ be a positively definite section of $T^{*} M \otimes T^{*} M$ over $V$; it remains to show that the process

$$
\begin{equation*}
N:=\int(d f)_{X}\left(U d Z^{\text {drift }}\right)+\frac{1}{2} \int g(d X, d X) \tag{1.7.10}
\end{equation*}
$$

is almost surely isotone. To this end, we construct a family $\left(N^{\delta}\right)_{\delta>0}$ of isotone processes with the property that $N^{\delta} \rightarrow N$ almost surely as $\delta \downarrow 0$, uniformly on compact time intervals of the form $[0, t]$. At this place the local richness of convex functions comes into effect, as by Lemma 1.7.22, to each point $a \in V$ there exists an open neighbourhood $V_{a}$ of $a$ and a strictly convex function $\varphi^{a}$ on $V_{a}$ such that

$$
\varphi^{a}(a)=0, \quad\left(d \varphi^{a}\right)_{a}=(d f)_{a}, \quad\left(\nabla d \varphi^{a}\right)_{a}=g_{a}
$$

and such that in addition, for fixed $\delta>0$, possibly after shrinking of $V_{a}$, it holds that

$$
\begin{align*}
& \sup _{x \in V_{a}}\left|d\left(\varphi^{a} \circ h^{-1}\right)-d\left(f \circ h^{-1}\right)\right|(h(x)) \leq \delta,  \tag{1.7.11}\\
& \sup _{x \in V_{a}}\left|\left(\nabla d \varphi^{a}\right)_{x}\left(\partial_{i}, \partial_{j}\right)-g_{x}\left(\partial_{i}, \partial_{j}\right)\right| \leq \delta
\end{align*}
$$

where $\partial_{i}=\frac{\partial}{\partial h^{i}}$ with respect to a fixed chart $(h, V)$. For a given $\delta>0$ then $V$ is already covered by finitely many $V_{a}$ 's, and according to Lemma 1.3.1 we can find a sequence $\left(\tau_{n}\right)_{n \geq 0}$ of stopping times such that

$$
0=\tau_{0} \leq \tau_{1} \leq \tau_{2} \leq \ldots \quad \text { and } \quad \sup _{n} \tau_{n}=\infty
$$

and such that on each interval $\left[\tau_{n}, \tau_{n+1}[\right.$ the process $X$ takes values only in one (of the finitely many) $V_{a(n)}$. Therewith we finally define the process

$$
N^{\delta}=\int \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left(d \varphi^{a(n)}\right)(X)\left(U d Z^{\mathrm{drift}}\right)+\frac{1}{2} \int \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left(\nabla d \varphi^{a(n)}\right)(d X, d X)
$$

which is almost surely isotone, since by construction $N \delta$ satisfies monotonicity on each subinterval $\left[\tau_{n}, \tau_{n+1}[\right.$. We want to verify the convergence $\delta \rightarrow 0$ almost surely and uniformly $[0, t]$ as $\delta \downarrow 0$. But now we have

$$
\begin{align*}
N_{t}^{\delta}-N_{t} & =\int_{0}^{t} \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left(d \varphi^{a(n)}-d f\right)(X)\left(U d Z^{\mathrm{drift}}\right) \\
& +\frac{1}{2} \int_{0}^{t} \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left(\nabla d \varphi^{a(n)}-g\right)(d X, d X) \tag{1.7.12}
\end{align*}
$$

and (1.7.11) can be used to estimate. For the first term we have the estimate

$$
\begin{aligned}
& \left|\int_{0}^{t} \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left(d \varphi^{a(n)}-d f\right)(X)\left(U d Z^{\mathrm{drift}}\right)\right| \\
& =\left|\int_{0}^{t} \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left[d\left(\varphi^{a(n)} \circ h^{-1}\right)-d\left(f \circ h^{-1}\right)\right](d h)(X)\left(U d Z^{\mathrm{drift}}\right)\right| \\
& \quad \leq \delta \sum_{i} \sup _{[0, t]}\left|(d h)(X)\left(U e_{i}\right)\right| \int_{0}^{t}\left|d\left(Z^{\mathrm{drift}}\right)^{i}\right| .
\end{aligned}
$$

In a similar way, letting $X^{i}=h^{i}(X)$, again with (1.7.11), we obtain for the second term in (1.7.12) the estimate

$$
\begin{aligned}
& \left|\int_{0}^{t} \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left(\nabla d \varphi^{a(n)}-\nabla d f\right)(d X, d X)\right| \\
& \quad=\left|\sum_{i, j} \int_{0}^{t} \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left[\nabla d \varphi^{a(n)}\left(\partial_{i}, \partial_{j}\right)-\nabla d f\left(\partial_{i}, \partial_{j}\right)\right] d\left[X^{i}, X^{j}\right]\right| \\
& \quad \leq \delta \sum_{i, j} \int_{0}^{t}\left|d\left[X^{i}, X^{j}\right]\right|
\end{aligned}
$$

which completes the proof.

### 1.8. Convergence and Confluence of Martingales

In this Section we want to elaborate and develop further the theory of martingales on manifolds. One of the difficulties of the theory relies in the fact that on manifolds there is no counterpart of the linear concept of taking conditional expectations. This apparent drawback is due to the nature of the subject and makes martingales on manifolds to an interesting non-linear instrument. The close connection between the behavior of martingales on manifolds and questions of convex geometry will quickly become apparent, for instance, questions of approximability of the Riemannian distance function on a manifold by convex functions. Such questions are known to be closely linked to the curvature of a Riemannian manifold.

One of basic tools of scalar martingale theory is the martingale convergence theorem which guarantees, for instance, that bounded martingales on $\mathbb{R}^{n}$ converge, i.e. have an almost sure limit as $t \rightarrow \infty$. In this form the convergence theorem obviously does not carry over to manifolds, as martingales on $M$ taking values in a relatively compact subset do not need to converge which can already seen from simple examples, like Brownian motions on compact Riemannian manifolds or one-dimensional Brownian motions $X=\gamma(B)$ where $\gamma$ is a closed geodesic curve.

As well-known [38], for real-valued continuous local martingale $X$, the following sets coincide modulo nullsets:

$$
\left\{\lim _{t \rightarrow \infty} X_{t} \text { exists in } \mathbb{R}\right\}, \quad\left\{[X, X]_{\infty}<\infty\right\}, \quad\left\{\sup _{t \in \mathbb{R}_{+}} X_{t}<\infty\right\}, \quad\left\{\inf _{t \in \mathbb{R}_{+}} X_{t}>-\infty\right\}
$$

On the other hand, the concept of quadratic variation of $M$-valued semimartingales provides a notion to quantify the "oscillation" of $M$-valued martingales. Since each $M$-valued martingale $X$ comes via stochastic development from an $\mathbb{R}^{n}$-valued local martingale $Z$ and since for Riemannian manifolds the Riemannian quadratic variation of $X$ coincides with the quadratic variation of $Z$,

$$
g(d X, d X)=\sum_{i, j} g\left(U e_{i}, U e_{j}\right) d Z^{i} d Z^{j}=\sum_{i} d\left[Z^{i}, Z^{i}\right]=d[Z, Z]
$$

it is not surprising that convergence of martingales on manifold can be expressed in terms of finiteness of the quadratic variation as $t \rightarrow \infty$.

Before entering into details we collect some notations.
Notation 1.8.1. For an adapted continuous process $A$ which is pathwise locally of bounded variation, we call $V_{A}=\int|d A|$ the variation process of $A$. For an $M$-valued continuous semimartingale $X$ and a bilinear form $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$ let $\int b(d X, d X)$
be the $b$-quadratic variation and $\int|b(d X, d X)|$ its variation process. For $b, g \in \Gamma\left(T^{*} M \otimes\right.$ $T^{*} M$ ), we write $b \leq g$ if $g-b$ is positive semidefinite which means that the bilinear form $(g-b)_{x} \in T_{x}^{*} M \otimes T_{x}^{*} M$ is positive semidefinite for each $x \in M$.

REMARK 1.8.2. Let $X$ be an $M$-valued semimartingale. If $S, T$ are two $\mathbb{R}_{+}$-valued random variables (not necessarily stopping times) with the property that $X \mid[S, T$ [ takes its values in an open set $U$ in $M$, and if $b \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$ such that $b \geq 0$ on $U$ (i.e. $b_{x}$ positive semidefinite for $x \in U$ ), then the process $\int b(d X, d X)$ is almost surely isotone on $\left[S, T\left[\right.\right.$. If in addition $-g \leq b \leq g$ on $U$ where $g \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$, then $\int|b(d X, d X)| \leq \int g(d X, d X)$ a.s. on $[S, T[$.

LEMMA 1.8.3. Let $M$ be a manifold and $\nabla$ a torsion-free linear connection on $M$. Every point in $M$ has an open neighbourhood $U$ such that each $M$-valued $\nabla$-martingale $X$ converges almost surely on the set

$$
\begin{aligned}
\Omega_{0}: & =\left\{X_{t} \in U \text { eventually }\right\} \\
& \equiv\left\{\omega \in \Omega: \exists t(\omega) \in \mathbb{R}_{+} \text {such that } X_{s}(\omega) \in U \text { for all } s \geq t(\omega)\right\}
\end{aligned}
$$

Proof. For $x \in M$ we choose a sufficiently small open neighbourhood $U$ of $x$ such that by Lemma 1.7.22 we can find a bounded function $\varphi=\left(\varphi^{1}, \ldots, \varphi^{n}\right) \in C^{\infty}\left(M ; \mathbb{R}^{n}\right)$ with the following properties:
(a) $\varphi^{i} \mid U$ is convex for $1 \leq i \leq n=\operatorname{dim} M$,
(b) $(\varphi \mid U, U)$ defines a chart for $M$ about $x$.

Since $X$ is a $\nabla$-martingale, by the Geometric Itô formula (1.6.32), we get

$$
\varphi^{i}(X)-\varphi^{i}\left(X_{0}\right)=M^{i}+A^{i}
$$

where $M^{i} \in \mathscr{M}$ and $d A^{i}=\frac{1}{2} \nabla d \varphi^{i}(d X, d X)$. By construction $\varphi^{i} \mid U$ is convex, and hence the process $A^{i}$ is almost surely eventually isotone on $\Omega_{0}$, and in particular pathwise bounded from below on $\Omega_{0}$. Since the functions $\varphi^{i}$ are bounded, we observe that for each index $i$ the local martingale

$$
M^{i}=\varphi^{i}(X)-\varphi^{i}\left(X_{0}\right)-A^{i}
$$

is almost surely pathwise bounded from above on $\Omega_{0}$, and hence convergent on $\Omega_{0}$. Conversely this shows however that each $A^{i}$ is actually almost surely bounded on $\Omega_{0}$ and (since eventually isotone on $\Omega_{0}$ ) also convergent.

REMARK 1.8.4. The proof above actually shows that $\varphi^{i}(X)$ is even a semimartingale up to $\infty$ on the set $\Omega_{0}=\left\{X_{t} \in U\right.$ eventually $\}$, i.e., if $\varphi^{i}(X)=\varphi^{i}\left(X_{0}\right)+M^{i}+A^{i}$ denotes the Doob-Meyer decomposition of $\varphi^{i}(X)$, then $M_{\infty}^{i}$ and $A_{\infty}^{i}$ exist almost surely on $\Omega_{0}$, and for almost all $\omega \in \Omega_{0}$, the map $[0, \infty] \rightarrow \mathbb{R}, t \mapsto A_{t}^{i}(\omega)$, is of bounded variation. The last claim comes from the fact that on $\Omega_{0}$, for sufficiently large $s$, it holds:

$$
\int_{0}^{\infty}\left|d A^{i}\right|=\int_{0}^{s}\left|d A^{i}\right|+A_{\infty}^{i}-A_{s}^{i} .
$$

We want to explain the notion of a semimartingale up to $\infty$ also for $M$-valued semimartingales.

DEFINITION 1.8.5 ( $M$-valued semimartingale up to $\infty$ ). Let $\Omega_{0} \subset \Omega$ be a measurable subset. An $M$-valued semimartingale $X$ is called a semimartingale up to $\infty$ on $\Omega_{0}$ if for any $\varphi \in C^{\infty}(M)$ the composition $\varphi(X)$ is a semimartingale up to $\infty$ on $\Omega_{0}$.

REMARK 1.8.6. A real-valued semimartingale $Y$ is obviously a semimartingale up to $\infty$ on all of $\Omega$ if and only if $\lim _{t \rightarrow \infty} Y_{t}$ exists almost surely and the time-changed process $\tilde{Y}$,

$$
\tilde{Y}_{t}:= \begin{cases}Y_{t /(1-t)} & \text { for } 0 \leq t<1 \\ Y_{\infty} & \text { for } t \geq 1\end{cases}
$$

is a semimartingale (with respect to the filtration $\tilde{\mathscr{F}}_{t}=\mathscr{F}_{t /(1-t)}$ for $0 \leq t<1$ and $\tilde{\mathscr{F}}_{t}=\mathscr{F}_{\infty}$ for $t \geq 1$ ).

LEMMA 1.8.7. Let $M$ be a differentiable manifold and $\nabla$ a torsion-free linear connection on $M$. Every $M$-valued $\nabla$-martingale $X$ is a semimartingale up to $\infty$ on the set

$$
\left\{\lim _{t \rightarrow \infty} X_{t} \text { exists in } M\right\}
$$

Proof. Let $\left(U_{n}\right)_{n \in \mathbb{N}}$ be a covering of $M$ by open subsets $U_{n}$ with the property as in Lemma 1.8.3. As explained in the proof to Lemma 1.8.3, to each $U_{n}$ there is then a function $\varphi_{n} \in C^{\infty}\left(M ; \mathbb{R}^{n}\right)$ such that $\left(\varphi_{n} \mid U_{n}, U_{n}\right)$ defines a chart for $M$ and such that $\varphi_{n}(X)$ is a semimartingale up to $\infty$ on the set $\Omega_{n}:=\left\{X_{\infty}\right.$ exists in $\left.U_{n}\right\}$. But then $X$ itself is a semimartingale up to $\infty$ on $\Omega_{n}$, and hence also on $\bigcup_{n} \Omega_{n} \equiv\left\{X_{\infty}\right.$ exists in $\left.M\right\}$.

THEOREM 1.8.8 (Convergence Theorem of Darling-Zheng). Let $M$ be a manifold, endowed with a torsion-free linear connection $\nabla$, and $X$ be a $\nabla$-martingale on $M$. Let $g$ be an arbitrary Riemannian metric on $M$ and $[X, X]=\int g(d X, d X)$ the $g$-quadratic variation of $X$. Then (modulo sets of measure 0) the following inclusions hold true:

$$
\left\{X_{\infty} \text { exists in } M\right\} \subset\left\{[X, X]_{\infty}<\infty\right\} \subset\left\{X_{\infty} \text { exists in } \hat{M} \equiv M \cup\{\infty\}\right\}
$$

Proof. The first inclusion is a direct consequence of Lemma 1.8.7 which assures that $X$ is a semimartingale up to $\infty$ on the subset $\Omega_{0}:=\left\{X_{\infty}\right.$ exists in $\left.M\right\}$ of $\Omega$. This implies $\int_{0}^{\infty} \nabla d \varphi(d X, d X)<\infty$ almost surely on $\Omega_{0}$ for each $\varphi \in C^{\infty}(M)$, and then also $\int_{0}^{\infty} g(d X, d X)<\infty$ almost surely on $\Omega_{0}$.

For the verification of the second inclusion we note that modulo nullsets

$$
\left\{X_{\infty} \text { exists in } \hat{M}\right\}=\left\{\varphi(X) \text { converges in } \mathbb{R} \text { for each } \varphi \in C_{c}^{\infty}(M)\right\}
$$

(where $C_{c}^{\infty}(M)$ denotes again the space of test functions on $M$ ). For a fixed test function $\varphi \in C_{c}^{\infty}(M)$, by compactness reasons, there is a constant $c>0$ such that

$$
-c g \leq \nabla d \varphi \leq c g \quad \text { and } \quad d \varphi \otimes d \varphi \leq c g
$$

This allows to estimate:

$$
\begin{aligned}
& \int_{0}^{t}|\nabla d \varphi(d X, d X)| \leq c \int_{0}^{t} g(d X, d X)=c[X, X]_{t}, \quad \text { as well as } \\
& {[\varphi(X), \varphi(X)]_{t}=\int_{0}^{t}(d \varphi \otimes d \varphi)(d X, d X) \leq c[X, X]_{t} .}
\end{aligned}
$$

Let now $\varphi(X)=\varphi\left(X_{0}\right)+N+A$ denote the Doob-Meyer decomposition of $\varphi(X)$. Then both $[N, N]=[\varphi(X), \varphi(X)]$ as well as $A=\frac{1}{2} \int \nabla d \varphi(d X, d X)$ have an almost-sure limit on the set $\left\{[X, X]_{\infty}<\infty\right\}$ as $t \rightarrow \infty$, and consequently also $\varphi(X)$ itself converges on $\left\{[X, X]_{\infty}<\infty\right\}$ almost surely.

Corollary 1.8.9. Let $(M, g)$ be a Riemannian manifold and $X$ a Brownian motion $\mathrm{BM}(M, g)$ of maximal lifetime $\zeta$. For each predictable stopping time $\xi$ such that $\xi \leq \zeta$ almost surely, the following inclusions hold modulo $\mathbb{P}$-nullsets:

$$
\begin{equation*}
\left\{X_{\xi-} \text { exists in } M\right\} \subset\{\xi<\infty\} \subset\left\{X_{\xi-} \text { exists in } \hat{M}\right\} \tag{1.8.1}
\end{equation*}
$$

Proof. By means of a time change (see Remark 1.1.15) which transforms the stochastic interval $[0, \xi[$ to $[0, \infty[$, the Brownian motion $X \mid[0, \xi[$ transforms to a martingale $\hat{X}$ defined on all of $\mathbb{R}_{+}$. But then we have $[\hat{X}, \hat{X}]_{\infty}=[X, X]_{\xi}=n \xi($ where $n=\operatorname{dim} M)$, and the Convergence Theorem 1.8.8 of Darling-Zheng gives the claim.

DEFINITION 1.8.10. Let $(M, g)$ be a Riemannian manifold and $\nabla$ the Levi-Civita connection on $M$.
(i) $(M, g)$ is called stochastically complete if for each $M$-valued $\nabla$-martingale $X$,

$$
\left\{[X, X]_{\infty}<\infty\right\} \subset\left\{X_{\infty} \text { exists in } M\right\}, \quad \text { modulo } \mathbb{P} \text {-nullsets. }
$$

(ii) $(M, g)$ is called BM-complete (or complete for Brownian motions) if for each predictable stopping time $\xi>0$ and every $M$-valued Brownian motion $X$ defined on $[0, \xi[$,

$$
\{\xi<\infty\} \subset\left\{X_{\xi-} \text { exists in } M\right\}, \quad \text { modulo } \mathbb{P} \text {-nullsets. }
$$

(iii) $(M, g)$ is said to be metrically complete (or geodesically complete) if for any $x \in M$ and $v \in T_{x} M$ the unique geodesic curve $\gamma$ with $\gamma(0)=x$ and $\dot{\gamma}(0)=v$ is defined on all of $\mathbb{R}$.

REMARK 1.8.11. BM-completeness of a Riemannian manifold ( $M, g$ ) means that each Brownian motion of maximal lifetime on $(M, g)$ has actually infinite lifetime and cannot explode in finite time. Stochastic completeness of $(M, g)$ means that martingales $X$ on $(M, g)$ with finite "intrinsic time" $T_{t}=\int_{0}^{t} g(d X, d X)$ cannot explode.

Compact Riemannian manifolds $(M, g)$ are always metrically complete, and also stochastically complete by the Martingale convergence Theorem 1.8.8 of Darling-Zheng. Trivially, stochastic completeness implies BM-completeness, but not vice versa: for instance, $M=\mathbb{R}^{2} \backslash\{(1,0)\}$ is BM-complete but not stochastically complete, as can be seen from the example $X=\left(X^{1}, X^{2}\right)$ with $X^{1}$ a $\operatorname{BM}(\mathbb{R})$ and $X^{2}=0$.

REMARK 1.8.12. Stochastically complete Riemannian manifolds are metrically complete.

Proof. Assuming that $(M, g)$ is metrically incomplete, we find a geodesic $\gamma:] a, b[\rightarrow$ $M$ where $-\infty \leq a<0<b \leq \infty$ such that its domain $] a, b[$ is a proper subset of $\mathbb{R}$ which can not further be extended. Let now $Y \in \mathscr{M}$ be a convergent $] a, b[$-valued local martingale such that $Y_{0}=0$ and $Y_{\infty} \in\{a, b\} \cap \mathbb{R}$ almost surely (constructed for instance from a stopped $\mathrm{BM}(\mathbb{R})$ via time change). In particular, we have then $[Y, Y]_{\infty}<\infty$ almost surely. The composition $X:=\gamma(Y)$ is by Theorem 1.6.53 (ii) a (one-dimensional) $M$-valued martingale with the property that $\mathbb{P}\left\{X_{t}\right.$ converges for $\left.t \rightarrow \infty\right\}=0$. On the other hand, by means of pullback formula (1.3.4), we obtain

$$
[X, X]_{\infty} \equiv \int_{0}^{\infty} g(d X, d X)=\int_{0}^{\infty}\left|\dot{\gamma}\left(Y_{s}\right)\right|^{2} d[Y, Y]=|\dot{\gamma}(0)|^{2}[Y, Y]_{\infty}<\infty
$$

which shows that $(M, g)$ not stochastically complete.
The converse in Remark 1.8.12 is false in general: metrically complete Riemannian manifolds are not even BM-complete. Brownian motions on metrically complete Riemannian manifolds may explode in finite time, as will be shown in a later section. Also BM-completeness does not imply metric completeness, as can be seen from Brownian motion on $\mathbb{R}^{n} \backslash\{$ point $\}$ for $n \geq 2$.

DEFINITION 1.8.13. An exhaustion function on a differentiable manifold $M$ is a proper map $\varphi \in C^{\infty}\left(M ; \mathbb{R}_{+}\right)$. The map $\varphi$ is called proper if all sublevel sets $\{\varphi \leq c\}$ are compact, or in other words, if $\varphi(x) \rightarrow \infty$ as $x \rightarrow \infty$ in $\hat{M}$.

THEOREM 1.8.14. Let $(M, g)$ be a Riemannian manifold which carries an exhaustion function $\varphi \in C^{\infty}\left(M ; \mathbb{R}_{+}\right)$with bounded gradient (i.e. $|\operatorname{grad} \varphi| \leq$ const). Then:
(i) $(M, g)$ is metrically complete;
(ii) $(M, g) \mathrm{BM}$-complete if in addition $\Delta \varphi$ is bounded from above;
(iii) $(M, g)$ is stochastically complete if in addition $\nabla d \varphi$ is bounded from above (i.e. $\nabla d \varphi \leq c g$ for some $c>0$ ).

Proof. (i) Assume there is a geodesic curve $\gamma:[0, b[\rightarrow M$ in $M$ which cannot be extended beyond $b$. Then we have

$$
\left|(\varphi \circ \gamma)^{\prime}\right|=\left|\left\langle(\operatorname{grad} \varphi)_{\gamma(t)}, \dot{\gamma}(t)\right\rangle\right| \leq \text { const }|\dot{\gamma}(t)|=\text { const }|\dot{\gamma}(0)|<\infty
$$

so that $(\varphi \circ \gamma) \mid[0, b[$ is bounded, and since by assumption the function $\varphi$ is proper, $\gamma([0, b[)$ will be relatively compact. However this leads immediately to a contradiction: there is a sequence $\left(t_{n}\right)$ in $\left[0, b\left[\right.\right.$ such that $t_{n} \uparrow b$ with the property that $\dot{\gamma}\left(t_{n}\right)$ has a limit $v_{0} \in T M$; but then there exists a neighbourhood $V$ of $v_{0}$ in $T M$ and $\varepsilon>0$ such that each geodesic curve $\alpha$ with $\alpha(0)=\pi(v)$ and $\dot{\alpha}(0)=v \in V$, is well-defined on the interval $]-\varepsilon, \varepsilon[$; thus choosing $t_{n_{0}}>b-\varepsilon$ with $\dot{\gamma}\left(t_{n_{0}}\right) \in V$, we see that $\gamma$ can be extended beyond $b$.
(ii) Let $\Delta \varphi$ be bounded from above and let $X$ be a $\operatorname{BM}(M, g)$ of maximal lifetime $\zeta$. We want to show that $\mathbb{P}\{\zeta=\infty\}=1$. To this end denote by $\varphi(X)=\varphi\left(X_{0}\right)+N+A$ the Doob-Meyer decomposition of $\varphi(X)$. Then, in particular,

$$
[N, N]=[\varphi(X), \varphi(X)]=\int|\operatorname{grad} \varphi|^{2}(X) d t, \quad A=\frac{1}{2} \int \Delta \varphi(X) d t
$$

from where we conclude that $[N, N]_{\zeta} \leq$ const $\times \zeta$ and $\lim \sup _{t \uparrow \zeta} A_{t} \leq$ const $\times \zeta$. Hence, we have $\mathbb{P}$-a.s. the inclusion

$$
\{\zeta<\infty\} \subset\left\{\varphi(X)_{\zeta-} \text { exists in } \mathbb{R}\right\}
$$

But $\zeta$ is the maximal lifetime of $X$ and thus $X_{t} \rightarrow \infty$ in $\hat{M}$ almost surely on $\{\zeta<\infty\}$ as $t \uparrow \zeta$, and consequently $\varphi\left(X_{t}\right) \rightarrow \infty$ from where we conclude that $\mathbb{P}\{\zeta<\infty\}=0$.
(iii) Assume now $\nabla d \varphi$ to be bounded from above and let $X$ be a martingale on $(M, g)$. By the Convergence Theorem 1.8.8 of Darling-Zheng, it is sufficient to show that almost surely

$$
\left\{[X, X]_{\infty}<\infty\right\} \subset\{\varphi(X) \nrightarrow \infty\}
$$

However, by assumption, there is a constant $c>0$ such that $d \varphi \otimes d \varphi \leq c g$ and $\nabla d \varphi \leq c g$. Hence, denoting by $\varphi(X)=\varphi\left(X_{0}\right)+N+A$ the Doob-Meyer decomposition of $\varphi(X)$, we conclude

$$
\begin{aligned}
{[N, N] } & =\int(d \varphi \otimes d \varphi)(d X, d X) \leq c[X, X], \quad \text { and } \\
A & =\frac{1}{2} \int \nabla d \varphi(d X, d X) \leq \frac{1}{2} c[X, X]
\end{aligned}
$$

from where the claim follows.
We want to note already at this point how Theorem 1.8 .14 is usually applied. On a connected Riemannian manifold $(M, g)$ one constructs an exhaustion function $\varphi$ via a suitable smoothing of the distance function $\varphi_{0}=d_{M}\left(x_{0}, \cdot\right)$ to a given point $x_{0}$ in $M$.

Recall that for two points $x_{0}, x_{1}$ in $M$ the distance $d_{M}\left(x_{0}, x_{1}\right)$ is defined as the infimum length of all (piecewise) differentiable curves connecting $x_{0}$ and $x_{1}$ (cf. Definition 1.5.3). We shall see that $\left|\operatorname{grad} \varphi_{0}\right|=1$ at points where $\varphi_{0}$ is differentiable, and that there in addition $\Delta \varphi_{0}$, respectively $\nabla d \varphi_{0}$, can be controlled by curvature bounds (Hessian Comparison Theorem).

As already explained, the concept of $\nabla$-martingales covers the class of local martingales on the real line; on manifolds however a distinction of local versus true martingales is meaningless. In the scalar theory however this point is by no means only of a technical nature, for instance when it comes to questions of whether the knowledge of the state $X_{t}$ for a fixed $t>0$, together with the filtration $\left(\mathscr{F}_{s}\right)_{0 \leq s \leq t}$ up to time $t$, allows to reconstruct the whole process $X \mid[0, t]$. In scalar martingale theory, the "size" of a (local) martingale is controlled by the quadratic variational process; this aspect of the theory can be carried over to manifolds through the notion of quadratic variation of a martingale.

DEFINITION 1.8.15 ( $H^{p}$-martingale). Let $(M, g)$ be a Riemannian manifold, $\nabla$ the Levi-Civita connection on $M$ and $1 \leq p \leq \infty$. A $\nabla$-martingale $X$ on $M$ with Riemannian quadratic variation $[X, X]=\int g(d X, d X)$ is called $H^{p}$-martingale if $[X, X]_{\infty}^{1 / 2} \in L^{p}(\mathbb{P})$, i.e., $\mathbb{E}\left[[X, X]_{\infty}^{p / 2}\right]<\infty$.

Note that for stochastically complete manifolds, by the Convergence Theorem of Darling-Zheng, the condition " $[X, X]_{\infty}<\infty$ almost surely" characterizes convergent martingales $X$ on $M$. As already noted, each $\nabla$-martingale $X$ on $M$ comes by stochastic development from an $\mathbb{R}^{n}$-valued local martingale $Z$ and the Riemannian quadratic variation $[X, X]$ of $X$ coincides with the quadratic variation $[Z, Z] \equiv \sum_{i}\left[Z^{i}, Z^{i}\right]$ of $Z$, hence in particular $\mathbb{E}\left([X, X]_{\infty}^{p / 2}\right)=\mathbb{E}\left([Z, Z]_{\infty}^{p / 2}\right)$. Definition 1.8 .15 thus corresponds to the general approach to carry over $\mathbb{R}^{n}$-valued concepts to manifolds via stochastic development.

EXAMPLE 1.8.16. In the special case $M=\mathbb{R}^{n}$, a martingale $X \in \mathscr{M}_{0}\left(\mathbb{R}^{n}\right)$ is an $H^{2}$ martingale (i.e., $\left.\mathbb{E}[X, X]_{\infty}<\infty\right)$ if and only if $X_{\infty}^{*} \in L^{2}\left(\mathbb{P} ; \mathbb{R}^{n}\right)$ where $X^{*}$ denotes the (componentwise) maximal process of $X$. This condition is equivalent to $X_{t}=\mathbb{E}^{\mathscr{F}_{t}}\left[X_{\infty}\right]$ for some $X_{\infty} \in L^{2}\left(\mathbb{P} ; \mathbb{R}^{n}\right)$. Thus on $\mathbb{R}^{n}$, the $H^{2}$-martingales coincide with the class of $L^{2}$-bounded martingales. Consequently, on an $n$-dimensional Riemannian manifold, the $H^{2}$-martingales are exactly those martingales which come from an $L^{2}$-bounded $\mathbb{R}^{n}$-valued martingale via stochastic development.

REMARK 1.8.17. Since the class of $H^{p}$-martingales is invariant under time-change, Definition 1.8.15 extends in an obvious way to martingales that are only defined on a finite time interval $[0, t]$ or up to some predictable stopping time $\xi$.

THEOREM 1.8.18. Let $(M, g)$ be a Riemannian manifold and $\nabla$ the Levi-Civita connection on $M$. Let $U \subset M$ be an open subset, $\lambda: U \rightarrow \mathbb{R}_{+}$a continuous function and $\varphi \in C^{\infty}(U)$ such that $\alpha \leq \varphi \leq \beta$ (for some constants $\left.\alpha, \beta>0\right)$ such that

$$
\nabla d \varphi+2 \lambda \varphi g \leq 0
$$

Then, for each $\nabla$-martingale $X$ on $M$ taking values in $U$, it holds

$$
\mathbb{E}\left[\exp \left(\int_{0}^{\infty} \lambda(X) d[X, X]\right)\right]<\infty
$$

If in addition $\lambda \geq \varepsilon>0$ for some constant $\varepsilon$, then $X$ is an $H^{p}$-martingale for any $1 \leq p<\infty$, and in particular almost surely convergent.

Proof. Let $X$ be $\nabla$-martingale taking values in $U$. Then the real process

$$
S:=\varphi(X) \exp \left(\int \lambda(X) d[X, X]\right)
$$

is a local supermartingale, as can be seen directly by Itô's formula:

$$
\begin{aligned}
d S & =\exp (\ldots) d(\varphi(X))+\varphi(X) \exp (\ldots) \lambda(X) d[X, X] \\
& \stackrel{\mathrm{m}}{=} \exp (\ldots)\left(\frac{1}{2} \nabla d \varphi(d X, d X)+\varphi(X) \lambda(X) g(d X, d X)\right) \\
& =d(\text { decreasing process })
\end{aligned}
$$

By means of a localizing sequence of stopping times $\tau_{n} \uparrow \infty$ for $S$, we obtain for any $t \geq 0$ the estimate $\mathbb{E}\left[S_{t}^{\tau_{n}}\right] \leq \mathbb{E}\left[S_{0}^{\tau_{n}}\right]=\mathbb{E}\left[S_{0} \leq \beta\right]$, and hence by Fatou's Lemma

$$
\beta \geq \liminf _{n \rightarrow \infty} \mathbb{E}\left[S_{t}^{\tau_{n}}\right] \geq \mathbb{E}\left[\liminf _{n \rightarrow \infty} S_{t}^{\tau_{n}}\right]=\mathbb{E}\left[S_{t}\right] \geq \alpha \mathbb{E}\left[\exp \left(\int_{0}^{t} \lambda(X) d[X, X]\right)\right]
$$

This completes the proof.
THEOREM 1.8.19. Let $(M, g)$ be a Riemannian manifold and $\nabla$ be the Levi-Civita connection on $M$. Suppose that $K$ is a compact subset of $M$ such that there is a strictly convex $C^{\infty}$-function defined on an open neighbourhood of $K$. Then each $\nabla$-martingale on $M$ taking its values in $K$ is an $H^{p}$-martingale for $1 \leq p<\infty$, and hence almost surely convergent.

Proof. By assumption there is an open set $U$ containing $K$ and carrying a strictly convex function $\varphi \in C^{\infty}(U)$. Multiplying $\varphi$ by -1 we have $\nabla d \varphi<0$ on $U$. Without restrictions we may assume that $\varphi$ is bounded and, if necessary by adding a positive constant, that $\alpha \leq \varphi \leq \beta$ with $\alpha, \beta>0$. By compactness reasons, we may assume, possibly after reducing the size of $U$, that even $\nabla d \varphi+2 \varepsilon \varphi g \leq 0$ holds on $U$ for some sufficiently small $\varepsilon>0$. The claim then follows from Theorem 1.8.18 with $\lambda \equiv \varepsilon$.

Note that Theorem 1.8 .19 covers the well-known fact that bounded $\mathbb{R}^{n}$-valued local martingales converge almost surely. However, as already mentioned, manifold-valued martingales taking values in a compact set are not at all convergent in general.

We want to discuss another well-known property of continuous real martingales in the case of $M$-valued martingales. For a real martingale $X$, the knowledge of $X_{t}$ at a fixed time $t>0$, together with the filtration $\left(\mathscr{F}_{s}\right)_{0 \leq s \leq t}$, already determines the martingale $\left(X_{s}\right)_{0 \leq s \leq t}$ up to time $t$, namely as $X_{s}=\mathbb{E}^{\mathscr{F}_{s}}\left[X_{t}\right]$ almost surely. An equivalent formulation of this property is that if $X$ and $Y$ are continuous real martingales adapted to the same filtration and if $X_{t}=Y_{t}$ for some $t>0$, then already $X|[0, t]=Y|[0, t]$ modulo indistinguishability. We call this property non-confluence of real martingales. This leads to the question to what extent it is possible to have confluence of non-identical martingales on manifolds at a certain time.

THEOREM 1.8.20 (Minimum principle). Let $(M, g)$ be a Riemannian manifold and $\nabla$ the Levi-Civita connection on $M$. Let $\varphi \in C(M ; \mathbb{R})$ and $U:=\{\varphi>0\}, D:=\{\varphi \leq 0\}$. Furthermore, let $\lambda: U \rightarrow \mathbb{R}_{+}$be a continuous function. Suppose that $\varphi$ is bounded on $U$ and $\varphi \mid U \in C^{\infty}(U)$, and in addition $\nabla d \varphi+2 \lambda \varphi g \geq 0$ on $U$. Extending $\lambda$ to all of $M$ by $d \lambda \mid D:=0$, the following statement holds: If $a \nabla$-martingale $X$ with the property

$$
\begin{equation*}
\mathbb{E}\left[\exp \left(\int_{0}^{\infty} \lambda(X) d[X, X]\right)\right]<\infty \tag{1.8.2}
\end{equation*}
$$

converges to a $D$-valued random variable, then $X$ lives completely in $D$.

Proof. Let $X$ be a $\nabla$-martingale with the property (1.8.2) such that $X_{t} \rightarrow X_{\infty}$ almost surely for some $D$-valued random variable $X_{\infty}$. As $X$ has continuous paths and $U$ is open, it is sufficient to show that $\mathbb{P}\left\{X_{t_{0}} \in U\right\}=0$ for each fixed $t_{0} \geq 0$.

To this end let $S:=\varphi(X) \exp \left(\int \lambda(X) d[X, X]\right)$ and $\tau=\inf \left\{t \geq t_{0}: X_{t} \notin U\right\}$. As in the proof of Theorem 1.8.18 one verifies that

$$
Y_{t}:=1_{\left\{X_{t_{0}} \in U\right\}} S_{t_{0}+(t \wedge \tau)}, \quad t \geq 0
$$

defines a non-negative local submartingale (with respect to the filtration $\left(\hat{\mathscr{F}}_{t}\right)_{t \geq 0}$ where $\left.\hat{\mathscr{F}}_{t}:=\mathscr{F}_{t_{0}+t}\right)$. By assumption, setting $\alpha:=\sup (\varphi \mid U)$, we have

$$
Y_{t} \leq \alpha \exp \left(\int_{0}^{\infty} \lambda(X) d[X, X]\right) \in L^{1}(\mathbb{P})
$$

so that the process $Y$ is uniformly integrable. From the fact that $Y_{\infty}=0$ almost surely, it follows that $0 \leq \mathbb{E}\left[Y_{0}\right] \leq \mathbb{E}\left[Y_{\infty}\right]=0$. Since $Y_{0} \mid\left\{X_{t_{0}} \in U\right\}>0$ we conclude that $\mathbb{P}\left\{X_{t_{0}} \in U\right\}=0$.

DEFINITION 1.8.21 (Convex geometry). Let $M$ be a manifold equipped with a torsionfree linear connection. An open subset $V$ of $M$ is said to have convex geometry if there exists a non-negative convex smooth function

$$
\phi: V \times V \rightarrow \mathbb{R}_{+}
$$

which vanishes exactly on the diagonal $\Delta=\{(x, x): x \in V\}$.
Convexity of $\phi$ in Definition 1.8.21 is understood with respect to the direct sum connection on $M \times M$, see Remark 1.7.1.

REMARK 1.8 .22 . Let $\nabla, \nabla^{\prime}$ be linear connections on differentiable manifolds $M$ resp. $M^{\prime}$. The direct sum connection on the product $\bar{M}=M \times M^{\prime}$ is given as follows: Each vector field $\bar{A} \in \Gamma\left(\bar{\alpha}^{*} T \bar{M}\right)$ along a curve $\bar{\alpha}=\left(\alpha, \alpha^{\prime}\right)$ on $\bar{M}$ decomposes as $\bar{A}=\left(A, A^{\prime}\right)$ with a vector field $A$ along $\alpha$ on $M$ and a vector field $A^{\prime}$ along $\alpha^{\prime}$ on $M^{\prime}$. The covariant derivative of $\bar{A}$ along $\bar{\alpha}$ is given by

$$
\begin{equation*}
\bar{\nabla}_{D} \bar{A}:=\left(\nabla_{D} A, \nabla_{D} A^{\prime}\right) \tag{1.8.3}
\end{equation*}
$$

An immediate consequence is that a curve $\bar{\gamma}=\left(\gamma, \gamma^{\prime}\right): I \rightarrow \bar{M}$ is a geodesic if and only if $\gamma: I \rightarrow M$ and $\gamma^{\prime}: I \rightarrow M^{\prime}$ are both geodesic curves. By Corollary 1.7.10, a differentiable function $\phi: M \times M^{\prime} \rightarrow \mathbb{R}$ is hence convex if for all geodesics $\gamma$ on $M$ and $\gamma^{\prime}$ on $M^{\prime}$ the curve $t \mapsto \phi\left(\gamma(t), \gamma^{\prime}(t)\right)$ is convex, i.e., a curve with non-negative second derivative.

Combined with Corollary 1.7 .11 this shows that both the projections pr: $\bar{M} \rightarrow M$ and $\mathrm{pr}^{\prime}: \bar{M} \rightarrow M^{\prime}$, as well as the canonical embeddings $\iota_{x}: M^{\prime} \rightarrow \bar{M}$ and $\iota_{x^{\prime}}: M \rightarrow \bar{M}$ for $x \in M$ resp. $x^{\prime} \in M^{\prime}$ are affine maps. From the probabilistic perspective, according to Theorem 1.7.15 (i), a process $\bar{X}=\left(X, X^{\prime}\right)$ is hence a $\bar{\nabla}$-martingale on $\bar{M}$ if $X$ and $X^{\prime}$ are martingales on $M$, resp. $M^{\prime}$ (adapted to the same filtration).

If $M$ and $M^{\prime}$ are in addition Riemannian manifolds with Riemannian metrics $g$, resp. $g^{\prime}$, then canonically also $\bar{M}=M \times M^{\prime}$ is a Riemannian manifold where the metric $\bar{g}$ is given by

$$
\bar{g}_{\bar{x}}(\bar{v}, \bar{w}):=g_{x}(v, w)+g_{x^{\prime}}^{\prime}\left(v^{\prime}, w^{\prime}\right), \quad \bar{x}=\left(x, x^{\prime}\right), \bar{v}=\left(v, v^{\prime}\right), \bar{w}=\left(w, w^{\prime}\right)
$$

and (1.8.3) defines the Levi-Civita connection on $\bar{M}$.

REMARK 1.8.23. We want to give a description of the direct sum connection $\bar{\nabla}$ via the corresponding frame bundles. Let $P:=\operatorname{Iso}_{G}\left(\mathbb{R}^{n} ; E\right) \rightarrow M$ be the principal $G$-bundle associated to the tangent bundle $E:=T M \rightarrow M$, i.e., $P=\mathrm{L}(T M)$ with $G=\mathrm{GL}(n ; \mathbb{R})$, resp. $P=\mathrm{O}(T M)$ with $G=\mathrm{O}(n)$ in the case of a Riemannian manifold $(M, g)$. With the analogous notations we have the principal $G^{\prime}$-bundle $P^{\prime}:=\operatorname{Iso}_{G^{\prime}}\left(\mathbb{R}^{n^{\prime}} ; E^{\prime}\right) \rightarrow M^{\prime}$ over $M^{\prime}$. Induced by the linear connections in $E$ and $E^{\prime}$ we have a $G$-connection $H$ in $P$ and a $G^{\prime}$-connection $H^{\prime}$ in $P^{\prime}$. The product $\tilde{P}:=P \times P^{\prime} \rightarrow \bar{M}$ is then a principal $\tilde{G}$-bundle with $\tilde{G}:=G \times G^{\prime}$ and

$$
\tilde{H}:=H \times H^{\prime} \subset T P \times T P^{\prime}=T \tilde{P}
$$

gives a $\tilde{G}$-connection in $\tilde{P}$. The corresponding connection form and canonical one-form on $\tilde{P}$ are given by $\tilde{\omega}=\left(\omega, \omega^{\prime}\right) \in \Gamma\left(T^{*} \tilde{P} \otimes \tilde{\mathfrak{g}}\right)$, respectively $\left(\vartheta, \vartheta^{\prime}\right) \in \Gamma\left(T^{*} \tilde{P} \otimes \mathbb{R}^{n+n^{\prime}}\right)$. Letting now $\bar{E}:=E \times E^{\prime}$ and $\bar{n}:=n+n^{\prime}$, and $\bar{G}=\mathrm{GL}(\bar{n} ; \mathbb{R})$ resp. $\bar{G}=\mathrm{O}(\bar{n})$, there is a canonical homomorphism of Lie groups

$$
\alpha: \tilde{G} \rightarrow \bar{G}, \quad\left(A_{1}, A_{2}\right) \mapsto\left(\begin{array}{cc}
A_{1} & 0  \tag{1.8.4}\\
0 & A_{2}
\end{array}\right) .
$$

On the other hand, to $\left(u, u^{\prime}\right) \in \tilde{P}$ we get a linear isomorphism $\bar{u}: \mathbb{R}^{\bar{n}} \rightarrow \bar{E}_{\bar{\pi}\left(u, u^{\prime}\right)}$ in an obvious way, and the induced map over $\bar{M}$,

$$
\sigma: \tilde{P} \rightarrow \bar{P}, \quad \tilde{u}=\left(u, u^{\prime}\right) \mapsto \bar{u}
$$

is given in bundle charts by $\alpha$. Obviously $\sigma$ is equivariant: $\sigma(\tilde{u} \tilde{g})=\sigma(\tilde{u}) \alpha(\tilde{g})$ for $\tilde{u} \in \tilde{P}$ and $\tilde{g} \in \tilde{G}$, more precisely,

$$
\sigma\left(u, u^{\prime}\right)=\left(d \iota_{x^{\prime}}\right)_{x} \circ u \circ \operatorname{pr}+\left(d \iota_{x}\right)_{x^{\prime}} \circ u^{\prime} \circ \operatorname{pr}^{\prime} \quad \text { for }\left(u, u^{\prime}\right) \in \tilde{P}_{\left(x, x^{\prime}\right)}
$$

The $\bar{G}$-connection $\bar{H}$ in $\bar{P}$ induced by $\tilde{H}$ is finally given by

$$
\bar{H}_{\sigma(\tilde{u}) g}=\left(R_{g} \circ \sigma\right)_{*} \tilde{H}_{\tilde{u}}, \quad g \in \bar{G} .
$$

For the connection form $\bar{\omega} \in \Gamma\left(T^{*} \bar{P} \otimes \overline{\mathfrak{g}}\right)$ induced by $\bar{H}$ and the canonical one-form $\bar{\vartheta} \in \Gamma\left(T^{*} \bar{P} \otimes \mathbb{R}^{\bar{n}}\right)$ on $\bar{P}$ then obviously

$$
\begin{equation*}
\sigma^{*} \bar{\omega}=(d \alpha)_{1} \circ\left(\omega, \omega^{\prime}\right) \quad \text { and } \quad \sigma^{*} \bar{\vartheta}=\left(\vartheta, \vartheta^{\prime}\right) \tag{1.8.5}
\end{equation*}
$$

LEMMA 1.8.24. Let $\nabla$ and $\nabla^{\prime}$ be linear connections on differentiable manifolds $M$, resp. $M^{\prime}$. Suppose that $X$ and $X^{\prime}$ are semimartingales taking values in $M$, resp. $M^{\prime}$ (adapted to the same filtration); let $U$ and $U^{\prime}$ be the corresponding horizontal lifts in $P$, resp. $P^{\prime}$, as well as $Z=\int_{U} \vartheta$ and $Z^{\prime}=\int_{U^{\prime}} \vartheta^{\prime}$ the anti-developments taking values in $\mathbb{R}^{n}$, resp. $\mathbb{R}^{n^{\prime}}$. Suppose that $\bar{M}=M \times M^{\prime}$ is equipped with the canonical direct sum connection. Then $\bar{U}:=\sigma\left(U, U^{\prime}\right)$ taking values in $\mathrm{L}(\bar{M})$, resp. $\mathrm{O}(\bar{M})$, is a horizontal lift of the semimartingale $\bar{X}:=\left(X, X^{\prime}\right)$, and the $\mathbb{R}^{\bar{n}}$-valued anti-development of $\bar{X}$ is given by $\left(Z, Z^{\prime}\right)$.

Proof. Since $\bar{\pi} \circ \bar{U}=\left(\pi, \pi^{\prime}\right) \circ\left(U, U^{\prime}\right)=\left(X, X^{\prime}\right)$, we calculate by means of (1.8.5)

$$
\begin{aligned}
& \int_{\bar{U}} \bar{\omega}=\int_{\left(U, U^{\prime}\right)} \sigma^{*} \bar{\omega}=(d \alpha)_{1}\left(\int_{\left(U, U^{\prime}\right)}\left(\omega, \omega^{\prime}\right)\right)=(d \alpha)_{1}\left(\int_{U}(\omega, 0)+\int_{U^{\prime}}\left(0, \omega^{\prime}\right)\right)=0 \quad \text { and } \\
& \int_{\bar{U}} \bar{\vartheta}=\int_{\left(U, U^{\prime}\right)} \sigma^{*} \bar{\vartheta}=\int_{\left(U, U^{\prime}\right)}\left(\vartheta, \vartheta^{\prime}\right)=\int_{U}(\vartheta, 0)+\int_{U^{\prime}}\left(0, \vartheta^{\prime}\right)=(Z, 0)+\left(0, Z^{\prime}\right)=\left(Z, Z^{\prime}\right)
\end{aligned}
$$

which verifies the claim.

After these technical remarks we now return to the notion of convex geometry introduced in Definition 1.8.21.

THEOREM 1.8.25 (Non-confluence of $\nabla$-martingales). Let $M$ be a differentiable manifold endowed with a torsion-free linear connection and $K \subset M$ be a compact subset such that an open neighbourhood $V$ of $K$ has convex geometry. If then, for a given filtration, $X$ and $X^{\prime}$ are $\nabla$-martingales on $M$ taking values in $K$, which both converge and such that $X_{\infty}=X_{\infty}^{\prime}$ almost surely, then already $X=X^{\prime}$ modulo indistinguishability.

Proof. Suppose that $\phi: V \times V \rightarrow \mathbb{R}_{+}$describes the convex geometry of $V$, where without loss of generality we may assume $\phi$ to be bounded. For Riemannian manifolds the claim follows from the minimum principle (Theorem 1.8.20) with $\lambda=0$ (applied on the product manifold $V \times V)$. In the general case the proof of Theorem 1.8.20 carries over verbatim with $\lambda=0$. Indeed then $\left(X, X^{\prime}\right)$ is a martingale on $M \times M$ taking values in $V \times V$ and $Y:=\phi\left(X, X^{\prime}\right)$ a bounded non-negative submartingale. As by assumption $Y_{\infty}=0$, almost surely, we conclude that for each $t \geq 0$,

$$
0 \leq Y_{t} \leq \mathbb{E}^{\mathscr{F}_{t}}\left[Y_{\infty}\right]=0 \quad \text { almost surely }
$$

and hence $X_{t}=X_{t}^{\prime}$ almost surely, since $\phi$ vanishes only on the diagonal.
DEFINITION 1.8.26 (Totally geodesic submanifold). Let $M_{0}$ and $M$ be differentiable manifolds, equipped with torsion-free linear connections, and let $M_{0} \stackrel{\iota}{\longleftrightarrow} M$ be an embedding. The manifold $M_{0}$ is called a totally geodesic submanifold of $M$ if the inclusion map $\iota$ is affine.

Without loss of generality we may consider $M_{0}$ as subspace of $M$.
REMARK 1.8 .27 . An embedding $M_{0} \stackrel{\iota}{\longleftrightarrow} M$ is obviously a totally geodesic submanifold $M_{0}$ of $M$ if and only if $\iota$ transfers $M_{0}$-geodesics into $M$-geodesics, or equivalently: if $x_{0} \in M_{0}$ and $v_{0} \in T_{x_{0}} M_{0}$, as well as $\gamma$ a geodesic curve in $M$ such that $\gamma(0)=\iota\left(x_{0}\right)$ and $\dot{\gamma}(0)=\iota_{*} v_{0}$, then $\gamma(]-\varepsilon, \varepsilon[) \subset \iota\left(M_{0}\right)$ for some $\varepsilon>0$.

EXAMPLE 1.8.28. Let $M$ be a connected differentiable manifold and $\varphi: M \rightarrow \mathbb{R}$ an affine function. Then $M_{0}=\{\varphi=c\}$ defines a totally geodesic submanifold of $M$ for each $c \in \mathbb{R}$.

Proof. At first, for affine functions $\varphi$, we remark that $(d \varphi)_{x}=0$ for some $x \in M$ already implies $d \varphi=0$ locally about $x$. Indeed, for a geodesic curve $\gamma$ in $M$ with $\gamma(0)=x$ the composition $\varphi \circ \gamma$ defines a straight line with slope $(\varphi \circ \gamma)^{\prime}(0)=0$; hence $\varphi$ is constant locally about $x$. Thus, since by assumption $M$ is connected, the existence of a critical point of $\varphi$ means that $\varphi \equiv$ const. Hence assume now $(d \varphi)_{x} \neq 0$ for each $x \in M$ and consider $M_{0}=\varphi^{-1}\{c\}$. In particular, $c$ is then a regular value for $\varphi$. For $x_{0} \in M_{0}$ and $v \in T_{x_{0}} M_{0}$ let $\gamma$ be the geodesic in $M$ determined by $\gamma(0)=x_{0}$ and $\dot{\gamma}(0)=v$. We conclude that then $(\varphi \circ \gamma)^{\prime}(0)=0$ and consequently $\varphi \circ \gamma \equiv \varphi \circ \gamma(0)=c$ which shows that $\gamma$ lies in $M_{0}$.

Lemma 1.8.29. Let $M$ be a differentiable manifold equipped with a torsion-free linear connection. Let $M_{0}$ be a totally geodesic submanifold of $M$ and $x_{0} \in M_{0}$. There exists an open neighbourhood $V$ of $x_{0}$ in $M$ and a convex function $\varphi \in C^{\infty}(V)$ such that

$$
\{\varphi=0\}=V \cap M_{0} \quad \text { and } \quad\{\varphi>0\}=V \backslash M_{0} .
$$

Proof. Denoting by $\operatorname{codim}\left(M_{0}\right)=\operatorname{dim} M-\operatorname{dim} M_{0}=n-n_{0}$ the codimension of $M_{0}$, we may assume without loss of generality that $0<\operatorname{codim}\left(M_{0}\right)<n$. As in the proof of Lemma 1.7.22 we introduce normal coordinates $(h, V)$ for $M$ about $x_{0}$ via the
exponential map. As $M_{0}$ is totally geodesic, $h$ can be chosen such that $h=(\phi, \psi)$ and $V \cap M_{0}=\{\psi=0\}$. Then all Christoffel symbols $\Gamma_{i j}^{k}$ vanish at the point $x_{0}$. On the other hand, as $M_{0}$ is totally geodesic, all geodesic curves $t \mapsto \gamma(t)$ on $M$ with initial condition $\gamma(0)=x \in M_{0}$ and $\dot{\gamma}(0) \in T_{x} M_{0}$ stay in $M_{0}$ for small values of $t$, which in addition implies for $x \in V \cap M_{0}$,

$$
\begin{equation*}
\Gamma_{i j}^{k}(x)=0, \quad 1 \leq i, j \leq n_{0}, \quad n_{0}+1 \leq k \leq n \tag{1.8.6}
\end{equation*}
$$

as can be seen from the description of geodesic curves in coordinates,

$$
\ddot{\gamma}^{k}(t)+\sum_{i, j} \Gamma_{i j}^{k}(\gamma(t)) \dot{\gamma}^{i}(t) \dot{\gamma}^{j}(t)=0
$$

We are going to show that, with an appropriate choice of the constant $c>0$, the function

$$
\begin{equation*}
\varphi:=\frac{1}{2}\left(c+|\phi|^{2}\right)|\psi|^{2} \tag{1.8.7}
\end{equation*}
$$

satisfies the claim (after possibly shrinking of $V$ ). To this end, we have to verify that $\varphi$ is convex on a neighbourhood of $x_{0}$.

We start by calculating $(\nabla d \varphi)_{i j}=\partial_{i} \partial_{j} \varphi-\sum_{k} \Gamma_{i j}^{k} \partial_{k} \varphi$ in the chart $(h, V)$; see (1.5.8). Denoting the components of $h$ by $\left(\phi^{1}, \ldots, \phi^{n_{0}}, \psi^{n_{0}+1}, \ldots, \psi^{n}\right)$, it holds that

$$
\partial_{i} \partial_{j} \varphi= \begin{cases}\delta_{i j}|\psi|^{2} & \text { for } 1 \leq i, j \leq n_{0} \\ 2 \phi^{i} \psi^{j} & \text { for } 1 \leq i \leq n_{0}, \quad n_{0}+1 \leq j \leq n \\ \delta_{i j}\left(c+|\phi|^{2}\right) & \text { for } n_{0}+1 \leq i, j \leq n\end{cases}
$$

as well as

$$
\sum_{k=1}^{n} \Gamma_{i j}^{k} \partial_{k} \varphi=|\psi|^{2} \sum_{k=1}^{n_{0}} \Gamma_{i j}^{k} \phi^{k}+\left(c+|\phi|^{2}\right) \sum_{k=n_{0}+1}^{n} \Gamma_{i j}^{k} \psi^{k}
$$

Using the abbreviation $H_{i j}:=(\nabla d \varphi)_{i j}$, we have to show that $H$ on $V \backslash M_{0}=V \cap\{\psi \neq 0\}$ is positive definite for some sufficiently small open neighbourhood $V$ of the point $x_{0}$. In terms of the decomposition of $\{1, \ldots, n\}$ into $I=\left\{1, \ldots, n_{0}\right\}$ and $J=\left\{n_{0}+1, \ldots, n\right\}$, we see that $H$ is positive definite on $V \cap\{\psi \neq 0\}$ if and only if

$$
H^{*}:=\left(\begin{array}{cr}
\frac{1}{| |^{2}}\left(H_{i j}\right)_{(i, j) \in I \times I} & \frac{1}{|\psi|}\left(H_{i j}\right)_{(i, j) \in I \times J} \\
\frac{1}{|\psi|}\left(H_{i j}\right)_{(i, j) \in J \times I} & \left(H_{i j}\right)_{(i, j) \in J \times J}
\end{array}\right)
$$

is positive definite on $V \cap\{\psi \neq 0\}$. However, as easily seen, it can be achieved that $H^{*}$ on $V \cap\{\psi \neq 0\}$ is arbitrarily close to the matrix

$$
\left(\begin{array}{rr}
1 & 0 \\
0 & c
\end{array}\right)
$$

by sufficiently shrinking the neighbourhood $V$ of $x_{0}$ and choosing the constant $c>0$ in (1.8.7) small enough. To see this for the coefficients of the first quadrant of $H^{*}$, we use that on a sufficiently small neighbourhood of $x_{0}$, as a consequence of (1.8.6), one can estimate

$$
\left|\Gamma_{i j}^{k}\right| \leq C\left|\psi^{k}\right|, \quad 1 \leq i, j \leq n_{0}, \quad n_{0}+1 \leq k \leq n
$$

with a constant $C>0$.
COROLLARY 1.8.30. Let the manifold $M$ be equipped with a torsion-free linear connection $\nabla$ and suppose that $M_{0}$ is a totally geodesic submanifold of $M$. Then every point in $M_{0}$ possesses an open neighbourhood $V$ in $M$ with the property: if $X$ is a $\nabla$-martingale on $M$ taking values in $V$ such that for some $t>0$, the variable $X_{t}$ takes almost surely its values in $M_{0}$, then $X \mid[0, t]$ lives entirely in $M_{0}$.

Proof. We choose $\varphi$ as in Lemma 1.8.29 where without loss of generality we may assume that $\varphi$ is bounded. Then the composition $\varphi(X)$ defines a bounded non-negative submartingale with the property that $\varphi\left(X_{t}\right)=0$ almost surely, and thus necessarily already $\varphi(X) \equiv 0$ on $[0, t]$ modulo indistinguishability.

THEOREM 1.8.31. Let $M$ be a differentiable manifold and $\nabla$ be a torsion-free linear connection on $M$. Then each point has a neighbourhood $V$ of convex geometry. In particular, each point has a neighbourhood $V$ with the following property: if $X$ and $X^{\prime}$ are two $V$-valued $\nabla$-martingales on $M$ (adapted to the same filtration) such that $X_{t}=X_{t}^{\prime}$ almost surely, for some $t>0$, then already $X=X^{\prime}$ on $[0, t]$ modulo indistinguishability.

Proof. The first part of the claim follows from Lemma 1.8.29, applied to the diagonal manifold $M \stackrel{\iota}{\longleftrightarrow} M \times M, x \mapsto(x, x)$, considered as totally geodesic submanifold of the $M \times M$ where the product $M \times M$ is equipped with the direct sum connection. The second part follows from Corollary 1.8.30 or also directly from Theorem 1.8.25.

The results of this section show that the behavior of martingales on manifolds can be controlled in domains which are "small" in the sense that they support convex functions with specific properties. Aspects of global martingale theory, however, such as interaction with global geometry, are unaffected by this.

The local existence of suitable convex functions is ensured by Theorem 1.8.25 and Theorem 1.8.31, where the question naturally arises as to how large the domains of such "convexity areas" can be chosen in concrete cases. In the next Chapter, among other things, we will connect such questions with the concept of the curvature of a Riemannian manifold.

### 1.9. Stochastic Differentials and Second Order Tangent Spaces

This Section is not intended for the introduction of new concepts; the aim is rather to regard the methods of Stochastic Analysis as developed so far from a different perspective.

The symbol $d X$ for an $M$-valued semimartingale $X$ has so far not been considered as a mathematical object by its own; it has been used as a formal notation which received a precise meaning only by composition with scalar-valued functions. The object $d X$ is however interesting as it does not behave as a tangent vector, for instance in the intuitive sense of " $d X_{t} \in T_{X_{t}} M$ ", as one might think, rather it shares the formal properties of a second order tangent vector. In this Section we want to investigate how to interpret $d X$ as a "section of the second order tangent bundle $T^{2} M \rightarrow M$ along $X$ ", giving a meaning to $d X_{t} \in T_{X_{t}}^{2} M$; see [12], [36], [35], as well as [39], [40] for more detailed expositions in this direction.

As differential geometry of second order is not commonly widespread, the use of concepts of second order concepts in Stochastic Analysis gives sometimes the impression that Stochastic Differential Geometry requires a revision of standard differential geometry. We tried so far to counteract this impression by a consequent use of standard geometric notations. Nevertheless it is interesting to note that the simple concept of a second order tangent space permits a geometric description of the transformation behaviour given by Itô's formula.

In this Section, $M$ will always be a differentiable manifold, for $x \in M, \mathscr{E}_{x}(M)$ will denote the ring of germs of $C^{\infty}$ functions on $M$ at the point $x$, and $\mathfrak{m}_{x}(M) \subset \mathscr{E}_{x}(M)$ the maximal ideal of germs $[\varphi]$ with the property that $\varphi(x)=0$.

DEFINITION 1.9.1 (Tangent space of order $k$ ). Let $M$ be a manifold and $k \in \mathbb{N}$. The finite dimensional real vector space

$$
\begin{aligned}
T_{x}^{k} M: & =\left(\mathfrak{m}_{x}(M) / \mathfrak{m}_{x}(M)^{k+1}\right)^{*} \\
& \equiv\left\{v \in \mathscr{E}_{x}(M)^{*}: v(1)=0 \text { and } v\left(\varphi^{k+1}\right)=0, \text { if } \varphi(x)=0\right\}
\end{aligned}
$$

is called tangent space of order $k$ to $M$ at the point $x$.
Tangent spaces of first order in the sense of Definition 1.9.1 coincide with the usual tangent spaces.

REMARK 1.9.2. For any manifold $M$ it holds $T_{x}^{1} M=T_{x} M$ for all $x \in M$.
Proof. The inclusion $T_{x} M \subset T_{x}^{1} M$ is obvious, since by definition $T_{x}^{1} M$ is the real vector space of derivations at $x$, i.e.

$$
T_{x} M=\left\{v \in \mathscr{E}_{x}(M)^{*}: v(\varphi \psi)=\varphi(x) v(\psi)+\psi(x) v(\varphi)\right\}
$$

Conversely, assume that $v \in T_{x}^{1} M$; we want to show that $v(\varphi \psi)=\varphi(x) v(\psi)+\psi(x) v(\varphi)$ for all $\varphi, \psi \in \mathscr{E}_{x}(M)$. By the equation $2 \varphi \psi=(\varphi+\psi)^{2}-\varphi^{2}-\psi^{2}$ it is sufficient to verify $v\left(\varphi^{2}\right)=2 \varphi(x) v(\varphi)$ for $\varphi \in \mathscr{E}_{x}(M)$. By assumption we have $v(\varphi-\varphi(x))^{2}=0$, from where
$0=v\left(\varphi^{2}-2 \varphi(x) \varphi+\varphi(x)^{2}\right)=v\left(\varphi^{2}\right)-2 \varphi(x) v(\varphi)+v\left(\varphi(x)^{2}\right)=v\left(\varphi^{2}\right)-2 \varphi(x) v(\varphi)$
follows. This shows the inclusion $T_{x}^{1} M \subset T_{x} M$.
DEFINITION 1.9.3 (Cotangent space of order $k$, differential of order $k$ ). Let $M$ be a manifold and $k \in \mathbb{N}$. The finite dimensional real vector space

$$
T_{x}^{* k} M:=\left(T_{x}^{k} M\right)^{*}=\mathfrak{m}_{x}(M) / \mathfrak{m}_{x}(M)^{k+1}
$$

is called cotangent space of order $k$ to $M$ at the point $x$. For $\varphi \in C^{\infty}(M)$ resp., $\varphi \in$ $\mathscr{E}_{x}(M)$, we denote $\left(d^{k} \varphi\right)_{x}:=[\varphi-\varphi(x)] \in \mathfrak{m}_{x}(M) / \mathfrak{m}_{x}(M)^{k+1}$ the differential of order $k$ of $\varphi$ at the point $x$.

Obviously $\left(d^{1} \varphi\right)_{x} \equiv d \varphi_{x} \in T_{x}^{*} M$. In addition, it is easy to see that $T^{k} M \rightarrow M$ and $T^{* k} M \rightarrow M$ constitute vector bundle over $M$. Differentiable sections of these bundles, i.e. elements of $\Gamma\left(T^{k} M\right)$ resp. $\Gamma\left(T^{* k} M\right)$, are called vector fields of order $k$, resp. differential forms of order $k$.

REMARK 1.9.4 (Push-forward and pull-back). Every differentiable map $f: M \rightarrow N$ between manifolds induces canonically vector bundle homomorphisms

$$
f_{*}: T^{k} M \rightarrow f^{*} T^{k} N \quad \text { resp., } \quad f^{*}: f^{*}\left(T^{* k} N\right) \rightarrow T^{* k} M
$$

namely $\left(f_{*}\right)_{x}: T_{x}^{k} M \rightarrow T_{f(x)}^{k} N, v \longmapsto\left(f_{*}\right)_{x} v$ where $\left(f_{*}\right)_{x} v(\varphi):=v(\varphi \circ f)$ for $\varphi \in$ $\mathscr{E}_{f(x)}(N)$, and $\left(f^{*}\right)_{x}: T_{f(x)}^{* k} N \rightarrow T_{x}^{* k} M, \vartheta \longmapsto \vartheta \circ\left(f_{*}\right)_{x}$.

In the sequel we focus on the case $k=2$ and we want first to check that vector fields of second order correspond to differential operators without constant term of order at most two. For $L \in \Gamma\left(T^{2} M\right)$ and $\varphi \in \mathscr{E}_{x}(M)$ let $(L \varphi)(x):=L_{x} \varphi$ where $L_{x} \in T_{x}^{2} M$, i.e. $L_{x} \in \mathscr{E}_{x}(M)^{*}$ with $L_{x} 1=0$ and $L_{x} \varphi^{3}=0$ if $\varphi(x)=0$. Writing $\varphi=\bar{\varphi} \circ h$ in a chart $h$ at $x$ with $h(x)=0$ and $\bar{\varphi} \in \mathscr{E}_{0}\left(\mathbb{R}^{n}\right)$, then $\bar{\varphi}$ can be represented by Taylor's formula as

$$
\bar{\varphi}(y)=\bar{\varphi}(0)+\sum_{i}\left(D_{i} \bar{\varphi}\right)(0) y^{i}+\sum_{i, j} \gamma_{i j}(y) y^{i} y^{j}
$$

where $\gamma_{i j} \in \mathscr{E}_{0}\left(\mathbb{R}^{n}\right)$ is such that $\gamma_{i j}(0)=\frac{1}{2}\left(D_{i} D_{j} \bar{\varphi}\right)(0)$. Hence it holds

$$
\begin{aligned}
&(L \varphi)(x)= L_{x}\left(\sum_{i}\left(D_{i} \bar{\varphi}\right)(0) h^{i}\right. \\
&\left.=\sum_{i, j}\left(\gamma_{i j} \circ h\right) h^{i} h^{j}\right) \\
&=\sum_{i}\left(D_{i} \bar{\varphi}\right)(0) L_{x} h^{i}+\frac{1}{2} \sum_{i, j}\left(D_{i} D_{j} \bar{\varphi}\right)(0) L_{x}\left(h^{i} h^{j}\right) \\
&+\sum_{i, j} L_{x}\left(\left[\gamma_{i j} \circ h-\frac{1}{2}\left(D_{i} D_{j} \bar{\varphi}\right)(0)\right] h^{i} h^{j}\right)
\end{aligned}
$$

But we have $L_{x}(f g h)=0$ for functions $f, g, h$ defined locally about $x$ with the property that $f(x)=g(x)=h(x)=0$ which implies $L_{x}\left(\left[\gamma_{i j} \circ h-\frac{1}{2}\left(D_{i} D_{j} \bar{\varphi}\right)(0)\right] h^{i} h^{j}\right)=0$. Hence, letting $b^{i}(x)=L_{x} h^{i}$ and $a^{i j}(x)=\frac{1}{2} L_{x}\left(h^{i} h^{j}\right)$, we obtain, as wanted, a representation of the form

$$
\begin{equation*}
L_{x}=\sum_{i} b^{i}(x)\left(\frac{\partial}{\partial h^{i}}\right)_{x}+\sum_{i, j} a^{i j}(x)\left(\frac{\partial}{\partial h^{i}}\right)_{x}\left(\frac{\partial}{\partial h^{j}}\right)_{x} . \tag{1.9.1}
\end{equation*}
$$

This also shows that $T M \hookrightarrow T^{2} M$ is canonically a subbundle of $T^{2} M$.
REMARK 1.9.5. Functions $\varphi, \psi \in C^{\infty}(M)$ on a manifold $M$ induce canonically sections of $T^{* 2} M \rightarrow M$ through

$$
\begin{array}{rlrl}
d^{2} \varphi \in \Gamma\left(T^{* 2} M\right), & & L \mapsto\left(d^{2} \varphi\right) L=L \varphi \\
d \varphi \cdot d \psi & \in \Gamma\left(T^{* 2} M\right), & & L \mapsto \Gamma(\varphi, \psi) \equiv \frac{1}{2}[L(\varphi \psi)-\varphi L(\psi)-\psi L(\varphi)]
\end{array}
$$

where $L \in \Gamma\left(T^{2} M\right)$. In particular, for $A, B \in \Gamma(T M)$, it holds then

$$
\begin{align*}
& d^{2} \varphi(A)=A(\varphi), \quad d^{2} \varphi(A \cdot B)=A B(\varphi)  \tag{1.9.2}\\
& d \varphi \cdot d \psi(A)=0, \quad d \varphi \cdot d \psi(A \cdot B)=\frac{1}{2}[A(\varphi) B(\psi)+A(\psi) B(\varphi)] \tag{1.9.3}
\end{align*}
$$

where $d^{2} \varphi$ and $d \varphi \cdot d \psi$ are already determined by (1.9.1). Recall that $A \cdot B \in \Gamma\left(T^{2} M\right)$ denotes the composition of the derivations $A, B$.

For functions $\varphi, \psi \in C^{\infty}(M)$ it is easy to see that

$$
\begin{equation*}
d^{2}(\varphi \psi)=\varphi d^{2} \psi+\psi d^{2} \varphi+2 d \varphi \cdot d \psi \tag{1.9.4}
\end{equation*}
$$

Moreover note that $d^{2} \varphi \mid \Gamma(T M)=d \varphi$ and $d \varphi \cdot d \psi \mid \Gamma(T M)=0$.
Corollary 1.9.6. Let $M$ be a manifold and $(h, U)$ a chart at $x$. Then

$$
\begin{aligned}
& \left(\frac{\partial}{\partial h^{i}}, \frac{\partial^{2}}{\partial h^{i} \partial h^{i}}: 1 \leq i \leq n ; \quad \frac{\partial^{2}}{\partial h^{i} \partial h^{j}}: 1 \leq i<j \leq n\right), \quad \text { respectively, } \\
& \left(d^{2} h^{i}, d h^{i} \cdot d h^{i}: 1 \leq i \leq n ; \quad 2 d h^{i} \cdot d h^{j}: 1 \leq i<j \leq n\right)
\end{aligned}
$$

are frame systems for the second order tangent bundle $T^{2} M$ over $U$, respectively the second order cotangent bundle $T^{* 2} M$ over $U$, which are dual to each other.

Example 1.9.7. Let $L \in \Gamma\left(T^{2} M\right)$ and $\Gamma(\varphi, \psi)=(d \varphi \cdot d \psi) L$ for $\varphi, \psi \in C^{\infty}(M)$. As noted above, we interpret $L$ as PDO of second order via $(L \varphi)(x)=L_{x} \varphi, \varphi \in C^{\infty}(M)$.
(a) If $L=A_{0}+\sum_{i} A_{i}^{2}$ is a PDO in Hörmander form with $A_{0}, A_{i} \in \Gamma(T M)$, then

$$
\Gamma(\varphi, \psi)=\sum_{i}\left(A_{i} \varphi\right)\left(A_{i} \psi\right) .
$$

(b) If $(h, U)$ is a chart such that $L \left\lvert\, U=\sum_{i} b^{i} \frac{\partial}{\partial h^{i}}+\sum_{i, j} a^{i j} \frac{\partial}{\partial h^{i}} \frac{\partial}{\partial h^{j}}\right.$, then

$$
\Gamma(\varphi, \psi) \mid U=\sum_{i, j} a^{i j}\left(\partial_{i} \varphi\right)\left(\partial_{j} \psi\right) \quad \text { where } \partial_{i}=\frac{\partial}{\partial h^{i}} .
$$

More generally, each section $L \in \Gamma\left(T^{2} M\right)$ has a representation as $L=\sum_{\text {finite }} L^{\nu}$, where $L^{\nu}=A \in \Gamma(T M)$ or $L^{\nu}=A \cdot B$ with $A, B \in \Gamma(T M)$. It holds $\Gamma(\varphi, \psi)=0$ for all $\varphi, \psi \in C^{\infty}(M)$ if and only if $L \in \Gamma(T M)$.

Notation 1.9.8. For any two vector fields $A, B \in \Gamma(T M)$ we have $A \cdot B \in \Gamma\left(T^{2} M\right)$, defined as composition of the derivation $A$ and $B$; for two differential forms $\alpha, \beta \in$ $\Gamma\left(T^{*} M\right)$ we have (slightly more general than Remark 1.9.5) $\alpha \cdot \beta \in \Gamma\left(T^{* 2} M\right)$ well defined through $(\alpha \cdot \beta)_{x}:=(d \varphi \cdot d \psi)_{x}$ if $\alpha_{x}=(d \varphi)_{x}$ and $\beta_{x}=(d \psi)_{x}$.

Analogously to Lemma 1.3 .2 we have for differential forms of second order the following Lemma.

Lemma 1.9.9. On any manifold $M$ there exists a finite number of real-valued functions $\varphi^{1}, \ldots, \varphi^{k} \in C^{\infty}(M)$ such that the following properties hold:
(i) Each $\vartheta \in \Gamma\left(T^{* 2} M\right)$ writes as $\vartheta=\sum_{\nu=1}^{k} \vartheta_{\nu} d^{2} \varphi^{\nu}$ where $\vartheta_{\nu} \in C^{\infty}(M)$.
(ii) If $X$ is a continuous semimartingale on $M$, then each continuous adapted $T^{* 2} M$ valued process $\Theta$ over $X$ (i.e., $\Theta_{t} \in T_{X_{t}}^{* 2} M$ for $t \in \mathbb{R}_{+}$) has a representation of the form $\Theta=\sum_{\nu=1}^{k} \Theta_{\nu}\left(d^{2} \varphi^{\nu}\right)(X)$ with continuous adapted real-valued processes $\Theta_{\nu}$.
Proof. The proof proceeds along the usual scheme, as already exploited in the proof of Lemma 1.3.2. We first realize $M$ via a Whitney embedding $h: M \longleftrightarrow \mathbb{R}^{\ell}$ as closed submanifold of a suitable $\mathbb{R}^{\ell}$. There is a partition $\left(\phi_{\lambda}\right)_{\lambda \in \Lambda}$ of the unity on $M$ and a family $\left(I_{\lambda}\right)_{\lambda \in \Lambda}$ of subsets $I_{\lambda} \subset\{1, \ldots, \ell\}$ with the following property: For each $\lambda \in \Lambda$ the components $\left(h^{i}\right)_{i \in I_{\lambda}}$ define a chart for $M$ on an open neighbourhood of $\operatorname{supp}\left(\phi_{\lambda}\right)$.

For (i): By Corollary 1.9.6, we have

$$
\phi_{\lambda} \vartheta=\sum_{i=1}^{\ell} \vartheta_{i}^{\lambda} d^{2} h^{i}+\sum_{i, j=1}^{\ell} \vartheta_{i j}^{\lambda} d h^{i} \cdot d h^{j}
$$

where $\vartheta_{i}^{\lambda}, \vartheta_{i j}^{\lambda} \in C^{\infty}(M)$ are such that $\operatorname{supp}\left(\vartheta_{i}^{\lambda}\right), \operatorname{supp}\left(\vartheta_{i j}^{\lambda}\right) \subset \operatorname{supp}\left(\phi_{\lambda}\right)$ where $\vartheta_{i}^{\lambda}:=0$ for $i \notin I_{\lambda}$ and $\vartheta_{i j}^{\lambda}:=0$ for $\{i, j\} \not \subset I_{\lambda}$. Letting $\tilde{\vartheta}_{i}:=\sum_{\lambda} \vartheta_{i}^{\lambda}$ and $\tilde{\vartheta}_{i j}:=\sum_{\lambda} \vartheta_{i j}^{\lambda}$, this gives the representation

$$
\begin{aligned}
\vartheta & =\sum_{i=1}^{\ell} \tilde{\vartheta}_{i} d^{2} h^{i}+\sum_{i, j=1}^{\ell} \tilde{\vartheta}_{i j} d h^{i} \cdot d h^{j} \\
& \left.=\sum_{i=1}^{\ell} \tilde{\vartheta}_{i} d^{2} h^{i}+\frac{1}{2} \sum_{i, j=1}^{\ell} \tilde{\vartheta}_{i j}\left[d^{2}\left(h^{i} h^{j}\right)-h^{i} d^{2} h^{j}-h^{j} d^{2} h^{i}\right)\right]
\end{aligned}
$$

which shows that $\vartheta$ has a representation of the claimed form. Part (ii) is shown analogously.

We want to come back now to the initial question concerning the status of differentials of $M$-valued semimartingales. For an $M$-valued semimartingale $X$ we first define

$$
(d X)(\varphi):=d(\varphi(X)), \quad \varphi \in C^{\infty}(M)
$$

If $(h, U)$ is a chart for $M$ and $\sigma, \tau$ stopping times with the property that $X \mid[\sigma, \tau[$ takes only values in $U$, then with $\partial_{i}=\frac{\partial}{\partial h^{i}}$ and $X^{i}=h^{i}(X)$ it holds

$$
\begin{equation*}
1_{[\sigma, \tau[ } d(\varphi(X))=1_{[\sigma, \tau[ } \sum_{i}\left(\partial_{i} \varphi\right)(X) d X^{i}+1_{\left[\sigma, \tau\left[\frac{1}{2}\right.\right.} \sum_{i, j}\left(\partial_{i} \partial_{j} \varphi\right)(X) d\left[X^{i}, X^{j}\right] . \tag{1.9.5}
\end{equation*}
$$

As the left-hand side of (1.9.5) is coordinate invariant, also the right-hand side does not depend on the choice of the chart $h$ on $U$, and for $\varphi \in C^{\infty}(M)$ one finds

$$
\begin{aligned}
d(\varphi(X)) & =\sum_{i}\left(\partial_{i} \varphi\right)(X) d X^{i}+\frac{1}{2} \sum_{i, j}\left(\partial_{i} \partial_{j} \varphi\right)(X) d\left[X^{i}, X^{j}\right] \\
& \equiv \sum_{n} 1_{\left[\tau_{n}, \tau_{n+1}[ \right.}\left(\sum_{i}\left(\partial_{i} \varphi\right)(X) d X^{i}+\frac{1}{2} \sum_{i, j}\left(\partial_{i} \partial_{j} \varphi\right)(X) d\left[X^{i}, X^{j}\right]\right),
\end{aligned}
$$

where in the last line we choose to a countable covering of coordinate neighbourhoods the sequence $\left(\tau_{n}\right)_{n \geq 0}$ of stopping times according to Lemma 1.3.1.

Formally we may write this as

$$
d X=\sum_{i}\left(d X^{i}\right) \partial_{i}+\frac{1}{2} \sum_{i, j} d\left[X^{i}, X^{j}\right] \partial_{i} \partial_{j}
$$

from where we can already read off that, at least in a formal sense, the differential $d X$ behaves as a section of $T^{2} M \rightarrow M$ along $X$. The precise meaning of this heuristic argument is given by the following Theorem.

TheOrem 1.9.10 (Principle of Laurent Schwartz). Let $X$ be an $M$-valued semimartingale. There exists exactly one linear mapping

$$
\Theta \longmapsto \int\langle\Theta, d X\rangle
$$

from the real vector space of continuous adapted $T^{* 2} M$-valued processes $\Theta$ over $X$ (i.e., $\Theta_{t} \in T_{X_{t}}^{* 2} M$ for $t \in \mathbb{R}_{+}$) to $\mathscr{S}$ with the following properties:

$$
\begin{align*}
& d^{2} \varphi(X) \mapsto \varphi(X)-\varphi\left(X_{0}\right), \quad \varphi \in C^{\infty}(M)  \tag{1.9.6}\\
& K \Theta \mapsto \int K\langle\Theta, d X\rangle, \quad K \text { continuous, adapted, real-valued process. } \tag{1.9.7}
\end{align*}
$$

where by definition, $\langle\Theta, d X\rangle=d \int\langle\Theta, d X\rangle$.
Notation 1.9.11. We call $\int\langle\Theta, d X\rangle$ the integral of $\Theta$ along $X$. If in particular $\Theta=\vartheta(X)$ for some $\vartheta \in \Gamma\left(T^{* 2} M\right)$, we write also $\int\langle\vartheta, d X\rangle$ instead of $\int\langle\Theta, d X\rangle$.

Proof of Theorem 1.9.10. By Lemma 1.9.9 (ii)the process $\Theta$ has a representation of the form $\Theta=\sum_{\text {finite }} \Theta_{\nu}\left(d^{2} \varphi^{\nu}\right)(X)$; hence necessarily

$$
\begin{equation*}
\int\langle\Theta, d X\rangle=\sum_{\nu} \int \Theta_{\nu} d\left(\varphi^{\nu}(X)\right) \tag{1.9.8}
\end{equation*}
$$

It remains to show that $\int\langle\Theta, d X\rangle$ is well-defined by (1.9.8). Assuming for instance that $\sum_{\text {finite }} K_{\nu}\left(d^{2} \varphi^{\nu}\right)(X)=0$, we have to verify that already $\sum K_{\nu} d\left(\varphi^{\nu}(X)\right)=0$. Without restrictions we may replace here $K_{\nu}$ by $K_{\nu} 1_{[\sigma, \tau[ }$ and assume that $X$ takes on $[\sigma, \tau$ [ only values in the coordinate neighbourhood $U$ of a fixed chart $(h, U)$. In terms of $\varphi^{\nu} \mid U=$ $\bar{\varphi}^{\nu} \circ h$ one observes at first that over $U$

$$
\begin{aligned}
d^{2} \varphi^{\nu} & =d^{2}\left(\bar{\varphi}^{\nu} \circ h\right) \\
& =\sum_{i}\left(D_{i} \bar{\varphi}^{\nu} \circ h\right) d^{2} h^{i}+\sum_{i}\left(D_{i}^{2} \bar{\varphi}^{\nu} \circ h\right) d h^{i} \cdot d h^{i}+\sum_{i<j}\left(D_{i} D_{j} \bar{\varphi}^{\nu} \circ h\right) 2 d h^{i} \cdot d h^{j} \\
& =\sum_{i}\left(D_{i} \bar{\varphi}^{\nu} \circ h\right) d^{2} h^{i}+\sum_{i, j}\left(D_{i} D_{j} \bar{\varphi}^{\nu} \circ h\right) d h^{i} \cdot d h^{j},
\end{aligned}
$$

and hence for $\bar{X}:=h(X)$ by assumption

$$
\begin{aligned}
0 & =\sum_{\nu} K_{\nu}\left(d^{2} \varphi^{\nu}\right)(X) \\
& =\sum_{i} \sum_{\nu} K_{\nu}\left(D_{i} \bar{\varphi}^{\nu}\right)(\bar{X})\left(d^{2} h^{i}\right)(X)+\sum_{i, j} \sum_{\nu} K_{\nu}\left(D_{i} D_{j} \bar{\varphi}^{\nu}\right)(\bar{X})\left(d h^{i} d h^{j}\right)(X),
\end{aligned}
$$

from where we get $\sum_{\nu} K_{\nu}\left(D_{i} \bar{\varphi}^{\nu}\right)(\bar{X}) \equiv 0$ and $\sum_{\nu} K_{\nu}\left(D_{i} D_{j} \bar{\varphi}^{\nu}\right)(\bar{X}) \equiv 0$ almost surely for all $i, j$. On the other hand, this implies

$$
\begin{aligned}
& \sum_{\nu} K_{\nu} d\left(\varphi^{\nu}(X)\right)=\sum_{\nu} K_{\nu} d\left(\bar{\varphi}^{\nu}\right)(\bar{X}) \\
& \quad=\sum_{i}\left(\sum_{\nu} K_{\nu}\left(D_{i} \bar{\varphi}^{\nu}\right)(\bar{X})\right) d \bar{X}^{i}+\frac{1}{2} \sum_{i, j}\left(\sum_{\nu} K_{\nu}\left(D_{i} D_{j} \bar{\varphi}^{\nu}\right)(\bar{X})\right) d\left[\bar{X}^{i}, \bar{X}^{j}\right]=0,
\end{aligned}
$$

which shows the claim.
THEOREM 1.9.12 (Pullback formula). Let $\phi: M \rightarrow N$ be a differentiable map between manifolds and $\Theta$ be a continuous adapted $T^{* 2} N$-valued process over $\phi(X)$ (i.e., $\Theta_{t} \in T_{\phi \circ X_{t}}^{* 2} N$ for $t \in \mathbb{R}_{+}$). Then $\phi^{*} \Theta$ is a $T^{* 2} M$-valued process over $X$ and satisfies

$$
\begin{equation*}
\int\left\langle\phi^{*} \Theta, d X\right\rangle=\int\langle\Theta, d(\phi(X))\rangle \tag{1.9.9}
\end{equation*}
$$

In particular, for $\vartheta \in \Gamma\left(T^{* 2} N\right)$ then $\int\left\langle\phi^{*} \vartheta, d X\right\rangle=\int\langle\vartheta, d(\phi(X))\rangle$.
Proof. Because of $\phi^{*} d^{2} \varphi=d^{2}(\varphi \circ \phi)$, the left-hand side of (1.9.9) has the defining properties of the integral of $\Theta$ along $\phi \circ X$.

Before putting the integral $\int\langle\vartheta, d X\rangle$ of a second order differential form $\vartheta \in \Gamma\left(T^{* 2} M\right)$ along $X$ in perspective to the integrals treated in Section 1.3, e.g. $\int b(d X, d X)$ for $b \in$ $\Gamma\left(T^{*} M \otimes T^{*} M\right)$, respectively $\int_{X} \alpha$ for $\alpha \in \Gamma\left(T^{*} M\right)$, we want to note some further aspects of second order tangent spaces.

REMARK 1.9.13. Denoting for a manifold $M$ by $T M \odot T M$ the vector bundle over $M$ with the symmetric tensor products $T_{x} M \odot T_{x} M$ as fiber $x$, there is a $C^{\infty}(M)$ linear mapping

$$
\begin{equation*}
\wedge: \Gamma\left(T^{2} M\right) \rightarrow \Gamma(T M \odot T M), \quad L \mapsto \hat{L} \tag{1.9.10}
\end{equation*}
$$

determined by

$$
\hat{A}=0, \quad(A \cdot B)^{\wedge}=\frac{1}{2}(A \otimes B+B \otimes A) \equiv A \odot B, \quad A, B \in \Gamma(T M)
$$

Writing $L \in \Gamma\left(T^{2} M\right)$ in a chart $(h, U)$ as $L \mid U=\sum_{i} b^{i} \partial_{i}+\sum_{i, j} a^{i j} \partial_{i} \partial_{j}$ where $\partial_{i}=\frac{\partial}{\partial h^{i}}$, then obviously $\hat{L} \mid U=\sum_{i, j} a^{i j}\left(\partial_{i} \odot \partial_{j}\right)$. Note that the map (1.9.10) is characterized by the property

$$
(d \varphi \odot d \psi) \hat{L}=(d \varphi \cdot d \psi) L=\Gamma(\varphi, \psi), \quad \varphi, \psi \in C^{\infty}(M)
$$

REMARK 1.9.14. For a manifold $M$ we have the following exact sequence of vector bundles over $M$

$$
\begin{equation*}
0 \longrightarrow T M \hookrightarrow T^{2} M \xrightarrow{\wedge} T M \odot T M \longrightarrow 0 . \tag{1.9.11}
\end{equation*}
$$

By dualization of (1.9.11) we obtain the exact sequence

$$
\begin{align*}
& 0 T^{*} M \odot T^{*} M \xrightarrow{H} T^{* 2} M \xrightarrow{R} T^{*} M \longrightarrow 0  \tag{1.9.12}\\
& \alpha_{x} \odot \beta_{x} \longmapsto(\alpha \cdot \beta)_{x}
\end{align*}
$$

where $R$ represents the restriction of $T^{2} M$ to the subbundle $T M$.
TheOrem 1.9.15. For a manifold $M$ we have the exact sequence
(1.9.13) $0 \longrightarrow \Gamma\left(T^{*} M \odot T^{*} M\right) \xrightarrow{H} \Gamma\left(T^{* 2} M\right) \xrightarrow{R} \Gamma\left(T^{*} M\right) \longrightarrow 0$
where $H(\alpha \odot \beta)=\alpha \cdot \beta$, as well as $R\left(d^{2} \varphi\right)=d \varphi$ and $R(d \varphi \cdot d \psi)=0$.
The sequence (1.9.13) of $C^{\infty}(M)$-modules possesses an $\mathbb{R}$-linear splitting

$$
\begin{equation*}
0 \longrightarrow \Gamma\left(T^{*} M \odot T^{*} M\right) \xrightarrow{H} \Gamma\left(T^{* 2} M\right) \xrightarrow{R} \Gamma\left(T^{*} M\right) \longrightarrow 0 ; \tag{1.9.14}
\end{equation*}
$$

More precisely, there exists an $\mathbb{R}$-linear mapping $d: \Gamma\left(T^{*} M\right) \rightarrow \Gamma\left(T^{* 2} M\right)$ with the properties $d(d \varphi)=d^{2} \varphi$ and $d(\varphi \alpha)=d \varphi \cdot \alpha+\varphi \cdot d \alpha$ such that $R \circ d=\mathrm{id}$..

Proof. For $\alpha=\sum_{\nu} \varphi_{\nu} d h^{\nu} \in \Gamma\left(T^{*} M\right)$, we want to verify that

$$
d \alpha:=\sum_{\nu} d \varphi_{\nu} \cdot d h^{\nu}+\sum_{\nu} \varphi_{\nu} d^{2} h^{\nu}
$$

is well-defined. Assume that for instance $\alpha=\sum_{\nu} \varphi_{\nu} d h^{\nu}=0$. We then have to show that

$$
\vartheta:=\sum_{\nu} d \varphi_{\nu} \cdot d h^{\nu}+\sum_{\nu} \varphi_{\nu} d^{2} h^{\nu}=0 .
$$

To this end, it is sufficient to show that $\vartheta(L)=0$ for each section $L \in \Gamma\left(T^{2} M\right)$ where we may assume without restrictions that either $L=A$ or $L=A \cdot B$ with $A, B \in \Gamma(T M)$.
(1) If $L=A \in \Gamma(T M)$, then $\vartheta(A)=\sum_{\nu}\left(d \varphi_{\nu} \cdot d h^{\nu}\right) A+\sum_{\nu}\left(\varphi_{\nu} d^{2} h^{\nu}\right) A=0$ where the first term vanishes, since $d \varphi_{\nu} \cdot d h^{\nu} \mid \Gamma(T M)=0$, while the second term equals $\alpha(A)$ and vanishes since $\alpha=0$ by assumption.
(2) Let now $L=A \cdot B$ where $A, B \in \Gamma(T M)$ : At first we have $\vartheta(A B-B A)=$ $\vartheta([A, B])=0$ by (1), and thus

$$
\begin{aligned}
2 \vartheta(A \cdot B) & =\vartheta(A B+B A)+\vartheta([A, B]) \\
& =\sum_{\nu}\left(d \varphi_{\nu} \cdot d h^{\nu}\right)(A \cdot B+B \cdot A)+\sum_{\nu}\left(\varphi_{\nu} d^{2} h^{\nu}\right)(A \cdot B+B \cdot A) \\
& =\sum_{\nu}\left[\left(A \varphi_{\nu}\right)\left(B h^{\nu}\right)+\left(B \varphi_{\nu}\right)\left(A h^{\nu}\right)+\varphi_{\nu}(A \cdot B)\left(h^{\nu}\right)+\varphi_{\nu}(B \cdot A)\left(h^{\nu}\right)\right] \\
& =\sum_{\nu}\left[A\left(\varphi_{\nu}\left(B h^{\nu}\right)\right)+B\left(\varphi_{\nu}\left(A h^{\nu}\right)\right)\right]=A(\alpha(B))+B(\alpha(A))=0
\end{aligned}
$$

which gives the claim.
Lemma 1.9.16. Let $X$ be a semimartingale taking values in a manifold $M$.
(i) For $\varphi, \psi \in C^{\infty}(M)$ it holds $\int\langle d \varphi \cdot d \psi, d X\rangle=\frac{1}{2}[\varphi(X), \psi(X)]$.
(ii) For $\vartheta, \sigma \in \Gamma\left(T^{* 2} M\right)$ we have $R \vartheta, R \sigma \in \Gamma\left(T^{*} M\right)$, and it holds:

$$
\int\langle R \vartheta \cdot R \sigma, d X\rangle=\frac{1}{2}\left[\int\langle\vartheta, d X\rangle, \int\langle\sigma, d X\rangle\right]
$$

In particular, if $\vartheta \in \Gamma\left(T^{* 2} M\right)$ such that $R \vartheta=0$, then $\int\langle\vartheta, d X\rangle \in \mathscr{A}$.

Proof. To (i): By $d \varphi \cdot d \psi=\frac{1}{2}\left[d^{2}(\varphi \psi)-\varphi d^{2} \psi-\psi d^{2} \varphi\right]$ we have

$$
\begin{aligned}
& 2 \int\langle d \varphi \cdot d \psi, d X\rangle \\
& \quad=\left[(\varphi \psi)(X)-(\varphi \psi)\left(X_{0}\right)\right]-\int \varphi(X) d(\psi(X))-\int \psi(X) d(\varphi(X)) \\
& \quad=[\varphi(X), \psi(X)]
\end{aligned}
$$

To (ii): According to Lemma 1.9.9, $\vartheta$ and $\sigma$ have representations of the form $\vartheta=\sum_{\nu} \vartheta_{\nu} d^{2} \varphi^{\nu}$, respectively $\sigma=\sum_{\mu} \sigma_{\mu} d^{2} \psi^{\mu}$. Hence we have

$$
R \vartheta=\sum_{\nu} \vartheta_{\nu} d \varphi^{\nu}, \quad R \sigma=\sum_{\mu} \sigma_{\mu} d \psi^{\mu}, \quad R \vartheta \cdot R \sigma=\sum_{\nu, \mu} \vartheta_{\nu} \sigma_{\mu} d \varphi^{\nu} \cdot d \psi^{\mu}
$$

and by means of (i) we obtain

$$
\begin{aligned}
\int\langle R \vartheta \cdot R \sigma, d X\rangle & =\frac{1}{2} \sum_{\nu, \mu} \int \vartheta_{\nu}(X) \sigma_{\mu}(X) d\left[\varphi^{\nu}(X), \psi^{\mu}(X)\right] \\
& =\frac{1}{2}\left[\sum_{\nu} \int \vartheta_{\nu}(X) d\left(\varphi^{\nu}(X)\right), \sum_{\mu} \int \sigma_{\mu}(X) d\left(\psi^{\mu}(X)\right)\right] \\
& =\frac{1}{2}\left[\int\langle\vartheta, d X\rangle, \int\langle\sigma, d X\rangle\right]
\end{aligned}
$$

The additional claim follows from part (ii) with $\vartheta=\sigma$.
Theorem 1.9.17. Let $X$ be an $M$-valued semimartingale.
(i) For $b \in \Gamma\left(T^{*} M \odot T^{*} M\right)$ it holds that $\int\langle H b, d X\rangle=\frac{1}{2} \int b(d X, d X)$.
(ii) For $\alpha \in \Gamma\left(T^{*} M\right)$ it holds that $\int\langle d \alpha, d X\rangle=\int_{X} \alpha$.

Proof. It is sufficient to verify the defining properties.
To (i): For $b=d \varphi \odot d \psi$ we have $H b=d \varphi \cdot d \psi$ and by Lemma 1.9.16 (i) then $\int\langle H b, d X\rangle=\frac{1}{2}[\varphi(X), \psi(X)]$. On the other hand, we have $H(\varphi b)=\varphi H(b)$ from where the relation $\int\langle H(\varphi b), d X\rangle=\int \varphi(X) b(d X, d X)$ follows.

To (ii): If $\alpha=d \varphi$, then $d(d \varphi)=d^{2} \varphi$ and hence $\int\langle d \alpha, d X\rangle=\varphi(X)-\varphi\left(X_{0}\right)$. On the other hand we have $d(\varphi \alpha)=\varphi \cdot d \alpha+d \varphi \cdot \alpha$, and hence by means of Lemma 1.9.16 (ii), applied to $d \varphi \cdot \alpha=R\left(d^{2} \varphi\right) \cdot R(d \alpha)$,

$$
\begin{aligned}
\int\langle d(\varphi \alpha), d X\rangle & =\int\langle\varphi \cdot d \alpha, d X\rangle+\int\langle d \varphi \cdot \alpha, d X\rangle \\
& =\int \varphi(X)\langle d \alpha, d X\rangle+\frac{1}{2}\left[\int\left\langle d^{2} \varphi, d X\right\rangle, \int\langle d \alpha, d X\rangle\right] \\
& =\int \varphi(X) d\left(\int\langle d \alpha, d X\rangle\right)+\frac{1}{2}\left[\varphi(X), \int\langle d \alpha, d X\rangle\right] \\
& =\int \varphi(X) \circ d\left(\int\langle d \alpha, d X\rangle\right)
\end{aligned}
$$

which shows the defining properties and hence the claim.
EXAMPLE 1.9.18. Let $X$ be a $M$-valued semimartingale of locally bounded variation, in the sense that all compositions $\varphi(X)$ with functions $\varphi \in C^{\infty}(M)$ lie in $\mathscr{A}$. Then we have

$$
\begin{equation*}
\int\langle\Theta, d X\rangle=\int\langle d R \Theta, d X\rangle=\int(R \Theta)(\circ d X) \tag{1.9.15}
\end{equation*}
$$

for each continuous adapted $T^{* 2} M$-valued process $\Theta$ over $X$. Indeed, both $\int(R \Theta)(\circ d X)$ and $\int\langle d R \Theta, d X\rangle$ have the defining properties for $\int\langle\Theta, d X\rangle$. In particular, we then have

$$
\int\langle\vartheta, d X\rangle=\int_{X} R \vartheta, \quad \vartheta \in \Gamma\left(T^{* 2} M\right)
$$

a formula which uncovers why in classical Differential Geometry second order forms do not appear explicitly.

THEOREM 1.9.19. Let $M$ be a manifold. There is a one-to-one correspondence between the following objects:
(i) Torsion-free linear connections $\nabla$ on $M$.
(ii) Bundle homomorphisms $F: T^{2} M \rightarrow T M$ with $F \circ \iota=$ id (where $\iota$ denotes the canonical inclusion $T M \longleftrightarrow T^{2} M$ ), i.e., splittings of the following exact sequence of vector bundles over $M$
(iii) Bundle homomorphisms $G: T^{*} M \rightarrow T^{* 2} M$ with $R \circ G=\mathrm{id}$ (where $R$ denotes the restriction to $T M$ ), i.e., splittings of the following exact sequence of vector bundles over $M$

$$
\begin{equation*}
0 \longrightarrow T^{*} M \odot T^{*} M \xrightarrow{H} T^{* 2} \underset{\text { 久_-_G_-, }}{M} \xrightarrow{R} T^{*} M \longrightarrow 0 . \tag{1.9.17}
\end{equation*}
$$

Proof. Obviously (ii) and (iii) correspond to each other by dualization.
(i) $\rightarrow$ (ii): Let $\nabla$ be a torsion-free linear connection on $M$. Recall that by 1.4.30 torsion-freeness means that for any $\varphi \in C^{\infty}(M),(A, B) \mapsto \nabla d \varphi(A, B)$ is symmetric. We define $F: \Gamma\left(T^{2} M\right) \rightarrow \Gamma(T M)$ by

$$
\begin{equation*}
(F L)(\varphi)=L \varphi-\langle H \nabla d \varphi, L\rangle, \quad L \in \Gamma\left(T^{2} M\right), \varphi \in C^{\infty}(M) \tag{1.9.18}
\end{equation*}
$$

where $\langle\cdot, \cdot\rangle$ denotes the dual pairing of $T^{* 2} M$ and $T^{2} M$. Obviously, it holds that

$$
F L= \begin{cases}L & \text { for } L=A \in \Gamma(T M) \\ \nabla_{A} B & \text { for } L=A \cdot B \text { with } A, B \in \Gamma(T M)\end{cases}
$$

Indeed for $L=A \in \Gamma(T M)$ we have $\langle H \nabla d \varphi, A\rangle=0$; on the other hand for $L=A \cdot B$ we have

$$
F(A \cdot B)(\varphi)=(A \cdot B)(\varphi)-\langle H \nabla d \varphi, A \cdot B\rangle=(A \cdot B)(\varphi)-\nabla d \varphi(A, B)=\left(\nabla_{A} B\right)(\varphi)
$$

where we used that

$$
\langle H b, A \cdot B\rangle=\frac{1}{2}(b(A, B)+b(B, A))=b(A, B)
$$

for each $b \in \Gamma\left(T^{*} M \odot T^{*} M\right)$. This shows in particular that $F: \Gamma\left(T^{2} M\right) \rightarrow \Gamma(T M)$ is by Eq. (1.9.18) well-defined and that $F \mid \Gamma(T M)=\mathrm{id}$; moreover $F$ is $C^{\infty}(M)$-linear and hence defines a bundle homomorphism $F: T^{2} M \rightarrow T M$ with the wanted properties.
(ii) $\rightarrow$ (i): Now let be given a bundle homomorphism $F: T^{2} M \rightarrow T M$ such that $F \mid T M=$ id, and let $F: \Gamma\left(T^{2} M\right) \rightarrow \Gamma(T M)$ be the induced mapping at the level of sections. Inversely to (1.9.18), $F$ induces a linear connection $\nabla$ on $M$, namely as

$$
\begin{equation*}
\nabla_{A} B:=F(A \cdot B), \quad A, B \in \Gamma(T M) \tag{1.9.19}
\end{equation*}
$$

$\nabla$ is obviously $C^{\infty}(M)$-linear in $A$ and derivative in $B$, since $\nabla_{A}(\varphi B)=F(A \cdot(\varphi B))=$ $F(\varphi A \cdot B+A(\varphi) B)=\varphi F(A \cdot B)+A(\varphi) F(B)=\varphi \nabla_{A} B+A(\varphi) B$. Moreover, we observe that

$$
\nabla_{A} B-\nabla_{B} A=F(A B)-F(B A)=F(A B-B A)=F([A, B])=[A, B],
$$

which shows that $\nabla$ is torsion-free.

By symmetrization, the $C^{\infty}(M)$-linear map

$$
H: \Gamma\left(T^{*} M \odot T^{*} M\right) \rightarrow \Gamma\left(T^{* 2} M\right), \quad H(\alpha \odot \beta)=\alpha \cdot \beta
$$

can be extended to a $C^{\infty}(M)$-linear mapping

$$
H: \Gamma\left(T^{*} M \otimes T^{*} M\right) \rightarrow \Gamma\left(T^{* 2} M\right), \quad H(\alpha \otimes \beta):=H(\alpha \odot \beta)=\alpha \cdot \beta
$$

REMARK 1.9.20. Explicitly, the bundle homomorphism $G: T^{*} M \rightarrow T^{* 2} M$, induced from a torsion-free linear connection $\nabla$ on $M$ by Theorem 1.9.19 (iii), is given by

$$
\begin{equation*}
G \alpha=d \alpha-H \nabla \alpha, \quad \alpha \in \Gamma\left(T^{*} M\right) \tag{1.9.20}
\end{equation*}
$$

Proof. By construction, $G$ is determined by

$$
\langle G \alpha, L\rangle=\langle\alpha, F L\rangle:=\alpha(F L), \quad \alpha \in \Gamma\left(T^{*} M\right), L \in \Gamma\left(T^{2} M\right),
$$

with $F L$ being defined by Eq. (1.9.18). Since by $\nabla(\varphi \alpha)=d \varphi \otimes \alpha+\varphi \nabla \alpha$ the right-hand side of (1.9.20) is $C^{\infty}(M)$-linear in $\alpha$, it is sufficient to show Eq. (1.9.20) for $\alpha=d \varphi$ with $\varphi \in C^{\infty}(M)$. We have however

$$
\langle G d \varphi, L\rangle=\langle d \varphi, F L\rangle=(F L)(\varphi)=\left\langle d^{2} \varphi-H \nabla d \varphi, L\right\rangle, \quad L \in \Gamma\left(T^{2} M\right)
$$

from where the relation $G d \varphi=d^{2} \varphi-H \nabla d \varphi$ follows.
REMARK 1.9.21. Theorem 1.9.19 shows explicitly that torsion-free linear connections on $M$ are exactly the required extra structure to split differential operators $L \in \Gamma\left(T^{2} M\right)$ of second order canonically in a first order part (the drift of $L$ ), namely $F L$, and a part of purely second order, namely $L-F L=L-\iota F L$. We call a PDO $L \in \Gamma\left(T^{2} M\right)$ of purely second order if $F L \equiv 0$. Writing $L$ in a chart $(h, U)$ as $L \left\lvert\, U=\sum_{i} b^{i} \frac{\partial}{\partial h^{i}}+\right.$ $\sum_{i, j} a^{i j} \frac{\partial}{\partial h^{i}} \frac{\partial}{\partial h^{j}}$, then

$$
(F L) \left\lvert\, U=\sum_{k}\left(b^{k}+\sum_{i, j} a^{i j} \Gamma_{i j}^{k}\right) \frac{\partial}{\partial h^{k}}\right.
$$

is the corresponding first order part.
THEOREM 1.9.22. Let $M$ be a manifold, $\nabla$ a torsion-free linear connection on $M$ and $\gamma$ a differentiable curve taking values in $M$. The following conditions are equivalent:
(i) $\gamma$ is a geodesic curve.
(ii) $\ddot{\gamma} \equiv \gamma_{*}\left(\frac{d^{2}}{d t^{2}}\right)$ is purely second order, i.e., $F(\ddot{\gamma}) \equiv 0$.

Proof. Letting $(h, U)$ be a chart for $M$, then we have for $t \in \mathbb{R}$ such that $\gamma(t) \in U$

$$
\begin{aligned}
\dot{\gamma}(t) & =\sum_{i} \dot{\gamma}^{i}(t)\left(\frac{\partial}{\partial h^{i}}\right)_{\gamma(t)}, \\
\ddot{\gamma}(t) & =\sum_{i} \ddot{\gamma}^{i}(t)\left(\frac{\partial}{\partial h^{i}}\right)_{\gamma(t)}+\sum_{i, j} \dot{\gamma}^{i}(t) \dot{\gamma}^{j}(t)\left(\frac{\partial}{\partial h^{i}}\right)_{\gamma(t)}\left(\frac{\partial}{\partial h^{j}}\right)_{\gamma(t)},
\end{aligned}
$$

and consequently

$$
F(\ddot{\gamma}(t))=\sum_{k}\left(\ddot{\gamma}^{k}(t)+\sum_{i, j} \Gamma_{i j}^{k}(\gamma(t)) \dot{\gamma}^{i}(t) \dot{\gamma}^{j}(t)\right)\left(\frac{\partial}{\partial h^{k}}\right)_{\gamma(t)}
$$

from where we read off the claim with Eq. (1.4.10).

THEOREM 1.9.23. Let $M$ be a manifold and let $L: C^{\infty}(M) \rightarrow C^{\infty}(M)$ be a second order PDO without constant term, i.e. $L \in \Gamma\left(T^{2} M\right)$ where $(L \varphi)(x)=L_{x}(\varphi)$. Furthermore, let $X$ be an $M$-valued semimartingale which solves the martingale problem for $L$, i.e.,

$$
\begin{equation*}
d(\varphi(X))-(L \varphi)(X) d t \stackrel{\mathrm{~m}}{=} 0 \quad \text { for each } \varphi \in C^{\infty}(M) \tag{1.9.21}
\end{equation*}
$$

For a torsion-free linear connection $\nabla$ on $M$, the following two conditions are equivalent:
(i) $X$ is a $\nabla$-martingale on $M$
(ii) $L$ is of purely second order along $X$, i.e., $(F L)(X) \equiv 0 \mathbb{P}$-almost surely.

Proof. We want to check first of all that the property (1.9.21) implies

$$
\begin{equation*}
\langle\vartheta, d X\rangle \stackrel{\mathrm{m}}{=}\langle\vartheta(X), L(X)\rangle d t, \quad \vartheta \in \Gamma\left(T^{* 2} M\right) \tag{1.9.22}
\end{equation*}
$$

Indeed, any $\vartheta \in \Gamma\left(T^{* 2} M\right)$ has by Lemma 1.9.9 a representation of the form $\vartheta=\sum_{\nu} \vartheta_{\nu} d^{2} \varphi^{\nu}$ where $\vartheta_{\nu} \in C^{\infty}(M)$, and hence

$$
\langle\vartheta, d X\rangle=\sum_{\nu}\left(\vartheta_{\nu} \circ X\right) d\left(\varphi^{\nu}(X)\right) \stackrel{\mathrm{m}}{=} \sum_{\nu} \vartheta_{\nu}(X) L \varphi^{\nu}(X) d t=\langle\vartheta, L\rangle(X) d t
$$

Recall that by definition $X$ is a $\nabla$-martingale if

$$
d(\varphi(X)) \stackrel{\mathrm{m}}{=} \frac{1}{2}(\nabla d \varphi)(d X, d X), \quad \varphi \in C^{\infty}(M)
$$

By Theorem 1.9.17 (i) we have $\nabla d \varphi(d X, d X)=2\langle H \nabla d \varphi, d X\rangle$, and from relation (1.9.22) we get $\langle H \nabla d \varphi, d X\rangle \stackrel{\mathrm{m}}{=}\langle H \nabla d \varphi, L(X)\rangle d t$, so that $X$ is $\nabla$-martingale if and only if

$$
d(\varphi(X))-\langle(H \nabla d \varphi)(X), L(X)\rangle d t \stackrel{\mathrm{~m}}{=} 0, \quad \varphi \in C^{\infty}(M)
$$

Using $d(\varphi(X)) \stackrel{\mathrm{m}}{=}(L \varphi)(X) d t$ we conclude that $X$ is a $\nabla$-martingale if and only if for each $\varphi \in C^{\infty}(M)$ :

$$
(L \varphi-\langle H \nabla d \varphi, L\rangle)(X) d t=0
$$

which is because of $L \varphi-\langle H \nabla d \varphi, L\rangle=F L$ just the claim.
We finally want use relation (1.6.34) to define the Itô integral of one-forms along semimartingales in a more general context.

DEFINITION 1.9.24 (Itô integral along semimartingales). Let $M$ be a manifold, $\nabla$ a torsion-free linear connection on $M$ and $X$ a semimartingale taking values in $M$. For a $T^{*} M$-valued process $J$ over $X$, we call

$$
\int\langle J, F d X\rangle:=\int\langle G J, d X\rangle
$$

the Itô integral of $J$ along $X$. If in particular $J=\alpha(X)$ where $\alpha \in \Gamma\left(T^{*} M\right)$, then $\int\langle\alpha(X), F d X\rangle$ is also called Itô integral of $\alpha$ along $X$ and the following notations are used for it:

$$
(\nabla) \int_{X} \alpha=\int\langle\alpha, F d X\rangle=\int\langle G \alpha, d X\rangle
$$

REMARK 1.9.25. By (1.9.20) we have $G \alpha=d \alpha-H \nabla \alpha$ and $J \mapsto I_{J}:=\int\langle G J, d X\rangle$ is hence determined by the following properties:
(i) $I_{d \varphi(X)}=\varphi(X)-\varphi\left(X_{0}\right)-\frac{1}{2} \int \nabla d \varphi(d X, d X) \quad$ for $\varphi \in C^{\infty}(M)$;
(ii) $I_{K J}=\int K d I_{J} \quad$ for each continuous adapted $\mathbb{R}$-valued process $K$.

In particular, for each differential form $\alpha \in \Gamma\left(T^{*} M\right)$ the following relation between the Itô Integral and Stratonovich integral of $\alpha$ along $X$ holds:

$$
\begin{equation*}
(\nabla) \int_{X} \alpha=\int_{X} \alpha-\frac{1}{2} \int \nabla \alpha(d X, d X) \tag{1.9.23}
\end{equation*}
$$

THEOREM 1.9.26. Let $M$ be a manifold, $\nabla$ a torsion-free linear connection on $M$ and $X$ a $M$-valued semimartingale. The following statements are equivalent:
(i) $X$ is a $\nabla$-martingale;
(ii) $(\nabla) \int_{X} \alpha$ is a local martingale for any differential form $\alpha \in \Gamma\left(T^{*} M\right)$;
(iii) $\int\langle J, F d X\rangle$ is a local martingale for any continuous adapted $T^{*} M$-valued process $J$ above $X$.

Proof. For $\alpha=d \varphi$, respectively $J=\alpha(X)$, the assertions reduce to the definition of $\nabla$-martingales. The general case follows with Lemma 1.9.9.

The following Remark finally justifies the notion Itô integral, respectively Stratonovich integral of $\alpha$ along $X$.

THEOREM 1.9.27. Let $M$ be a manifold, $\nabla$ a torsion-free linear connection on $M$ and $X$ an $M$-valued semimartingale. Furthermore let $U$ be a horizontal lift of $X$ to $\mathrm{L}(T M)$ and $Z=\int_{U} \vartheta$ the anti-development of $X$ in $\mathbb{R}^{n}$. Then for the Itô integral, respectively Stratonovich integral of a differential form $\alpha \in \Gamma\left(T^{*} M\right)$ along $X$ the following formulas hold:

$$
\int_{X} \alpha=\sum_{i} \int \alpha(X) U e_{i} \circ d Z^{i}, \quad \text { as well as } \quad(\nabla) \int_{X} \alpha=\sum_{i} \int \alpha(X) U e_{i} d Z^{i}
$$

Proof. The first formula is already shown in Theorem 1.6.30 (ii); the second one reduces for $\alpha=d \varphi$ with Eq. (1.9.23) to the geometric Itô formula (1.6.32); the general case follows again with Lemma 1.9.9.

## CHAPTER 2

## Geometry of Brownian Motion

In this Chapter we focus on stochastic tools in Riemannian Geometry. We start by studying some questions concerning the geometry of Riemannian manifolds in connection with the long-term behaviour of Brownian motion. In particular, using a few selected problems, we want to illustrate the basic idea of stochastic Riemannian geometry, namely to relate differential geometric problems to stochastic questions and to deal with them using stochastic methods (see for instance, [32] and [23, 24]).

### 2.1. The Curvature Tensor and Jacobi Fields

The notion of curvature of a Riemannian manifold is one of the key concepts to control the asymptotic behaviour of Brownian motions for large times. Brownian motion is a sensitive instrument to measure curvature. Negative curvature amplifies the tendency of Brownian motion to exit compact sets and to drift off to $\infty$ if the topology of the manifold permits. Strongly divergent negative curvature, for instance, can have the effect that even on metrically complete manifolds $\mathrm{BM}(M, g)$ explodes in finite times.

Brownian motion on Riemannian manifolds can have other asymptotic properties not known from Euclidean Brownian motion. For instance, trajectories of $\mathrm{BM}(M, g)$ on certain negatively curved simply connected manifolds $M$, when considered in polar coordinates from some fixed point, stay with high probability in the entered angular sector and have an asymptotic direction.

Before discussing the concept of curvature we want to consider Riemannian manifolds under the aspect of metric spaces. Recall that the $d(x, y)$ of two points $x$ and $y$ is given by

$$
d(x, y):=\inf \left\{L(\alpha) \mid \alpha:[0,1] \rightarrow M \text { piecewise } C^{\infty} \text { with } \alpha(0)=x \text { and } \alpha(1)=y\right\}
$$

One of the most elementary questions about asymptotics of Brownian motions is the distance behaviour of $\mathrm{BM}(M, g)$ with respect to a given point $x \in M$, i.e. properties of "radial process" $d\left(x, X_{t}\right), t \geq 0$.

EXAMPLE 2.1.1 (Bessel process). Let $M=\mathbb{R}^{n}$ endowed with the Euclidean metric. The function $r(x):=d(0, x)=|x|$ is $C^{\infty}$ on $\mathbb{R}^{n} \backslash\{0\}$ with

$$
(\operatorname{grad} r)(x)=\frac{x}{|x|} \quad \text { and } \quad \Delta r=\frac{n-1}{r}
$$

Consequently, if $X$ is a $\operatorname{BM}\left(\mathbb{R}^{n}\right)$ such that $X_{0} \neq 0$ a.s. ( $n \geq 2$ ), then by Itô's formula

$$
d(r(X))=d N+d A=\sum_{i=1}^{n} \frac{X^{i}}{|X|} d X^{i}+\frac{1}{2} \frac{n-1}{r(X)} d t
$$

As $d[N, N]=\sum_{i, j} X^{i} X^{j}|X|^{-2} d X^{i} d X^{j}=\sum_{i, j} X^{i} X^{j}|X|^{-2} \delta_{i j} d t=d t$, the process $N$ is a one-dimensional Brownian motion $W$, and one gets

$$
r(X)=r\left(X_{0}\right)+W+\frac{1}{2} \int_{0}^{t} \frac{n-1}{r(X)} d s \quad \text { with } W \text { a } \operatorname{BM}(\mathbb{R})
$$

In other words, $r(X)$ is solution of an SDE of the type

$$
\begin{equation*}
d R=d W+\frac{n-1}{2} R^{-1} d t . \tag{2.1.1}
\end{equation*}
$$

Any such process is called $n$-dimensional Bessel process or Bessel process of index $n / 2-1$.
Example 2.1.1 rises the question to what extent in general for Brownian motions $X$ on $(M, g)$ and $x \in M$, radial processes of the form $r(X):=d(x, X)$ are semimartingales which can be described by Itô's formula. To this end, first questions concerning differentiability of the distance function $r(\cdot):=d(x, \cdot)$ on $M$ need to be clarified.

DEFINITION 2.1.2 (Variation of a curve). Let $(M, g)$ be a Riemannian manifold and $\gamma:[a, b] \rightarrow M$ a non-constant differentiable curve parametrized proportionally to arc length, i.e. $0<\ell:=|\dot{\gamma}|=$ const. A (free) variation of $\gamma$ is a differentiable map

$$
\alpha:[a, b] \times]-\varepsilon, \varepsilon[\rightarrow M
$$

such that $\alpha(\cdot, 0)=\gamma$.
In terms of the canonical vector fields $\frac{\partial}{\partial t}=D_{1}$ and $\frac{\partial}{\partial s}=D_{2}$ on $\left.[a, b] \times\right]-\varepsilon, \varepsilon[$, we consider to a variation $\alpha$ of $\gamma$ the vector fields along $\alpha$ :

$$
T:=\alpha_{*} D_{1} \equiv \frac{\partial}{\partial t} \alpha \in \Gamma\left(\alpha^{*} T M\right) \quad \text { and } \quad V:=\alpha_{*} D_{2} \equiv \frac{\partial}{\partial s} \alpha \in \Gamma\left(\alpha^{*} T M\right)
$$

Let $\gamma_{s}:=\alpha(\cdot, s)$ for $-\varepsilon<s<\varepsilon$; in particular $\gamma=\gamma_{0}$. Furthermore denote by $Y:=$ $V(\cdot, 0) \in \Gamma\left(\gamma^{*} T M\right)$ the "variational field" of $\alpha$ and by $\dot{\gamma}=T(\cdot, 0) \in \Gamma\left(\gamma^{*} T M\right)$ the tangential vector field along $\gamma$.

THEOREM 2.1.3 (First variation of arc length). Let $(M, g)$ be a Riemannian manifold with the Levi-Civita connection and let $\alpha:[a, b] \times]-\varepsilon, \varepsilon[\rightarrow M$ be the differentiable variation of a smooth curve $\gamma:[a, b] \rightarrow M$ such that $\ell=|\dot{\gamma}|=$ const $>0$. Let $Y=\left(\alpha_{*} D_{2}\right)(\cdot, 0) \in \Gamma\left(\gamma^{*} T M\right)$. Then the lengths $L\left(\gamma_{s}\right)=\int_{a}^{b}\left|\dot{\gamma}_{s}(t)\right| d t$ of the curves $\gamma_{s}=\alpha(\cdot, s):[a, b] \rightarrow M$ satisfy the "variational formula":

$$
\begin{equation*}
\left.\frac{d}{d s}\right|_{s=0} L\left(\gamma_{s}\right)=\frac{1}{\ell}\left\{\left.\langle Y, \dot{\gamma}\rangle\right|_{t=a} ^{t=b}-\int_{a}^{b}\left\langle Y, \nabla_{D} \dot{\gamma}\right\rangle d t\right\} . \tag{2.1.2}
\end{equation*}
$$

Proof. By Theorem 1.5.6 (iii) (on the characterization of Riemannian connections) we have $\frac{\partial}{\partial s}\langle T, T\rangle \equiv D_{2}\langle T, T\rangle=2\left\langle\nabla_{D_{2}} T, T\right\rangle$, and thus

$$
\begin{align*}
\frac{d}{d s} L\left(\gamma_{s}\right) & =\frac{d}{d s} \int_{a}^{b}\langle T, T\rangle^{1 / 2} d t=\int_{a}^{b} \frac{1}{2}\langle T, T\rangle^{-1 / 2} \frac{\partial}{\partial s}\langle T, T\rangle d t  \tag{2.1.3}\\
& =\int_{a}^{b} \frac{1}{|T|}\left\langle\nabla_{D_{2}} T, T\right\rangle d t=\int_{a}^{b} \frac{1}{|T|}\left\langle\nabla_{D_{1}} V, T\right\rangle d t  \tag{2.1.4}\\
& =\int_{a}^{b} \frac{1}{|T|}\left[D_{1}\langle V, T\rangle-\left\langle V, \nabla_{D_{1}} T\right\rangle\right] d t \tag{2.1.5}
\end{align*}
$$

Here the second to the last equality is a consequence of the first structural equation of Cartan (Theorem 1.4.27): $\nabla_{D_{1}}\left(\alpha_{*} D_{2}\right)-\nabla_{D_{2}}\left(\alpha_{*} D_{1}\right)=\alpha_{*}\left[D_{1}, D_{2}\right]=0$, whereas the
last equality comes from $D_{1}\langle V, T\rangle=\left\langle\nabla_{D_{1}} V, T\right\rangle+\left\langle V, \nabla_{D_{1}} T\right\rangle$ which is a consequence of Theorem 1.5.6 (iii).

For $s=0$ this gives the wanted equation

$$
\begin{equation*}
\left.\frac{d}{d s}\right|_{s=0} L\left(\gamma_{s}\right)=\frac{1}{\ell}\left\{\left.\langle V, \dot{\gamma}\rangle(t, 0)\right|_{t=a} ^{t=b}-\int_{a}^{b}\left\langle V, \nabla_{D} \dot{\gamma}\right\rangle(t, 0) d t\right\} \tag{2.1.6}
\end{equation*}
$$

In the particular case of a variation $\left(\gamma_{s}\right)_{-\varepsilon<s<\varepsilon}$ of $\gamma$ with fixed initial and end point, i.e. $\gamma_{s}(a)=\gamma(a)$ and $\gamma_{s}(b)=\gamma(b)$ for all $s$, we get $V_{(a, \cdot)}=0$ and $V_{(b, \cdot)}=0$, which combined with Eq. (2.1.2) gives the following characterization of geodesic curves:

Corollary 2.1.4 (Geodesics as critical points of the length functional). Let ( $M, g$ ) be a Riemannian manifold and $\gamma:[a, b] \rightarrow M$ a differentiable curve. Then $\gamma$ is a geodesic curve, i.e. $\nabla_{D} \dot{\gamma}=0$, if and only if $\left.\frac{d}{d s}\right|_{s=0} L\left(\gamma_{s}\right)=0$ for all variations $\left(\gamma_{s}\right)$ of $\gamma$ with fixed initial and end point.

We discuss first some local properties of geodesic curves. On a Riemannian manifold $(M, g)$, to each $x \in M$ and $v \in T_{x} M$ there exists exactly one geodesic curve $\gamma_{v}$ with $\gamma_{v}(0)=x$ and $\dot{\gamma}_{v}(0)=v$; since $\left|\dot{\gamma}_{v}\right| \equiv|v|$ the geodesic $\gamma_{v}$ is a normal geodesic in the sense that $\left|\dot{\gamma}_{v}\right| \equiv 1$, if and only if $|v|=1$.

We consider $\mathscr{O}(M):=\left\{v \in T M: \gamma_{v}\right.$ is defined for $\left.t=1\right\}$ and the exponential map of $(M, g)$

$$
\exp : \mathscr{O}(M) \rightarrow M \times M, \quad v \mapsto\left(\pi(v), \gamma_{v}(1)\right)
$$

As a consequence of the theory of ordinary differential equations, $\mathscr{O}(M)$ is open in $T M$, and contains obviously the zero section in $T M$. By Definition 1.8.10 the Riemannian manifold $(M, g)$ is metrically complete if and only if $\mathscr{O}(M)=T M$. As already observed (see the proof to Lemma 1.7.22), the differential of

$$
\exp _{x}=\exp \mid\left(T_{x} M \cap \mathscr{O}(M)\right): T_{x} M \cap \mathscr{O}(M) \rightarrow\{x\} \times M \equiv M
$$

at the zero element $0_{x} \in T_{x} M$ is given by the identity. In general, along the zero section of a vector bundle $E$ over $M$, the tangent spaces $T_{0_{x}} E$ decompose canonically as $T_{0_{x}} E=T_{x} M \oplus E_{x}$ with $E_{x}$ the part in fiber direction, and $(d \exp )_{0_{x}}$ read as map $(d \exp )_{0_{x}}: T_{x} M \oplus T_{x} M \rightarrow T_{x} M \oplus T_{x} M$ given by the matrix

$$
\left(\begin{array}{cc}
\text { id } & 0 \\
\text { id } & \text { id }
\end{array}\right) .
$$

Hence $d \exp$ has full rank at the zero section and by the inverse function theorem exp maps an open neighbourhood of the zero section in $T M$ locally diffeomorphically to an open neighbourhood of the diagonal in $M \times M$. For $x \in M$ let $V_{\varepsilon}(0):=\left\{v \in T_{x} M:|v|<\varepsilon\right\}$, then

$$
\left.\left.\varrho(x):=\sup \left\{\varepsilon>0: \exp _{x} \mid V_{\varepsilon}(0) \text { is an embedding }\right\} \in\right] 0, \infty\right]
$$

is called injectivity radius at $x$, and hence for $0<\varepsilon \leq \varrho(x)$ there are diffeomorphisms

$$
\exp _{x} \mid V_{\varepsilon}(0): V_{\varepsilon}(0) \xrightarrow{\sim} \exp _{x}\left(V_{\varepsilon}(0)\right)=: B_{\varepsilon}(x)
$$

Note that the map $\varrho: M \rightarrow \overline{\mathbb{R}}$ is lower semi-continuous: $\{\varrho>c\}$ is open in $M$ for any $c \geq 0$.

For $\varepsilon<\varrho(x)$ we may consider besides the normal coordinates $\left(h, B_{\varepsilon}(x)\right)$ at $x$ where $h=\left(\exp _{x} \mid V_{\varepsilon}(0)\right)^{-1}$ the so-called "geodesic polar coordinates" with center $x$, which is
the diffeomorphism

$$
\begin{align*}
B_{\varepsilon}(x) \backslash\{x\} & \left.\sim V_{\varepsilon}(0) \backslash\{0\} \longrightarrow\right] 0, \varepsilon\left[\times S^{n-1}\right. \\
x & \longmapsto\left(|x|, \frac{x}{|x|}\right) . \tag{2.1.7}
\end{align*}
$$

The inverse map to (2.1.7) is given by $\phi:(r, v) \mapsto \exp _{x}(r v)=\gamma_{v}(r)$ Note that through

$$
S_{r}(x):=\left\{\exp _{x}(v): v \in T_{x} M,|v|=r\right\} \equiv \phi\left(\{r\} \times S^{n-1}\right), \quad 0<r<\varrho(x)
$$

then hypersurfaces in $M$ (one-codimensional submanifolds) are given.
THEOREM 2.1.5 (Gauss Lemma). Let $(M, g)$ be a Riemannian manifold equipped with the Levi-Civita connection, $x \in M$ and $v \in T_{x} M \cap \mathscr{O}(M)$ such that $\exp _{x} v$ is defined. For any $w \in T_{x} M \cong T_{v}\left(T_{x} M\right)$ then

$$
\begin{equation*}
\left\langle\left(d \exp _{x}\right)_{v} v,\left(d \exp _{x}\right)_{v} w\right\rangle=\langle v, w\rangle \tag{2.1.8}
\end{equation*}
$$

In particular, the geodesics through the point $x$ are perpendicular on the hypersurfaces $S_{r}(x)$ for $0<r<\varrho(x)$.

Proof. Decomposing $w=w^{\prime}+w^{\perp}$ such that $w^{\prime}$ is parallel and $w^{\perp}$ orthogonal to $v$, the formula

$$
\left\langle\left(d \exp _{x}\right)_{v} v,\left(d \exp _{x}\right)_{v} w^{\prime}\right\rangle=\left\langle v, w^{\prime}\right\rangle
$$

is immediate from the Definition of $\exp _{x}$. By means of the linearity of $d \exp _{x}$, to verify (2.1.8) it is thus sufficient to consider the case $\langle v, w\rangle=0$.

We show the following: If $c:]-\varepsilon, \varepsilon\left[\rightarrow T_{x} M\right.$ is a curve $T_{x} M \cap \mathscr{O}(M)$ such that $|c(s)|=$ const, $c(0)=v$ and $\dot{c}(0)=w$, then for any $0<t_{0} \leq 1$ it holds that

$$
\left.\left.\frac{d}{d s}\right|_{s=0} \exp _{x}\left(t_{0} c(s)\right) \perp \frac{d}{d t}\right|_{t=t_{0}} \exp _{x}(t c(0))
$$

For $t_{0}=1$ this shows $\left(\exp _{x}\right)_{*} w \perp \dot{\gamma}_{v}(1)=\left(\exp _{x}\right)_{*} v$ as claimed.


Figure 2.1.1. Exponential function
Denting $\alpha_{s}(t)=\exp _{x}(t c(s))$ for $0 \leq t \leq 1$, then on one hand $L\left(\alpha_{s} \mid\left[0, t_{0}\right]\right)$ is independent of $s$ and by means of formula (2.1.2) (first variation of length) we have

$$
0=\left.\frac{d}{d s}\right|_{s=0} L\left(\alpha_{s} \mid\left[0, t_{0}\right]\right)=\left.\left\langle\left.\frac{d}{d s}\right|_{s=0} \alpha_{s}(t), \dot{\alpha}_{0}(t)\right\rangle\right|_{t=0} ^{t=t_{0}}
$$

$$
=\left\langle\left.\frac{d}{d s}\right|_{s=0} \exp _{x}\left(t_{0} c(s)\right), \dot{\gamma}_{v}\left(t_{0}\right)\right\rangle
$$

The second part of the claim is obvious: If $s \mapsto c(s)$ is a differentiable curve in $S^{n-1} \subset$ $T_{x} M$ and $\beta(s):=\exp _{x}(r c(s))$ die corresponding curve in $S_{r}(x)$, then as above $\dot{\beta}(s) \perp$ $\dot{\gamma}_{c(s)}(r)$.

THEOREM 2.1.6. Let $(M, g)$ be a Riemannian manifold and $x \in M$. Furthermore let $V_{\varepsilon}(0) \subset T_{x} M$ be an open $\varepsilon$-ball such that $\exp _{x} \mid V_{\varepsilon}(0)$ is an embedding. Then for any $v \in V_{\varepsilon}(0)$ the geodesic curve

$$
\begin{equation*}
\gamma_{v}:[0,1] \rightarrow M, \quad t \mapsto \exp _{x}(t v) \tag{2.1.9}
\end{equation*}
$$

has length $L\left(\gamma_{v}\right)=|v|=d\left(x, \exp _{x} v\right)$, and is modulo parametrization the only curve of length $d(x, y)$ connecting $x$ and $y:=\exp _{x} v$. In addition,

$$
\exp _{x}\left(V_{\varepsilon}(0)\right)=\{p \in M: d(x, p)<\varepsilon\}
$$

Notation 2.1.7. We call $B_{r}(x)=\{p \in M: d(x, p)<r\}$ geodesic ball about $x$ of radius $r$. For $\varepsilon<\varrho(x)$ we then have $B_{\varepsilon}(x)=\exp _{x}\left(V_{\varepsilon}(0)\right)$ and the hypersurface

$$
S_{\varepsilon}(x)=\{p \in M: d(x, p)=\varepsilon\}
$$

is called geodesic sphere about $x$ of radius $\varepsilon$. The geodesics in $B_{\varepsilon}(x)$ emanating from the center $x$ are called radial geodesics; by the Gauss Lemma they pass orthogonally through geodesic spheres about $x$.

Proof. (of Theorem 2.1.6): We use "geodesic polar coordinates" centered at $x$ on $B_{\varepsilon}(x) \backslash\{x\}=\exp _{x}\left(V_{\varepsilon}(0) \backslash\{0\}\right)$ and identify

$$
\phi:] 0, \varepsilon\left[\times S^{n-1} \xrightarrow{\sim} B_{\varepsilon}(x) \backslash\{x\}, \quad(r, \vartheta) \mapsto \exp _{x}(r \vartheta) .\right.
$$

By the Gauss Lemma we then have

$$
\phi^{*}\left(g \mid\left(B_{\varepsilon}(x) \backslash\{x\}\right)\right)=d r \otimes d r+h_{r}
$$

where $h_{r}$ denotes the Riemannian metric on $S^{n-1}$ defined by pullback under $\phi$ from the Riemannian metric on the geodesic $r$-sphere $S_{r}(x)$ induced by $g$.
(1) We show first that every piecewise differentiable curve $c:[0,1] \rightarrow M$ starting at $x$ which exits $B_{\varepsilon}(x)=\exp _{x}\left(V_{\varepsilon}(0)\right)$, has length $\geq \varepsilon$ hat.

To this end denote by $\left.\left.t_{1} \in\right] 0,1\right]$ the first time such that $c\left(t_{1}\right) \in \partial B_{\varepsilon}(x)=S_{\varepsilon}(x)$. Then $\left.c \mid] 0, t_{1}\right]$ has a unique representation of the form

$$
c(t)=\exp _{x}(r(t) \vartheta(t)) \equiv \phi(r(t), \vartheta(t)), \quad 0<t \leq t_{1}
$$

with piecewise differentiable curves $t \mapsto \vartheta(t)$ in $S^{n-1} \subset T_{x} M$ and $t \mapsto r(t)$ in $] 0, \infty[$ (without restriction we may assume that $c\left(t_{0}\right) \neq x$ for $\left.\left.t_{0} \in\right] 0, t_{1}\right]$; otherwise we neglect the interval $\left[0, t_{0}[\right.$ which only decreases the length of $c$. Then we have (up to a finite number of points)

$$
|\dot{c}(t)|^{2}=|\dot{r}(t)|^{2}+h_{r(t)}(\dot{\vartheta}(t), \dot{\vartheta}(t)),
$$

and we may estimate

$$
L(c) \geq \int_{0}^{t_{1}}|\dot{c}(t)| d t=\int_{0}^{t_{1}}\left[|\dot{r}(t)|^{2}+h_{r(t)}(\dot{\vartheta}(t), \dot{\vartheta}(t))\right]^{1 / 2} d t \geq \int_{0}^{t_{1}}|\dot{r}(t)| d t \geq \varepsilon
$$

(2) Let now $\gamma=\gamma_{v}$ as in 2.1.9 the geodesic from $x$ to $y:=\exp _{x} v$ and let $c:[0,1] \rightarrow$ $M$ be any piecewise differentiable curve connecting $x$ and $y$. Then $L(c) \geq L(\gamma)$ with equality if and only if $c=\gamma$ modulo parametrization. Indeed by (1) we assume that $c$ stays
entirely in $B_{\varepsilon}(x)$ and hence takes the form $c(t)=\exp _{x}(r(t) \cdot \vartheta(t))$ as in (1). This implies again

$$
L(c)=\int_{0}^{1}\left[|\dot{r}(t)|^{2}+h_{r(t)}(\dot{\vartheta}(t), \dot{\vartheta}(t))\right]^{1 / 2} d t \geq \int_{0}^{1}|\dot{r}(t)| d t \geq r(1)-r(0)=L(\gamma)
$$

with equality if and only if $\dot{\vartheta}(t) \equiv 0$ and $t \mapsto r(t)$ isotone.
Geodesic curves, which realize the distance between two points are called minimal geodesics. Theorem 2.1.6 shows in particular that any curve which realizes the distance between its end-points, is (after reparametrization) a minimal geodesic.

COROLLARY 2.1.8 (Geodesics as locally shortest curves). Let $(M, g)$ be a Riemannian manifold, $I \subset \mathbb{R}$ an open interval, and let $c: I \rightarrow M$ be a differentiable curve in $M$ parametrized proportional to arc length. The curve $c$ is a geodesic if and only if for each $t \in I$ there exists $\varepsilon>0$ such that $d(c(t), c(t+\varepsilon))=L(c \mid[t, t+\varepsilon])$.

Altogether, we can already give the following partial answer to the mentioned question concerning differentiability of the distance function $d(x, \cdot)$ : If $V_{\varepsilon}(0) \subset T_{x} M$ is an open $\varepsilon$-ball with $\varepsilon \leq \varrho(x)$, then

$$
d(x, \cdot)\left|B_{\varepsilon}(x)=| | \circ\left(\exp _{x} \mid V_{\varepsilon}(0)\right)^{-1}\right.
$$

is differentiable on the punctured geodesic ball $B_{\varepsilon}(x) \backslash\{x\}$ about $x$ of radius $\varepsilon$. The question, how large $\varepsilon$ can be chosen, requires hence information about the injectivity radius $\varrho(x)$ at $x$.

Before turning to such questions we note some facts about the metric structure of Riemannian manifolds. It is easy to verify that the distance function $d: M \times M \rightarrow \mathbb{R}_{+}$ defines indeed a metric on $M$ and that the topology of $M$ coincides with the metric topology of $(M, d)$.

THEOREM 2.1.9 (Hopf-Rinow). For a connected Riemannian manifold $(M, g)$ the following conditions are equivalent:
(i) $(M, d)$ is a complete metric space (i.e., every Cauchy sequence in $M$ is convergent).
(ii) $(M, g)$ is metrically complete (i.e., the domain of any geodesic can be extended to all of $\mathbb{R}$ ).
(iii) The exponential function $\exp _{x}$ is defined on all of $T_{x} M$ for at least one $x \in M$.

All three conditions imply that any two points in $M$ can be connected by a minimal geodesic.

Proof. (i) $\Rightarrow$ (ii): Otherwise there is a maximal geodesic $\gamma:] a, b[\rightarrow M$ such that $b<\infty$; without restrictions let $|\dot{\gamma}| \equiv 1$. We choose a monotonic sequence $\left(t_{n}\right)_{n \in \mathbb{N}}$ in $] a, b[$ such that $t_{n} \rightarrow b$. Since $d\left(\gamma\left(t_{m}\right), \gamma\left(t_{n}\right)\right) \leq\left|t_{m}-t_{n}\right|$ the sequence $\left(\gamma\left(t_{n}\right)\right)_{n \in \mathbb{N}}$ is then Cauchy, so that $x:=\lim _{n \rightarrow \infty} \gamma\left(t_{n}\right)$ exists by assumption.

About the point $x$ we choose a geodesic ball $B_{\varepsilon}(x)$ of radius $\varepsilon>0$ such that the injectivity radius $\varrho$ satisfies $\varrho \mid B_{\varepsilon}(x)>2 \varepsilon$. For large $n$ we then have $\gamma\left(t_{n}\right) \in B_{\varepsilon}(x)$, and $\gamma \mid\left[t_{n}, t_{n+k}\right]$ is the minimal geodesic connecting $\gamma\left(t_{n}\right)$ and $\gamma\left(t_{n+k}\right)$, in particular then

$$
d\left(\gamma\left(t_{n+k}\right), \gamma\left(t_{n}\right)\right)=t_{n+k}-t_{n}
$$

and with $k \rightarrow \infty$ then $d\left(x, \gamma\left(t_{n}\right)\right)=b-t_{n}$. Hence $\gamma \mid\left[t_{n}, t_{n+k}\right]$ is a curve with length $t_{n+k}-t_{n}$, which starts on the geodesic sphere $S_{b-t_{n}}(x)$ of radius $b-t_{n}$ and ends on the
geodesic sphere $S_{b-t_{n+k}}(x)$ of Radius $b-t_{n+k}$. By Theorem 2.1.6, $\gamma \mid\left[t_{n}, t_{n+k}\right]$ lies on a radial geodesic starting at $x$, and it follows that

$$
\gamma(t)=\exp _{x}((b-t) v), \quad b-\varepsilon<t<b
$$

for some $v \in T_{x} M$ with $|v|=1$. This shows that $\gamma$ can be extended beyond $b$ via

$$
\gamma(t)=\exp _{x}((b-t) v), \quad b \leq t<b+\varepsilon
$$

in contradiction to the maximality of $\gamma$.
(ii) $\Rightarrow$ claim of the addition: Let $x, y \in M$ such that $d(x, y)=r>0$; we want to show that $x$ and $y$ can be joined by a geodesic of length $r$. We fix a geodesic ball $B_{\varepsilon}(x)$ about $x$ of radius $\varepsilon<\varrho(x)$ and $y \notin B_{\varepsilon}(x)$. As a consequence of the compactness of $S_{\varepsilon}(x)$, there exists $x_{0} \in S_{\varepsilon}(x)$ of minimal distance to $y$; we get $x_{0}=\exp _{x}(\varepsilon v)$ for some $v \in T_{x} M$ with $|v|=1$.

Consider the geodesic curve $\gamma(t)=\exp _{x}(t v)$, by assumption defined for $t \in \mathbb{R}_{+}$. We want to show that $y=\exp _{x}(r v)$. Since $r=d(x, y) \leq d(x, \gamma(t))+d(\gamma(t), y) \leq$ $t+d(\gamma(t), y)$ we have $d(\gamma(t), y) \geq r-t$. We show that for any $t$ with $\varepsilon \leq t \leq r$ even

$$
\begin{equation*}
d(\gamma(t), y)=r-t \tag{2.1.10}
\end{equation*}
$$

holds which then gives the claim for $t=r$. First we verify (2.1.10) for $t=\varepsilon$ : indeed $d\left(x_{0}, y\right)=r-\varepsilon$ since

$$
r=d(x, y)=\min _{x^{\prime} \in S_{\varepsilon}(x)}\left(d\left(x, x^{\prime}\right)+d\left(x^{\prime}, y\right)\right)=\varepsilon+d\left(x_{0}, y\right)
$$

If (2.1.10) holds for $t \in[\varepsilon, r]$, then also for all $t^{\prime}$ with $\varepsilon \leq t^{\prime} \leq t$, since

$$
d\left(\gamma\left(t^{\prime}\right), y\right) \leq d\left(\gamma\left(t^{\prime}\right), \gamma(t)\right)+d(\gamma(t), y) \leq\left(t-t^{\prime}\right)+(r-t)=r-t^{\prime}
$$

Let now $t_{0}:=\sup \{t \in[\varepsilon, r]: d(\gamma(t), y)=r-t\}$; then (2.1.10) holds in particular also for $t_{0}$ and it remains to show that $t_{0}=r$. Supposing $t_{0}<r$, we may choose a sufficiently small geodesic ball $B_{\varepsilon_{1}}\left(\gamma\left(t_{0}\right)\right)$ about $\gamma\left(t_{0}\right)$ and $x_{1} \in S_{\varepsilon_{1}}\left(\gamma\left(t_{0}\right)\right)$ with minimal distance to $y$. Since $d\left(x_{1}, y\right)=d\left(\gamma\left(t_{0}\right), y\right)-\varepsilon_{1}=\left(r-t_{0}\right)-\varepsilon_{1}$, we then have

$$
d\left(x, x_{1}\right) \geq d(x, y)-d\left(x_{1}, y\right)=r-\left(r-t_{0}-\varepsilon_{1}\right)=t_{0}+\varepsilon_{1}
$$

Since the curve $\gamma \mid\left[0, t_{0}\right]$ from $x$ to $\gamma\left(t_{0}\right)$, prolongated by the radial geodesic from $\gamma\left(t_{0}\right)$ to $x_{1}$, has length $t_{0}+\varepsilon_{1}$, the curve realizes the distance $d\left(x, x_{1}\right)$; by the Corollary above it must be geodesic and hence coincide with $\gamma \mid\left[0, t_{0}+\varepsilon_{1}\right]$. In particular, then $x_{1}=\gamma\left(t_{0}+\varepsilon_{1}\right)$ and

$$
d\left(\gamma\left(t_{0}+\varepsilon_{1}\right), y\right)=d\left(x_{1}, y\right)=r-\left(t_{0}+\varepsilon_{1}\right)
$$

in contradiction to the Definition of $t_{0}$.
(ii) $\Rightarrow$ (iii) is a weakening.
(iii) $\Rightarrow$ (i): Since Cauchy sequences are bounded, it is sufficient to show that each bounded subset $A$ of $M$ is contained in a compact subset $K \subset M$. Let $x \in M$ be such that $\exp _{x}: T_{x} M \rightarrow M$ is well-defined on all of $T_{x} M$. Assume that $d(a, x) \leq r$ for all $a \in A$. By Theorem 2.1.9 then $a=\exp _{x} v$ for some $v \in T_{x} M$ with $|v| \leq r$, so that $A \subset \exp _{x}\left(\bar{V}_{r}(0)\right)$, where $K:=\exp _{x}\left(\bar{V}_{r}(0)\right)$ is compact as image of a compact set under a continuous map.

In the sequel we assume that $(M, g)$ is metrically complete and without restriction connected. For any $x \in M$ then $\exp _{x}$ is defined on all of $T_{x} M$ and defines for $r \leq \varrho(x)$ a diffeomorphism of $V_{r}(0)$ to $B_{r}(x)$. A point $x \in M$ is called pole for $(M, g)$ if $\varrho(x)=\infty$.

DEFINITION 2.1.10 (Cut locus). Let $(M, g)$ be a metrically complete Riemannian manifold. For $x \in M$ and $v \in T_{x} M$ with $|v|=1$ denote by $\gamma_{v}$ the geodesic curve starting at $x$ such that $\dot{\gamma}_{v}(0)=v$, i.e. $\gamma_{v}(t)=\exp _{x}(t v)$, and let

$$
\begin{aligned}
s(v): & =\sup \left\{t \geq 0: d\left(x, \gamma_{v}(t)\right)=t\right\} \\
& \left.\left.\equiv \sup \left\{t \geq 0: \gamma_{v} \mid[0, t] \text { is minimal geodesic }\right\} \in\right] 0, \infty\right]
\end{aligned}
$$

Then $C_{x}:=\left\{s(v) \cdot v: v \in T_{x} M,|v|=1, s(v)<\infty\right\}$ is called cut locus of $\exp _{x}$ in $x$ and die set

$$
\operatorname{cut}(x):=\exp _{x}\left(C_{x}\right) \subset M
$$

cut locus of $M$ with respect to $x$.
In case $s(v)<\infty$, the curve $\gamma_{v} \mid[0, t]$ stops to be the shortest connection between $x=\gamma_{v}(0)$ and $\gamma_{v}(t)$ for $t>s(v)$. The point $\gamma_{v}(s(v))$ is then also called cut point of $x$ along $\gamma_{v}$. By the Theorem of Hopf-Rinow, the curve $\gamma_{v}$ is cut at each point $\gamma_{v}(s(v)+\varepsilon)$ (for $\varepsilon>0$ ) by a shorter geodesic curve emanating from $x$.

We will show that always $\varrho(x)=d(x, \operatorname{cut}(x))$; by Theorem 2.1.6 it obvious that $\exp _{x} \mid V_{r}(0)$ for $r>d(x, \operatorname{cut}(x))$ is no longer an embedding, hence either no longer injective or it has critical points. We deal first with critical points of the exponential function.

DEFINITION 2.1.11 (Conjugate locus). Let $(M, g)$ be a metrically complete Riemannian manifold and $x \in M$. Critical points $v \in T_{x} M$ of $\exp _{x}: T_{x} M \rightarrow M$ are called vectors conjugate to $x$. Then

$$
K_{x}:=\left\{v \in T_{x} M: v \text { is a vector conjugate to } x\right\}
$$

is called conjugate locus of $\exp _{x}$ in $T_{x} M$ and the set

$$
\operatorname{Conj}(x):=\exp _{x}\left(K_{x}\right) \equiv\left\{y \in M: y \text { is critical value of } \exp _{x}\right\}
$$

the conjugate locus of $x$ in $M$. If $y \in \operatorname{Conj}(x)$ such that $y=\exp _{x} v$ with $v \in K_{x}$, one says that " $y$ is conjugate to $x$ along the geodesic $\gamma_{v}(t)=\exp _{x}(t v)(0 \leq t \leq 1)$ ".

The next theorem gives a basic characterization of the cut locus.
THEOREM 2.1.12. Let $(M, g)$ be a metrically complete Riemannian manifold, $\gamma$ a normal geodesic curve on $M$ and $\gamma\left(t_{0}\right)$ a cutting point of $x=\gamma(0)$ along $\gamma$. Then either
(i) $\gamma\left(t_{0}\right)$ is conjugate to $\gamma(0)$ along $\gamma$, or
(ii) there is a geodesic curve $\sigma \neq \gamma$ from $x$ to $\gamma\left(t_{0}\right)$ such that $L(\sigma)=L\left(\gamma \mid\left[0, t_{0}\right]\right)$.

Proof. Let $\gamma\left(t_{0}\right)$ be as described and $\left(\varepsilon_{n}\right)_{n \in \mathbb{N}}$ a sequence of real numbers such that $0<\varepsilon_{n} \rightarrow 0$. For any $n \in \mathbb{N}$ let $\sigma_{n}$ a normal minimal geodesic connecting $x$ and $\gamma\left(t_{0}+\right.$ $\left.\varepsilon_{n}\right)$. Then $\dot{\sigma}_{n}(0) \in T_{x} M$ and $\left|\dot{\sigma}_{n}(0)\right|=1$; by compactness of the unit sphere $S^{n-1} \subset$ $T_{x} M$ we may assume (after eventually passing to a subsequence) that $\dot{\sigma}_{n}(0)$ converges in $S^{n-1}$. Hence there is a geodesic curve $\sigma$ starting at $x$ such that $\dot{\sigma}_{n}(0) \rightarrow \dot{\sigma}(0)$. By the continuity of the exponential function, $\sigma$ is a minimal geodesic from $x$ to $\gamma\left(t_{0}\right)$ and hence $L\left(\sigma \mid\left[0, t_{0}\right]\right)=L\left(\gamma \mid\left[0, t_{0}\right]\right)$. If now $\sigma \neq \gamma$, then part (ii) of the claim is satisfied; it is hence sufficient to verify assertion (i) if $\sigma=\gamma$; thus we have to show that $d \exp _{x}$ is singular at $t_{0} \dot{\gamma}(0)$ if $\sigma=\gamma$ gilt.

Assume that $\dot{\sigma}(0)=\dot{\gamma}(0)$ and $d \exp _{x}$ not singular at $t_{0} \dot{\gamma}(0)$. Then there is an open neighbourhood $V$ of $t_{0} \dot{\gamma}(0)$ on which $\exp _{x}$ is an embedding. By Definition of $\sigma_{n}$ we have $\gamma\left(t_{0}+\varepsilon_{n}\right)=\sigma_{n}\left(t_{0}+\tilde{\varepsilon}_{n}\right)$ with $\tilde{\varepsilon}_{n} \leq \varepsilon_{n}$, since $\sigma_{n}$ is a minimal geodesic, and $\tilde{\varepsilon}_{n} \rightarrow 0$ for
$n \rightarrow \infty$. For $n$ sufficiently large then $\left(t_{0}+\tilde{\varepsilon}_{n}\right) \dot{\sigma}_{n}(0)$ und $\left(t_{0}+\varepsilon_{n}\right) \dot{\gamma}(0)$ lie in $V$, and we have

$$
\exp _{x}\left(\left(t_{0}+\varepsilon_{n}\right) \dot{\gamma}(0)\right)=\gamma\left(t_{0}+\varepsilon_{n}\right)=\sigma_{n}\left(t_{0}+\tilde{\varepsilon}_{n}\right)=\exp _{x}\left(\left(t_{0}+\tilde{\varepsilon}_{n}\right) \dot{\sigma}_{n}(0)\right)
$$

hence $\left(t_{0}+\varepsilon_{n}\right) \dot{\gamma}(0)=\left(t_{0}+\tilde{\varepsilon}_{n}\right) \dot{\sigma}_{n}(0)$ and then $\dot{\gamma}(0)=\dot{\sigma}_{n}(0)$ for large $n$, in contradiction to the Definition of $\sigma_{n}$.

If $(M, g)$ is a metrically complete Riemannian manifold, $x \in M$ a given point and $S^{n-1}=\left\{v \in T_{x} M:|v|=1\right\}$ the unit sphere in $T_{x} M$, one can show (e.g. [27], p. 98) that the map in Definition 2.1.10

$$
s: S^{n-1} \rightarrow \overline{\mathbb{R}}_{+}, \quad s(v)=\sup \left\{t \geq 0: \gamma_{v} \mid[0, t] \text { is a minimal geodesic }\right\}
$$

is continuous. Since in addition $s$ is strictly positive, the set

$$
U_{x}:=\left\{t v \in T_{x} M: v \in S^{n-1}, 0 \leq t<s(v)\right\}
$$

defines an open star-shaped neighbourhood of 0 in $T_{x} M$ with $\partial U_{x}=C_{x}$; according to Definition 2.1.10 then $\operatorname{cut}(x)=\exp _{x}\left(\partial U_{x}\right)$.

THEOREM 2.1.13. If $(M, g)$ is a metrically complete Riemannian manifold and $x \in$ M, then

$$
M=\exp _{x}\left(U_{x}\right) \dot{\cup} \operatorname{cut}(x)
$$

Proof. Let $y \in M$. By Theorem 2.1.9 (Hopf-Rinow) there is a minimal geodesic

$$
\gamma_{v}: \gamma_{v}(t)=\exp _{x}(t v), \quad|v|=1, \quad 0 \leq t \leq b
$$

connecting $x$ and $y$; hence $b \leq s(v)$. It rests to show the disjointness of the union. Suppose that $y \in \exp _{x}\left(U_{x}\right) \cap \operatorname{cut}(x)$, then $y=\exp _{x}\left(t_{0} v_{0}\right)=\exp _{x}\left(t_{1} v_{1}\right)$ with $v_{0}, v_{1} \in S^{n-1} \subset$ $T_{x} M$ such that $t_{0}<s\left(v_{0}\right)$ and $t_{1}=s\left(v_{1}\right)$. Both $\gamma_{v_{0}} \mid\left[0, t_{0}\right]$ and $\gamma_{v_{1}} \mid\left[0, t_{1}\right]$ are then minimal geodesics connecting $x$ and $y$, hence $t_{0}=t_{1}$. In addition, $\gamma_{v_{0}} \mid\left[0, t_{0}+\varepsilon\right]$ is still minimal for $\varepsilon>0$ sufficiently small. By the following Lemma, $\gamma_{v_{0}}$ can however not be minimal beyond the interval $\left[0, t_{0}\right]$.

Lemma 2.1.14. Let $\gamma: \mathbb{R} \rightarrow M$ be a geodesic on a Riemannian manifold $(M, g)$. If there is a geodesic curve $\sigma \neq \gamma$ connecting $\gamma(0)$ and $\gamma\left(t_{0}\right)$ with the same length as $\gamma$, then $\gamma \mid\left[0, t_{0}+\varepsilon\right]$ cannot be minimal for $\varepsilon>0$.

Proof. Assume that $\sigma:\left[0, t_{0}\right] \rightarrow M$ is a further geodesic connecting $\gamma(0)$ and $\gamma\left(t_{0}\right)$ of length $L(\sigma)=L\left(\gamma \mid\left[0, t_{0}\right]\right)$; furthermore suppose that $\gamma \mid\left[0, t_{0}+\varepsilon\right]$ is still minimal for some $\varepsilon>0$. Then also $c:\left[0, t_{0}+\varepsilon\right] \rightarrow M$ defined by $c\left|\left[0, t_{0}\right]=\sigma\right|\left[0, t_{0}\right]$ and $c \mid\left[t_{0}, t_{0}+\right.$ $\varepsilon]=\gamma \mid\left[t_{0}, t_{0}+\varepsilon\right]$, is a curve connecting $x$ and $\gamma\left(t_{0}+\varepsilon\right)$ with the length as $\gamma \mid\left[0, t_{0}+\varepsilon\right]$ as well. Thus also $c$ is a minimal geodesic curve which must coincide with $\gamma$ since $c \mid\left[t_{0}, t_{0}+\right.$ $\varepsilon]=\gamma \mid\left[t_{0}, t_{0}+\varepsilon\right]$. Consequently also $\gamma$ coincides with $\sigma$ on $\left[0, t_{0}\right]$.

Theorem 2.1.13 combined with Lemma 2.1.14 gives the following result.
Corollary 2.1.15. Let $(M, g)$ be a metrically complete Riemannian manifold and $x \in M$. Then, for each point $y \in M \backslash \operatorname{cut}(x)$ there is exactly one minimal geodesic connecting $x$ and $y$.

EXAMPLE 2.1.16. If $M=S^{n}$ denotes the $n$-dimensional sphere (considered as part of $\mathbb{R}^{n+1}$ with the induced canonical Riemannian metric), then for each point $x \in S^{n}$

$$
\operatorname{Conj}(x)=\operatorname{cut}(x)=\{-x\} .
$$

For $y \in \operatorname{Conj}(x)$ such that $y=\exp _{x} v$, one says that " $v$ lies in the first conjugate locus in $T_{x} M "$ if $\exp _{x}$ is regular at $t v$ for $0 \leq t<1$. We will see that each geodesic curve $\gamma$ in $M$ emanating from $x$ meets the cut locus cut $(x)$ along $\gamma$ not later than the first point conjugate to $x$. This then shows that $\exp _{x} \mid U_{x}$ is not only injective but a local diffeomorphism, and hence defines a diffeomorphism of $U_{x}$ to $\exp _{x}\left(U_{x}\right)$. Thus $M \backslash \operatorname{cut}(x)$ is diffeomorphic to an open ball in $\mathbb{R}^{n}$, and $\operatorname{cut}(x)$ itself is a strong deformation retract of $M \backslash\{x\}$. In this sense the cut locus cut $(x)$ contains the topology of $M$ and will hence in general have a complicated structure which indicates that Example 2.1.16 is not typical.

REMARK 2.1.17. $M \backslash \operatorname{cut}(x)$ can be characterized as the maximal open subset of $M$ with the property that each of its points can be uniquely joined with $x$ by a minimal geodesic curve.

Before turning to one of the fundamental questions of the theory, i.e. the problem when along a radial ray $t \mapsto t v$ in $T_{x} M$ the first conjugate vector in in $T_{x} M$ shows up, we insert a discussion of the general notion of curvature of a Riemannian manifold. The basic answer to the question above is then given by the so-called "comparison principle" which roughly speaking says that the first conjugate vector comes later the smaller curvature is.

REMARK 2.1.18. Let $(M, g)$ be a Riemannian manifold of of dimension at least 2 and $\nabla$ the Levi-Civita connection on $M$. By Definition 1.4.24, the Riemann curvature tensor $R \in \Gamma\left(T^{*} M^{\otimes 3} \otimes T M\right)$ is given by

$$
R(X, Y, Z) \equiv R(X, Y) Z:=\nabla_{X} \nabla_{Y} Z-\nabla_{Y} \nabla_{X} Z-\nabla_{[X, Y]} Z
$$

for $X, Y, Z \in \Gamma(T M)$. One may read $R$ equally either as $C^{\infty}(M)$-trilinear map

$$
\Gamma(T M)^{3} \rightarrow \Gamma(T M)
$$

or as

$$
\Gamma(T M \otimes T M) \longrightarrow \operatorname{Hom}_{C^{\infty}(M)}(\Gamma(T M), \Gamma(T M)) .
$$

Lemma 2.1.19 (curvature identities). For $X, Y, Z, U \in \Gamma(T M)$ one has:
(i) $\langle R(X, Y) Z, U\rangle=-\langle R(Y, X) Z, U\rangle=-\langle R(X, Y) U, Z\rangle$
(ii) $\langle R(X, Y) Z, U\rangle=\langle R(Z, U) X, Y\rangle$
(iii) $R(X, Y) Z+R(Y, Z) X+R(Z, X) Y=0$ (Bianchi identity).

Proof. (i): Anti-symmetry in the first and second argument is trivial, in the third and fourth argument it holds because of

$$
\begin{aligned}
\left\langle\nabla_{X} \nabla_{Y} Z, Z\right\rangle & =\frac{1}{2} X(Y\langle Z, Z\rangle)-\left\langle\nabla_{Y} Z, \nabla_{X} Z\right\rangle \\
\left\langle\nabla_{[X, Y]} Z, Z\right\rangle & =\frac{1}{2}[X, Y]\langle Z, Z\rangle
\end{aligned}
$$

(iii): The Bianchi identity follows from the fact that

$$
[X,[Y, Z]]+[Y,[Z, X]]+[Z,[X, Y]]=0
$$

for $X, Y, Z \in \Gamma(T M)$ (Jacobi identity for the Lie product of vector fields); on the other hand by torsion-freeness of $\nabla$ one has

$$
R(X, Y) Z+R(Y, Z) X+R(Z, X) Y=[X,[Y, Z]]+[Y,[Z, X]]+[Z,[X, Y]]=0
$$

(ii): According to the Bianchi identity one has

$$
\begin{aligned}
\langle R(X, Y) Z, U\rangle+\langle R(Y, Z) X, U\rangle+\langle R(Z, X) Y, U\rangle & =0 \\
-\langle R(X, Y) U, Z\rangle+\langle R(Y, U) X, Z\rangle+\langle R(U, X) Y, Z\rangle & =0
\end{aligned}
$$

$$
\begin{aligned}
-\langle R(Z, U) X, Y\rangle+\langle R(U, X) Z, Y\rangle+\langle R(X, Z) U, Y\rangle & =0 \\
\langle R(Z, U) Y, X\rangle+\langle R(U, Y) Z, X\rangle+\langle R(Y, Z) U, X\rangle & =0
\end{aligned}
$$

Addition of the four equations and taking into account i) leads to

$$
2\langle R(X, Y) Z, U\rangle-2\langle R(Z, U) X, Y\rangle=0
$$

which gives the claim.
Further curvature identities are obtained from the Riemannian curvature tensor by contraction.

DEFINITION 2.1.20 (Ricci curvature, scalar curvature). For a Riemannian manifold $(M, g)$ the tensor $\operatorname{Ric}^{M} \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$, defined by

$$
\operatorname{Ric}_{x}^{M}(u, v):=\operatorname{trace}\left(T_{x} M \rightarrow T_{x} M, w \mapsto R(w, u, v)\right) \equiv \sum_{i=1}^{n}\left\langle R\left(e_{i}, u\right) v, e_{i}\right\rangle
$$

where $\left(e_{1}, \ldots, e_{d}\right)$ denotes an orthonormal basis of $T_{x} M$, is called Ricci tensor of $(M, g)$; the symmetric bilinear form

$$
\operatorname{Ric}_{x}^{M}: T_{x} M \times T_{x} M \rightarrow \mathbb{R}
$$

is called Ricci curvature at $x$. The real-valued function $k^{M}$ finally,

$$
k^{M}(x):=\operatorname{trace} \operatorname{Ric}_{x}^{M}=\sum_{j=1}^{n} \operatorname{Ric}_{x}^{M}\left(e_{j}, e_{j}\right) \equiv \sum_{i=1}^{n} \sum_{j=1}^{n}\left\langle R\left(e_{i}, e_{j}\right) e_{j}, e_{i}\right\rangle
$$

where $\left(e_{1}, \ldots, e_{d}\right)$ denotes again an orthonormal basis of $T_{x} M$, is called scalar curvature of $(M, g)$.

DEFINITION 2.1.21 (Sectional curvature). Let $(M, g)$ be a Riemannian manifold with $\operatorname{dim} M \geq 2$. Furthermore let $G_{2} T M \rightarrow M$ be the Grassmann 2-bundle, defined as set by $G_{2} T M:=\bigcup_{x \in M} G_{2} T_{x} M$ where

$$
G_{2} T_{x} M:=\left\{E \subset T_{x} M: E \text { two-dimensional real subspace }\right\} .
$$

Then the map Riem ${ }^{M}: G_{2} T M \rightarrow \mathbb{R}$,

$$
\operatorname{Riem}^{M} \mid G_{2} T_{x} M \equiv \operatorname{Riem}_{x}^{M}: E=\operatorname{span}\{u, v\} \mapsto \frac{\langle R(u, v) v, u\rangle}{|u|^{2}|v|^{2}-\langle u, v\rangle^{2}}
$$

is well-defined (i.e. independent of the choice of $u$ and $v$ ) and is called Riemannian sectional curvature of $M$. Here $|u|^{2}|v|^{2}-\langle u, v\rangle^{2}=|u \wedge v|^{2}$ is the squared area of the parallelogram spanned by $u$ and $v$.

REMARK 2.1.22. The sectional curvature determines the Riemann curvature tensor.
Proof. Indeed, for $X, Y \in \Gamma(T M)$ at first $k(X, Y):=\langle R(X, Y) Y, X\rangle$ is uniquely determined by Riem ${ }^{M}$. By means of the curvature identities it holds however

$$
\begin{aligned}
6\langle R(X, Y) Z, U\rangle=k & (X+U, Y+Z)-k(X+U, Y)-k(X+U, Z) \\
& -k(X, Y+Z)-k(U, Y+Z)+k(X, Z)+k(U, Y) \\
& -k(Y+U, X+Z)+k(Y+U, X)+k(Y+U, Z) \\
& +k(Y, X+Z)+k(U, X+Z)-k(Y, Z)-k(U, X),
\end{aligned}
$$

so that $R$ is determined by $k$.

DEfinition 2.1.23. A Riemannian manifold $(M, g)$ is said to have constant (resp., positive, negative) curvature, if the sectional curvature Riem ${ }^{M}$ is constant (resp., positive, negative). The Riemannian manifold $(M, g)$ is said to be flat if Riem ${ }^{M} \equiv 0$ (equivalently, $R \equiv 0$ ). A Riemannian manifold $(M, g)$ is called Einstein manifold if $\mathrm{Ric}^{M}=c g$ for some real constant $c$.

We now turn again to the conjugacy behaviour of the exponential map. Let $\gamma$ be a geodesic curve on $M$ and $\gamma\left(t_{0}\right)$ a cut point of $x=\gamma(0)$ along $\gamma$. Then $\gamma \mid\left[0, t_{0}+\varepsilon\right]$ is for $\varepsilon>0$ no longer the shortest connection of $x=\gamma(0)$ and $\gamma\left(t_{0}+\varepsilon\right)$, which can mean either that a deformation of $\gamma$ provides shorter curves, or that there exist non-neighbouring curves of shorter length connecting $\gamma(0)$ and $\gamma\left(t_{0}+\varepsilon\right)$ for $\varepsilon>0$. We shall see that in the first case $\gamma\left(t_{0}\right)$ will be conjugates to $x$ along $\gamma$, i.e., $\gamma\left(t_{0}\right) \in \operatorname{Conj}(x)$. In general it will turn out that a geodesic curve $\gamma$ emanating from $x$, which does not hit the conjugate locus $\operatorname{Conj}(x)$ up to time $t_{0}$, is the shortest connection of $x$ and $\gamma\left(t_{0}\right)$, compared to all (piecewise differentiable) curves from $x$ to $\gamma\left(t_{0}\right)$ which are sufficiently close to $\gamma \mid\left[0, t_{0}\right]$.

In Corollary 2.1.8 we characterized geodesics as critical points of the length functional under smooth variation of curves. This point of view motivates to consider in addition to the first derivative (first variation) also the second derivative of the length functional.

THEOREM 2.1.24 (Second variation of arc length). Let $(M, g)$ be a Riemannian manifold with the Levi-Civita connection, $\gamma:[a, b] \rightarrow M$ a normal geodesic and let

$$
\alpha:[a, b] \times]-\varepsilon, \varepsilon[\rightarrow M
$$

be a differentiable variation of $\gamma$. In terms of $\gamma_{s}=\alpha(\cdot, s)$, along with $T=\frac{\partial}{\partial t} \alpha \equiv$ $\alpha_{*} D_{1} \in \Gamma\left(\alpha^{*} T M\right)$ and $V=\frac{\partial}{\partial s} \alpha \equiv \alpha_{*} D_{2} \in \Gamma\left(\alpha^{*} T M\right)$, for the second derivative of the length functional $L(s):=L\left(\gamma_{s}\right) \equiv \int_{a}^{b}\left|\dot{\gamma}_{s}(t)\right| d t$ at $s=0$ the so-called Synge formula holds:
$L^{\prime \prime}(0)=\left.\left\langle\nabla_{D_{2}} V, T\right\rangle(t, 0)\right|_{t=a} ^{t=b}+\int_{a}^{b}\left\{\left|\nabla_{D_{1}} V\right|^{2}-\langle R(V, T) T, V\rangle-\left(D_{1}\langle V, T\rangle\right)^{2}\right\}(t, 0) d t$.
Proof. Recall that by (2.1.3) we have $L^{\prime}(s)=\frac{d}{d s} L\left(\gamma_{s}\right)=\int_{a}^{b} \frac{1}{|T|}\left\langle\nabla_{D_{1}} V, T\right\rangle d t$. By means of Cartan's structural equations (Theorem 1.4.27) and the characterization of Riemannian connections in Theorem 1.5.6 (iii), this gives

$$
\begin{aligned}
L^{\prime \prime}(s) & =\int_{a}^{b}\left\{\frac{D_{2}\left\langle\nabla_{D_{1}} V, T\right\rangle}{|T|}-\frac{1}{2} \frac{\left\langle\nabla_{D_{1}} V, T\right\rangle D_{2}\langle T, T\rangle}{|T|^{3}}\right\} d t \\
& =\int_{a}^{b}\left\{\frac{\left\langle\nabla_{D_{2}} \nabla_{D_{1}} V, T\right\rangle+\left\langle\nabla_{D_{1}} V, \nabla_{D_{2}} T\right\rangle}{|T|}-\frac{\left\langle\nabla_{D_{1}} V, T\right\rangle\left\langle\nabla_{D_{2}} T, T\right\rangle}{|T|^{3}}\right\} d t \\
& =\int_{a}^{b}\left\{\frac{\langle R(V, T) V, T\rangle+\left\langle\nabla_{D_{1}} \nabla_{D_{2}} V, T\right\rangle+\left\langle\nabla_{D_{1}} V, \nabla_{D_{2}} T\right\rangle}{|T|}-\frac{\left\langle\nabla_{D_{1}} V, T\right\rangle^{2}}{|T|^{3}}\right\} d t .
\end{aligned}
$$

Now since $\nabla_{D_{1}} T=0$ and $|T|=1$ along $\gamma$, we obtain for $s=0$ the formula

$$
\begin{aligned}
L^{\prime \prime}(0) & =\int_{a}^{b}\left\{\left\langle\nabla_{D_{1}} V, \nabla_{D_{2}} T\right\rangle-\langle R(V, T) T, V\rangle+D_{1}\left\langle\nabla_{D_{2}} V, T\right\rangle-\left(D_{1}\langle V, T\rangle\right)^{2}\right\}(t, 0) d t \\
= & \left.\left\langle\nabla_{D_{2}} V, T\right\rangle(t, 0)\right|_{t=a} ^{t=b}+\int_{a}^{b}\left\{\left|\nabla_{D_{1}} V\right|^{2}-\langle R(V, T) T, V\rangle-\left(D_{1}\langle V, T\rangle\right)^{2}\right\}(t, 0) d t
\end{aligned}
$$

which proves the claim.

Notation 2.1.25. Let $\gamma:[a, b] \rightarrow M$ be a geodesic curve and $\alpha:[a, b] \times]-\varepsilon, \varepsilon[\rightarrow$ $M$ a differentiable variation of $\gamma$. The variation of $\gamma$ is called Jacobi variation if all neighbouring curves $\gamma_{s}=\alpha(\cdot, s)$ to $\gamma$ are geodesics. For $t \in[a, b]$ we say that $\alpha$ varies geodesically at $t$ if the induced curve $\alpha(t, \cdot):]-\varepsilon, \varepsilon[\rightarrow M$ is a geodesic.

If in the situation of Theorem 2.1.24 $\alpha:[a, b] \times]-\varepsilon, \varepsilon[\rightarrow M$ is a variation of $\gamma:[a, b] \rightarrow$ $M$ with fixed initial and end point (i.e., $\alpha(a, s) \equiv \alpha(a, 0)$ and $\alpha(b, s) \equiv \alpha(b, 0)$ for $-\varepsilon<s<\varepsilon$ ), or more generally, if $\alpha$ varies geodesically at the end points $t=a$ and $t=b$, then the term $\left.\left\langle\nabla_{D_{2}} V, T\right\rangle(t, 0)\right|_{t=a} ^{t=b}$ vanishes in the Synge formula for the second variation of the length, and Theorem 2.1.24 gives the following Corollary.

Corollary 2.1.26. Let $(M, g)$ be a Riemannian manifold, $\gamma:[a, b] \rightarrow M$ a normal geodesic curve and $\alpha:[a, b] \times]-\varepsilon, \varepsilon[\rightarrow M$ a variation of $\gamma$ which varies geodesically at $a$ and $b$. Denoting by $Y=\alpha_{*} D_{2}(\cdot, 0) \in \Gamma\left(\gamma^{*} T M\right)$ the corresponding variational field along $\gamma$ and $Y^{\perp}:=Y-\langle Y, \dot{\gamma}\rangle \dot{\gamma} \in \Gamma\left(\gamma^{*} T M\right)$ its orthogonal part, then for the second variation of the length the following formula holds:

$$
\begin{equation*}
L^{\prime \prime}(0)=\int_{a}^{b}\left\{\left|\nabla_{D} Y^{\perp}\right|^{2}-\left\langle R\left(Y^{\perp}, \dot{\gamma}\right) \dot{\gamma}, Y^{\perp}\right\rangle\right\} d t \tag{2.1.11}
\end{equation*}
$$

Proof. Indeed we have $\nabla_{D} Y^{\perp}=\nabla_{D} Y-\left\langle\nabla_{D} Y, \dot{\gamma}\right\rangle \dot{\gamma}=\left(\nabla_{D} Y\right)^{\perp}$ and hence

$$
\left|\nabla_{D} Y^{\perp}\right|^{2}=\left|\nabla_{D} Y\right|^{2}-\left\langle\nabla_{D} Y, \dot{\gamma}\right\rangle^{2}
$$

from where the claim follows since $\langle R(Y, \dot{\gamma}) \dot{\gamma}, Y\rangle=\left\langle R\left(Y^{\perp}, \dot{\gamma}\right) \dot{\gamma}, Y^{\perp}\right\rangle$.
REMARK 2.1.27. Let $\alpha:[a, b] \times]-\varepsilon, \varepsilon[\rightarrow M$ be a Jacobi variation of the geodesic $\gamma:[a, b] \rightarrow M$ and let $Y=\left(\alpha_{*} D_{2}\right)(\cdot, 0) \in \Gamma\left(\gamma^{*} T M\right)$. Then $D\langle Y, \dot{\gamma}\rangle$ is constant along $\gamma=\alpha(\cdot, 0)$. Hence, if $\langle Y, \dot{\gamma}\rangle$ vanishes at the end points $a$ and $b$, then $\langle Y, \dot{\gamma}\rangle$ vanishes already identically on $[a, b]$.

Proof. Let again $T=\alpha_{*} D_{1} \in \Gamma\left(\alpha^{*} T M\right)$ and $V=\alpha_{*} D_{2} \in \Gamma\left(\alpha^{*} T M\right)$. From $\nabla_{D_{1}} V-\nabla_{D_{2}} T=\alpha_{*}\left[D_{1}, D_{2}\right]=0$ and $\nabla_{D_{1}} T=0$ it follows first that

$$
\begin{equation*}
\nabla_{D_{1}} \nabla_{D_{1}} V=\nabla_{D_{1}} \nabla_{D_{2}} T=\nabla_{D_{1}} \nabla_{D_{2}} T-\nabla_{D_{2}} \nabla_{D_{1}} T=R(T, V) T \tag{2.1.12}
\end{equation*}
$$

and then

$$
D_{1} D_{1}\langle V, T\rangle=D_{1}\left\langle\nabla_{D_{1}} V, T\right\rangle=\left\langle\nabla_{D_{1}} \nabla_{D_{1}} V, T\right\rangle=\langle R(T, V) T, T\rangle=0
$$

where the last equality comes from Lemma 2.1.19 (i).
Hence if $\gamma:[a, b] \rightarrow M$ is a normal geodesic curve and $\alpha$ a Jacobi variation of $\gamma$ varying geodesically at the end points such that $\langle Y, \dot{\gamma}\rangle(a)=\langle Y, \dot{\gamma}\rangle(b)=0$ holds for the variational vector field $Y=\left(\alpha_{*} D_{2}\right)(\cdot, 0) \in \Gamma\left(\gamma^{*} T M\right)$, than by Remark 2.1.27 the Synge formula simplifies to

$$
\begin{equation*}
\left.\frac{d^{2}}{d s^{2}}\right|_{s=0} L\left(\gamma_{s}\right)=\int_{a}^{b}\left\{\left|\nabla_{D} Y\right|^{2}-\langle R(Y, T) T, Y\rangle\right\} d t \tag{2.1.13}
\end{equation*}
$$

The idea is now to bilinearize the Synge formula (2.1.13) for the second variation of the arc length which leads to the notion of the index form of $\gamma$.

DEFINITION 2.1.28 (Index form). Let $(M, g)$ be a Riemannian manifold, $\gamma:[a, b] \rightarrow$ $M$ a normal geodesic and $\Gamma^{\perp}\left(\gamma^{*} T M\right)$ the real vector space of piecewise differentiable vector fields $X$ along $\gamma$ such that $\langle X, \dot{\gamma}\rangle \equiv 0$. Then

$$
\begin{equation*}
I(X, Y):=\int_{a}^{b}\left\{\left\langle\nabla_{D} X, \nabla_{D} Y\right\rangle-\langle R(X, \dot{\gamma}) \dot{\gamma}, Y\rangle\right\} d t \tag{2.1.14}
\end{equation*}
$$

defines a symmetric bilinear form on $\Gamma^{\perp}\left(\gamma^{*} T M\right)$, the so-called index form of $\gamma$. The nullspace of $I$ is the linear subspace of $X \in \Gamma^{\perp}\left(\gamma^{*} T M\right)$ with the property that $I(X, Y)=$ 0 for all $Y \in \Gamma^{\perp}\left(\gamma^{*} T M\right)$.

The index form $I$ of a normal geodesic curve $\gamma:[a, b] \rightarrow M$ thus assigns to each differentiable vector field $Y \in \Gamma^{\perp}\left(\gamma^{*} T M\right)$ the second variation $L^{\prime \prime}(0)$ of the length $L$ with respect to the following variation of $\gamma$ (induced by $Y$ ),

$$
\begin{equation*}
\alpha:[a, b] \times]-\varepsilon, \varepsilon\left[\rightarrow M, \quad \alpha(t, s)=\exp _{\gamma(t)}\left(s Y_{t}\right),\right. \tag{2.1.15}
\end{equation*}
$$

that is $L^{\prime \prime}(0)=I(Y, Y)$ and $\left(\alpha_{*} D_{2}\right)(\cdot, 0)=Y$. Note that (2.1.15) is well-defined for $\varepsilon>0$ sufficiently small by the compactness of the interval $[a, b]$.

REMARK 2.1.29. If $I(Y, Y)<0$ for a differentiable vector field $Y \in \Gamma^{\perp}\left(\gamma^{*} T M\right)$ with $Y_{a}=0$ and $Y_{b}=0$, then there are curves arbitrarily close to $\gamma$ connecting $\gamma(a)$ and $\gamma(b)$ with a shorter length than $\gamma$. If however $I$ is positively definite on the subspace of differentiable vector fields $Y \in \Gamma^{\perp}\left(\gamma^{*} T M\right)$ vanishing at the end points, then the length of $\gamma$ is minimal compared to all variational curves sufficiently close to $\gamma$ with the same end points.

DEFINITION 2.1.30 (Jacobi field). Let $\gamma:[a, b] \rightarrow M$ be a geodesic on a Riemannian manifold $(M, g)$. A vector field $J \in \Gamma\left(\gamma^{*} T M\right)$ along $\gamma$ is said to be a Jacobi field along $\gamma$ if it satisfies the "Jacobi equation"

$$
\begin{equation*}
\nabla_{D} \nabla_{D} J+R(J, \dot{\gamma}) \dot{\gamma}=0 \tag{2.1.16}
\end{equation*}
$$

A Jacobi field $J$ along $\gamma$ is called proper, if in addition $\langle J, \dot{\gamma}\rangle=0$ holds.
It is easy to see that the Jacobi equation (2.1.16) is equivalent to a second order system of linear differential equations. Fixing a parallel section $e$ along $\gamma$ in $\mathrm{O}(T M)$, then $\left(e_{1}(t), \ldots, e_{d}(t)\right)$ is an orthonormal basis for $T_{\gamma(t)} M$ and $J$ writes as $J=\sum_{i}\left\langle J, e_{i}\right\rangle e_{i}$.

For the scalar functions $\left\langle J, e_{i}\right\rangle$ we have then $\left\langle J, e_{i}\right\rangle^{\prime} \equiv D\left\langle J, e_{i}\right\rangle=\left\langle\nabla_{D} J, e_{i}\right\rangle$ and $\left\langle J, e_{i}\right\rangle^{\prime \prime} \equiv D D\left\langle J, e_{i}\right\rangle=\left\langle\nabla_{D} \nabla_{D} J, e_{i}\right\rangle$, and the Jacobi equation (2.1.16) is equivalent to the system of linear differential equations

$$
\begin{equation*}
\left\langle J, e_{j}\right\rangle^{\prime \prime}=\sum_{i=1}^{n}\left\langle R\left(\dot{\gamma}, e_{i}\right) \dot{\gamma}, e_{j}\right\rangle\left\langle J, e_{i}\right\rangle, \quad j=1, \ldots, d \tag{2.1.17}
\end{equation*}
$$

By the theory of ordinary linear differential equations the system (2.1.17) has a $2 n$-dimensional space of solutions, and to each initial value and first derivative, corresponding to the data of $\left.J\right|_{t=t_{0}}$ and $\left.J^{\prime}\right|_{t=t_{0}}:=\left(\nabla_{D} J\right)\left(t_{0}\right)$ for some $t_{0}$, there is exactly one solution.

Since $\nabla_{D} \dot{\gamma}=0$ we observe in addition $\langle J, \dot{\gamma}\rangle^{\prime \prime}=\left\langle\nabla_{D} \nabla_{D} J, \dot{\gamma}\right\rangle=\langle R(\dot{\gamma}, J) \dot{\gamma}, \dot{\gamma}\rangle=0$. Each Jacobi field $J$ along $\gamma$ has hence a unique representation as

$$
\begin{equation*}
J=J^{\perp}+\left(c_{1}+t c_{2}\right) \dot{\gamma} \tag{2.1.18}
\end{equation*}
$$

with $J^{\perp}$ a proper Jacobi field (i.e. $\left\langle J^{\perp}, \dot{\gamma}\right\rangle=0$ ) and real constants $c_{1}, c_{2}$.
REMARK 2.1.31. Let $(M, g)$ be a Riemannian manifold, $\gamma:[0, b] \rightarrow M$ a geodesic and $\alpha:[0, b] \times]-\varepsilon, \varepsilon[\rightarrow M$ a Jacobi variation of $\gamma$. Then the "variational field"

$$
\begin{equation*}
J:=\left.\left(\alpha_{*} D_{2}\right)(\cdot, 0) \equiv \frac{\partial}{\partial s}\right|_{s=0} \alpha(\cdot, s) \in \Gamma\left(\gamma^{*} T M\right) \tag{2.1.19}
\end{equation*}
$$

is a Jacobi field along $\gamma$, and all Jacobi fields along $\gamma$ are obtained in this way.

Proof. That (2.1.19) defines a Jacobi field along $\gamma$ is a consequence of formula (2.1.12) for Jacobi variations. Conversely, let $J$ be an arbitrary Jacobi field along $\gamma$. We want to show that $J$ results from a variation of geodesic curves (Jacobi variation). To this end, we fix a curve $c:]-\varepsilon, \varepsilon\left[\rightarrow M\right.$ with $c(0)=\gamma(0)$ and $\dot{c}(0)=J_{0}$. Along $c$ we choose a vector field $W$ such that $W(0)=\dot{\gamma}(0)$ and $\left(\nabla_{D} W\right)_{0}=\left(\nabla_{D} J\right)_{0}$ (for instance, $W_{s}:=/ /_{0, s} \dot{\gamma}(0)+s / /_{0, s}\left(\nabla_{D} J\right)_{0} \in \Gamma\left(c^{*} T M\right)$ with $/ /_{0, s}$ the parallel transport along $c$ from $T_{c(0)} M$ to $\left.T_{c(s)} M\right)$. Then $\alpha(t, s):=\exp _{c(s)}\left(t W_{s}\right)$ defines a Jacobi variation of $\gamma$ and hence $\bar{J}=\left(\alpha_{*} D_{2}\right)(\cdot, 0)$ a Jacobi field along $\gamma$. But we have $\bar{J}_{0}=J_{0}$ and $\left(\nabla_{D} \bar{J}\right)_{0}=\left(\nabla_{D} J\right)_{0}$ (this follows with $T=\alpha_{*} D_{1}$ and $V=\alpha_{*} D_{2}$ according to $\left.\left(\nabla_{D} \bar{J}\right)_{0}=\left(\nabla_{D_{1}} V\right)_{(0,0)}=\left(\nabla_{D_{2}} T\right)_{(0,0)}=\left(\nabla_{D} W\right)_{0}=\left(\nabla_{D} J\right)_{0}\right)$; hence necessarily $\bar{J}=J$ holds.

The proof of Remark 2.1.31 provides in particular a method to construct Jacobi fields. The special case described in the following example is of particular importance.

EXAMPLE 2.1.32. Let $(M, g)$ be a Riemannian manifold, $\gamma:[0, b] \rightarrow M$ a geodesic and $J \in \Gamma\left(\gamma^{*} T M\right)$ a Jacobi field along $\gamma$ with $J(0)=0$. Then $J$ is the variational field to the variation

$$
\alpha:[0, b] \times]-\varepsilon, \varepsilon\left[\rightarrow M, \quad \alpha(t, s)=\exp _{x}\left[t\left(\dot{\gamma}(0)+s J^{\prime}(0)\right)\right]\right.
$$

where $\gamma(0)=x$ and $J^{\prime}(0):=\left(\nabla_{D} J\right)_{0}$; in other words::

$$
J(t)=\left.\frac{\partial}{\partial s}\right|_{s=0} \alpha(t, s)=\left(d \exp _{x}\right)_{t \dot{\gamma}(0)}\left(t J^{\prime}(0)\right) \in T_{\gamma(t)} M .
$$

Consider now the index form $I$ defined in (2.1.14) on the vector space $\Gamma_{0}^{\perp}\left(\gamma^{*} T M\right)$ of piecewise differentiable vector fields $X$ along a normal geodesic $\gamma:[a, b] \rightarrow M$ with $\langle X, \dot{\gamma}\rangle \equiv 0$ satisfying in addition $X_{a}=0$ and $X_{b}=0$. For $X, Y \in \Gamma_{0}^{\perp}\left(\gamma^{*} T M\right)$ and $a=t_{0}<t_{1}<\ldots<t_{n}=b$ a subdivision of the interval $[a, b]$ such that $X$ and $Y$ are differentiable on the subintervals $\left[t_{i}, t_{i+1}\right]$, one has $\left\langle\nabla_{D} X, \nabla_{D} Y\right\rangle=D\left\langle\nabla_{D} X, Y\right\rangle-$ $\left\langle\nabla_{D} \nabla_{D} X, Y\right\rangle$ on $\left[t_{i}, t_{i+1}\right]$ and hence for the index form:

$$
\begin{aligned}
I(X, Y) & =\int_{a}^{b}\left\{\left\langle\nabla_{D} X, \nabla_{D} Y\right\rangle-\langle R(X, \dot{\gamma}) \dot{\gamma}, Y\rangle\right\} d t \\
& =\left.\sum_{i=1}^{n-1}\left\langle\nabla_{D} X, Y\right\rangle\right|_{t_{i}+} ^{t_{i}-}-\int_{a}^{b}\left\{\left\langle\nabla_{D} \nabla_{D} X, Y\right\rangle+\langle R(X, \dot{\gamma}) \dot{\gamma}, Y\rangle\right\} d t
\end{aligned}
$$

If $X \mid\left[t_{i}, t_{i+1}\right]$ is a Jacobi field, it follows that

$$
I(X, Y)=\sum_{i}\left\langle\Delta_{t_{i}}\left(\nabla_{D} X\right), Y_{t_{i}}\right\rangle:=\left.\sum_{i}\left\langle\nabla_{D} X, Y\right\rangle\right|_{t_{i}+} ^{t_{i}-}
$$

THEOREM 2.1.33 (Jacobi fields as nullspace of the index form). Let $(M, g)$ be a Riemannian manifold, $\gamma:[a, b] \rightarrow M$ a normal geodesic and I the index form on $\Gamma_{0}^{\perp}\left(\gamma^{*} T M\right)$. Then the nullspace of I contains exactly the Jacobi fields $J$ along $\gamma$ vanishing at the end points: $J(a)=0$ and $J(b)=0$.

Proof. It is sufficient to show: From $I(X, Y)=0$ for all $Y \in \Gamma_{0}^{\perp}\left(\gamma^{*} T M\right)$ follows that $X$ is a Jacobi field. Let $a=t_{0}<t_{1}<\ldots<t_{n}=b$ be a subdivision of $[a, b]$ such that $X$ is differentiable on $\left[t_{i}, t_{i+1}\right]$, and $\varphi:[a, b] \rightarrow \mathbb{R}$ a differentiable function vanishing exactly at the places $t_{i}$ for $i=0, \ldots, n$. With

$$
Y:=\varphi \cdot\left(\nabla_{D} \nabla_{D} X+R(X, \dot{\gamma}) \dot{\gamma}\right) \in \Gamma_{0}^{\perp}\left(\gamma^{*} T M\right)
$$

one obtains that each $X \mid\left[t_{i}, t_{i+1}\right]$ is a Jacobi field. Considering then $I\left(X, Y^{0}\right)$ for an arbitrary vector field $Y^{0} \in \Gamma_{0}^{\perp}\left(\gamma^{*} T M\right)$ with $Y^{0}\left(t_{i}\right)=\Delta_{t_{i}}\left(\nabla_{D} X\right)$, gives the claim.

The next Theorem finally connects Jacobi fields to the conjugacy behaviour of the exponential function.

THEOREM 2.1.34. Let $(M, g)$ be a Riemannian manifold and $x, y \in M$; furthermore let $v \in T_{x} M$ with $y=\exp _{x} v$ and $\gamma:[0,1] \rightarrow M, \gamma(t):=\exp _{x}(t v)$ the connecting geodesic segment. The following statements are equivalent:
(i) $y$ is conjugate to $x$ along $\gamma$ (i.e., $v$ is a critical point of $\exp _{x}$ ).
(ii) There exists a non-identically vanishing Jacobi field $J$ along $\gamma$ with $J(0)=0$ and $J(1)=0$.
In particular, $x \in \operatorname{Conj}(y)$ if and only if $y \in \operatorname{Conj}(x)$.
Proof. (i) $\Rightarrow$ (ii): Let $y=\exp _{x} v$ and $v$ a critical point of $\exp _{x}$; then there exists $w \in T_{v}\left(T_{x} M\right) \cong T_{x} M$ such that $\left(d \exp _{x}\right)_{v} w=0$. The variation

$$
\alpha:[a, b] \times]-\varepsilon, \varepsilon\left[\rightarrow M \quad \alpha(t, s):=\exp _{x}(t(v+s w)),\right.
$$

of $\gamma$ is a Jacobi variation, and hence $J:=\alpha_{*} D_{2}(\cdot, 0)$ defines a Jacobi field according to Remark 2.1.31 which satisfies $J(0)=0$ and

$$
J(1)=\left.\frac{\partial}{\partial s}\right|_{s=0} \exp _{x}(v+s w)=\left(d \exp _{x}\right)_{v} w=0
$$

(i) $\Rightarrow$ (i): Conversely, let now $0 \neq J \in \Gamma\left(\gamma^{*} T M\right)$ be a Jacobi field along $\gamma$ with $J(0)=0$ and $J(1)=0$. Then, by Example 2.1.32,

$$
J(t)=\left(d \exp _{x}\right)_{t \dot{\gamma}(0)}(t w), \quad w=J^{\prime}(0) \neq 0
$$

Since $J(1)=\left(d \exp _{x}\right)_{v} w=0$ then $v$ is a critical point of $\exp _{x}$, and thus $y=\exp _{x} v$ is conjugate to $x$.

COROLLARY 2.1.35. Let $(M, g)$ be a Riemannian manifold and $\gamma:[a, b] \rightarrow M a$ geodesic with the property that $\gamma(a)$ and $\gamma(b)$ are conjugate to each other along $\gamma$. Then each Jacobi field $J$ along $\gamma$ is uniquely determined by the values $J(a)$ and $J(b)$.

Proof. The difference of two Jacobi fields along $\gamma$ with identical boundary values defines a Jacobi field which vanishes at $a$ and $b$, and hence vanishes identically by Theorem 2.1.34.

The next Theorem, the so-called Index Lemma, will serve as a crucial tool. I shows that Jacobi fields minimize the index form in a certain sense.

In the proof we use the following elementary observation.
LEMMA 2.1.36. Let $I \subset \mathbb{R}$ be an open real interval containing 0 and $h: I \rightarrow \mathbb{R}$ a differentiable function. Then there exists a differentiable function $\phi: I \rightarrow \mathbb{R}$ such that $h(t)=h(0)+t \phi(t)$ for $t \in I$.

Proof. Indeed, the function $\phi(t)=\int_{0}^{1} h^{\prime}(s t) d s$ satisfies the claim.
We assume the following situation: $(M, g)$ is a Riemannian manifold, $\gamma:[0, b] \rightarrow M$ a normal geodesic curve and $\Gamma^{\perp}\left(\gamma^{*} T M\right)$ the real vector space of piecewise differentiable vector fields $X$ along $\gamma$ such that $\langle X, \dot{\gamma}\rangle \equiv 0$. By Definition 2.1.28, on $\Gamma^{\perp}\left(\gamma^{*} T M\right)$ the index form of $\gamma$ is given:

$$
I(X, Y)=\int_{0}^{b}\left\{\left\langle\nabla_{D} X, \nabla_{D} Y\right\rangle-\langle R(X, \dot{\gamma}) \dot{\gamma}, Y\rangle\right\} d t, \quad X, Y \in \Gamma^{\perp}\left(\gamma^{*} T M\right)
$$

THEOREM 2.1.37 (Index Lemma). Let $(M, g)$ be a Riemannian manifold and suppose that $\gamma:[0, b] \rightarrow M$ is a normal geodesic with no points conjugate to $\gamma(0)$ along $\gamma$. Let J be a Jacobi field along $\gamma$ with $\langle J, \dot{\gamma}\rangle=0$ and $X$ a vector field in $\Gamma^{\perp}\left(\gamma^{*} T M\right)$. Suppose that $J(0)=X(0)=0$ and $J(b)=X(b)$. Then

$$
I(J, J) \leq I(X, X)
$$

with equality if and only if $J=X$.
Proof. (1) The real vector space $\mathcal{J}$ of Jacobi fields $J$ along $\gamma$ with $J(0)=0$ and $\langle J, \dot{\gamma}\rangle=0$ is of dimension $n-1$ where $n=\operatorname{dim} M$. Let $\left(J_{1}, \ldots, J_{n-1}\right)$ be a basis of $\mathcal{J}$ so that $J=\sum_{i} \alpha_{i} J_{i}$ with real constants $\alpha_{1}, \ldots, \alpha_{n}$. Since there is no $t$ such that $\gamma(t)$ is conjugate to $\gamma(0)$ along $\gamma$, according to Theorem 2.1.34, $\left(J_{1}(t), \ldots, J_{n-1}(t)\right)$ forms a basis of the orthogonal complement $\{\dot{\gamma}(t)\}^{\perp}$ of $\dot{\gamma}(t)$ in $T_{\gamma(t)} M$ for each $\left.\left.t \in\right] 0, b\right]$. Consequently, for any $t \in] 0, b]$, the vector field $X$ has a representation as

$$
\begin{equation*}
X(t)=\sum_{i=1}^{n-1} f_{i}(t) J_{i}(t) \tag{2.1.20}
\end{equation*}
$$

with $f_{i}$ piecewise differentiable functions on $\left.] 0, b\right]$. We want to check first that each $f_{i}$ can be differentiably extended to $t=0$, and hence to a piecewise differentiable function on $[0, b]$. Lemma 2.1.36, applied to the components $\left\langle J_{i}, e_{k}\right\rangle$ with respect to a parallel orthonormal basis $e=\left(e_{1} \ldots, e_{n}\right) \in \Gamma\left(\gamma^{*} \mathbf{O}(T M)\right)$ along $\gamma$, gives $J_{i}(t)=t A_{i}(t)$ with vector fields $A_{i} \in \Gamma\left(\gamma^{*} T M\right)$. In particular, then $\left(\nabla_{D} J_{i}\right)(0)=A_{i}(0)$ which shows the linear independence of $\left(A_{1}(0), \ldots, A_{n-1}(0)\right)$. For any $t \in[0, b]$ hence $\left(A_{1}(t), \ldots, A_{n-1}(t)\right)$ is a basis for $\{\dot{\gamma}(t)\}^{\perp}$ in $T_{\gamma(t)} M$, and one has $X(t)=\sum_{i} g_{i}(t) A_{i}(t)$ for $t \in[0, b]$, where $g_{i}$ are piecewise differentiable functions on $[0, b]$ with $g_{i}(0)=0$. Applying Lemma 2.1.36 one more time gives $g_{i}(t)=t h_{i}(t)$ with $h_{i}$ piecewise differentiable functions on $[0, b]$. Since $f_{i}(t)=h_{i}(t)$ for $t \neq 0$, this shows the wanted continuability.
(2) Next we show that on the interior of each subinterval, on which the $f_{i}$ are differentiable, the following formula holds:

$$
\begin{align*}
\left\langle\nabla_{D} X, \nabla_{D} X\right\rangle & -\langle R(X, \dot{\gamma}) \dot{\gamma}, X\rangle \\
& =\left\langle\sum_{i} f_{i}^{\prime} J_{i}, \sum_{i} f_{i}^{\prime} J_{i}\right\rangle+D\left\langle\sum_{i} f_{i} J_{i}, \sum_{i} f_{i} \nabla_{D} J_{i}\right\rangle \tag{2.1.21}
\end{align*}
$$

To shorting the notation we write $\langle A, A\rangle+D\langle X, B\rangle$ for the right-hand side of (2.1.21) where $A:=\sum_{i} f_{i}^{\prime} J_{i}$ and $B:=\sum_{i} f_{i} \nabla_{D} J_{i}$. Firstly we have

$$
R(X, \dot{\gamma}) \dot{\gamma}=\sum_{i} f_{i} R\left(J_{i}, \dot{\gamma}\right) \dot{\gamma}=-\sum_{i} f_{i} \nabla_{D} \nabla_{D} J_{i}=-C
$$

with $C=\sum_{i} f_{i} \nabla_{D} \nabla_{D} J_{i}$, and hence for the left-hand side of (2.1.21):

$$
\begin{aligned}
\left\langle\nabla_{D} X, \nabla_{D} X\right\rangle-\langle R(X, \dot{\gamma}) \dot{\gamma}, X\rangle & =\langle A+B, A+B\rangle-\langle R(X, \dot{\gamma}) \dot{\gamma}, X\rangle \\
& =\langle A, A\rangle+\langle A, B\rangle+\langle B, A\rangle+\langle B, B\rangle+\langle C, X\rangle
\end{aligned}
$$

On the other hand, letting $Q:=\sum_{i} f_{i}^{\prime} \nabla_{D} J_{i}$, we get for the right-hand side of (2.1.21):

$$
\begin{aligned}
\langle A, A\rangle+D\langle X, B\rangle & =\langle A, A\rangle+\langle A+B, B\rangle+\langle X, Q+C\rangle \\
& =\langle A, A\rangle+\langle A, B\rangle+\langle B, B\rangle+\langle X, Q\rangle+\langle X, C\rangle
\end{aligned}
$$

To verify (2.1.21) it is hence sufficient to show $\langle B, A\rangle=\langle X, Q\rangle$, or equivalently:

$$
\begin{equation*}
\left\langle\sum_{i} f_{i} \nabla_{D} J_{i}, \sum_{i} f_{i}^{\prime} J_{i}\right\rangle=\left\langle\sum_{i} f_{i} J_{i}, \sum_{i} f_{i}^{\prime} \nabla_{D} J_{i}\right\rangle \tag{2.1.22}
\end{equation*}
$$

For the verification of (2.1.22) we consider for fixed indices $i, j$ the function

$$
h:[0, b] \rightarrow \mathbb{R}, \quad h:=\left\langle\nabla_{D} J_{i}, J_{j}\right\rangle-\left\langle J_{i}, \nabla_{D} J_{j}\right\rangle
$$

since $h(0)=0$ and

$$
\begin{aligned}
h^{\prime} & =\left\langle\nabla_{D} \nabla_{D} J_{i}, J_{j}\right\rangle+\left\langle\nabla_{D} J_{i}, \nabla_{D} J_{j}\right\rangle-\left\langle\nabla_{D} J_{i}, \nabla_{D} J_{j}\right\rangle-\left\langle J_{i}, \nabla_{D} \nabla_{D} J_{j}\right\rangle \\
& =\left\langle R\left(\dot{\gamma}, J_{i}\right) \dot{\gamma}, J_{j}\right\rangle-\left\langle J_{i}, R\left(\dot{\gamma}, J_{j}\right) \dot{\gamma}\right\rangle=0
\end{aligned}
$$

we have $h \equiv 0$ on $[0, b]$. This shows Eq. (2.1.22), Using

$$
\sum_{i, j} f_{i} f_{j}^{\prime}\left\langle\nabla_{D} J_{i}, J_{j}\right\rangle=\sum_{i, j} f_{i} f_{j}^{\prime}\left\langle J_{i}, \nabla_{D} J_{j}\right\rangle
$$

this shows Eq. (2.1.22), and completes the proof of formula (2.1.21).
(3) Integration of (2.1.21) gives

$$
I(X, X)=\left\langle\sum_{i} f_{i} J_{i}, \sum_{j} f_{j} \nabla_{D} J_{j}\right\rangle(b)+\int_{0}^{b}\left\langle\sum_{i} f_{i}^{\prime} J_{i}, \sum_{j} f_{j}^{\prime} J_{j}\right\rangle d t
$$

analogously one obtains for the Jacobi field $J$ the equation

$$
I(J, J)=\left\langle\sum_{i} \alpha_{i} J_{i}, \sum_{j} \alpha_{j} \nabla_{D} J_{j}\right\rangle(b)
$$

By assumption, we have $J(b)=X(b)$, and hence $\alpha_{i}=f_{i}(b)$, which implies

$$
\begin{equation*}
I(X, X)=I(J, J)+\int_{0}^{b}\left|\sum_{i} f_{i}^{\prime} J_{i}\right|^{2} d t \geq I(J, J) \tag{2.1.23}
\end{equation*}
$$

This completes the proof of the first part of the Index Lemma.
(4) If now $I(X, X)=I(J, J)$, then $\sum_{i} f_{i}^{\prime} J_{i}=0$ by (2.1.23).

By part (1) $\left(J_{1}(t), \ldots, J_{n-1}(t)\right)$ is linearly independent for each $\left.\left.t \in\right] 0, b\right]$ which gives first $f_{i}^{\prime}=0$ on $\left.] 0, b\right]$ and by continuity then also on $[0, b]$. This shows $f_{i}=\mathrm{const}$ for each $i$, since $f_{i}(b)=\alpha_{i}$ hence $f_{i} \equiv \alpha_{i}$ for each $i$, and hence $J=X$.

A first consequence from the Index Lemma is that geodesics $\gamma$ minimize the length up to the first conjugate point compared to sufficiently close neighbouring curves of $\gamma$ with the same end points. Indeed, considering the case $J(b)=0$ in the Index Lemma (without restrictions assume that $\gamma$ is normal), we read off the following: If $\gamma\left(t_{0}\right)$ is the first point conjugate to $\gamma(0)$ along $\gamma$, then $I(X, X)>0$ for any vector field $X \neq 0$ along $\gamma \mid[0, t]$ with $\langle X, \dot{\gamma}\rangle=0$, provided $t<t_{0}$ and $X$ vanishes at the end points, i.e., $X(0)=0$ and $X(t)=0$. This shows $L^{\prime \prime}(0)>0$ for all variations of $\gamma \mid[0, t]$ with fixed end points.

In addition the following conversion holds:
COROLLARY 2.1.38. Let $(M, g)$ be a Riemannian manifold, $\gamma:[0, \infty[\rightarrow M$ a geodesic curve such that $\gamma\left(t_{0}\right)$ is conjugate to $\gamma(0)$ along $\gamma$. Then $\gamma \mid[0, t]$ is not minimal for $t>t_{0}$.

Proof. Without restriction let $\gamma$ be normal; denote by $\gamma\left(t_{0}\right)$ the first point conjugate to $\gamma(0)$ on $\gamma$. By Theorem 2.1.34, there is a Jacobi field $J \neq 0$ along $\gamma \mid\left[0, t_{0}\right]$ with $J(0)=0$ and $J\left(t_{0}\right)=0$. We choose $\varepsilon>0$ sufficiently small so that no pair of conjugate points exists on $\gamma \mid\left[t_{0}-\varepsilon, t_{0}+\varepsilon\right]$, and extend $J$ to a piecewise differentiable vector field $X$ along $\gamma \mid\left[0, t_{0}+\varepsilon\right]$ via

$$
X \mid\left[0, t_{0}\right]=J \quad \text { and } \quad X \mid\left[t_{0}, t_{0}+\varepsilon\right]=0
$$

Besides $X$ we consider another piecewise differentiable vector field $Y$ along $\gamma \mid\left[0, t_{0}+\varepsilon\right]$ given by

$$
Y \mid\left[0, t_{0}-\varepsilon\right]=J \quad \text { and } \quad Y \mid\left[t_{0}-\varepsilon, t_{0}+\varepsilon\right]=\tilde{J}
$$

where $\tilde{J}$ is the unique Jacobi field along $\gamma \mid\left[t_{0}-\varepsilon, t_{0}+\varepsilon\right]$ such that $\tilde{J}\left(t_{0}-\varepsilon\right)=J\left(t_{0}-\varepsilon\right)$ and $\tilde{J}\left(t_{0}+\varepsilon\right)=0$. Note that $\langle X, \dot{\gamma}\rangle=\langle Y, \dot{\gamma}\rangle=0$. Since $X$ and $Y$ agree on $\left[0, t_{0}-\varepsilon\right]$, but $X \mid\left[t_{0}-\varepsilon, t_{0}+\varepsilon\right]$ is no Jacobi field, we obtain

$$
I(Y, Y)<I(X, X)=0
$$

by the Index Lemma 2.1.37. Since $Y$ induces a variation of $\gamma \mid\left[0, t_{0}+\varepsilon\right]$ with fixed end points according to (2.1.15) so that $L^{\prime \prime}(0)=I(Y, Y)$ for the corresponding second variation of the length, there exists a variation of $\gamma \mid\left[0, t_{0}+\varepsilon\right]$ which keeps the end point fixed and shortens the length of $\gamma \mid\left[0, t_{0}+\varepsilon\right]$.

Absolute values of Jacobi fields can be compared by means of curvature relations. This is the content of the Comparison Theorem of Rauch.

THEOREM 2.1.39 (Rauch Comparison Theorem). Let $(M, g)$ and $(\tilde{M}, \tilde{g})$ be Riemannian manifolds with $2 \leq \operatorname{dim} M \leq \operatorname{dim} \tilde{M}$ and $\gamma:[0, b] \rightarrow M$, respectively $\tilde{\gamma}:[0, b] \rightarrow \tilde{M}$ normal geodesic curves. Furthermore let $J$ and $\tilde{J}$ be Jacobi fields along $\gamma$, resp. $\tilde{\gamma}$ with $J(0), \tilde{J}(0)$ parallel to $\gamma(0)$, resp. $\tilde{\gamma}(0)$ such that:

$$
|J(0)|=|\tilde{J}(0)|, \quad\left\langle\nabla_{D} J(0), \dot{\gamma}(0)\right\rangle=\left\langle\nabla_{D} \tilde{J}(0), \dot{\tilde{\gamma}}(0)\right\rangle, \quad\left|\nabla_{D} J(0)\right|=\left|\nabla_{D} \tilde{J}(0)\right|
$$

Suppose that there are no points along $\tilde{\gamma}$ conjugate to $\tilde{\gamma}(0)$ and that the curvature of $M$ along $\gamma$ does not exceed the curvature of $\tilde{M}$ along $\tilde{\gamma}$, i.e., for any $t \in[0, b]$ and for all planes $E \subset T_{\gamma(t)} M$ with $\dot{\gamma}(t) \in E$, resp. $\tilde{E} \subset T_{\tilde{\gamma}(t)} \tilde{M}$ with $\dot{\tilde{\gamma}}(t) \in \tilde{E}$ the sectional curvatures of the planes $E, \tilde{E}$ satisfy the inequality $\operatorname{Riem}^{M}(E) \leq \operatorname{Riem}^{\tilde{M}}(\tilde{E})$. Then for all $t \in[0, b]$,

$$
|J(t)| \geq|\tilde{J}(t)|
$$

Proof. (1) It is sufficient to prove the statements for Jacobi fields $J, \tilde{J}$ such that $J(0)=0, \tilde{J}(0)=0$ and $\langle J, \dot{\gamma}\rangle=\langle\tilde{J}, \dot{\tilde{\gamma}}\rangle \equiv 0$, since by (2.1.18) one has

$$
J=J^{\perp}+\left(c_{1}+t c_{2}\right) \dot{\gamma} \quad \text { and } \quad \tilde{J}=\tilde{J}^{\perp}+\left(\tilde{c}_{1}+t \tilde{c}_{2}\right) \dot{\tilde{\gamma}} ;
$$

but by assumption $J_{0}^{\perp}=0, \tilde{J}_{0}^{\perp}=0$, as well as $c_{1}=|J(0)|=|\tilde{J}(0)|=\tilde{c}_{1}$ and

$$
c_{2}=D\langle J, \dot{\gamma}\rangle=\left\langle\nabla_{D} J, \dot{\gamma}\right\rangle=\left\langle\nabla_{D} J(0), \dot{\gamma}(0)\right\rangle=\left\langle\nabla_{D} \tilde{J}(0), \dot{\tilde{\gamma}}(0)\right\rangle=\left\langle\nabla_{D} \tilde{J}, \dot{\tilde{\gamma}}\right\rangle=\tilde{c}_{2}
$$

Hence if $\left|J^{\perp}(t)\right| \geq\left|\tilde{J}^{\perp}(t)\right|$ is shown, we have because of $\langle J, \dot{\gamma}\rangle(t)=\langle\tilde{J}, \dot{\tilde{\gamma}}\rangle(t)$ also $|J(t)| \geq|\tilde{J}(t)|$ for $t \in[0, b]$. On the other hand, since $\nabla_{D}\left(J^{\perp}\right)=\left(\nabla_{D} J\right)^{\perp}$, resp. $\nabla_{D}\left(\tilde{J}^{\perp}\right)=\left(\nabla_{D} \tilde{J}\right)^{\perp}$, it is easy to see that with $J$ and $\tilde{J}$ also $J^{\perp}$ and $\tilde{J}^{\perp}$ satisfy the assumptions of the theorem.

In addition, we may assume that $\left|\nabla_{D} J(0)\right|=\left|\nabla_{D} \tilde{J}(0)\right|>0$, since in the case $\left|\nabla_{D} J(0)\right|=\left|\nabla_{D} \tilde{J}(0)\right|=0$ we have $|J|=|\tilde{J}|=0$, and the claim trivially holds true
(2) Letting $h(t):=|J(t)|^{2}$ and $\tilde{h}(t):=|\tilde{J}(t)|^{2}$, then $h(t) / \tilde{h}(t)$ for $\left.\left.t \in\right] 0, b\right]$ is welldefined, since along $\tilde{\gamma}$ there are conjugate points to $\tilde{\gamma}(0)$. An application of l'Hospital's rule then gives

$$
\lim _{t \rightarrow 0} \frac{h(t)}{\tilde{h}(t)}=\lim _{t \rightarrow 0} \frac{h^{\prime \prime}(t)}{\tilde{h}^{\prime \prime}(t)}=\lim _{t \rightarrow 0} \frac{\left\langle\nabla_{D} \nabla_{D} J, J\right\rangle(t)+\left\langle\nabla_{D} J, \nabla_{D} J\right\rangle(t)}{\left\langle\nabla_{D} \nabla_{D} \tilde{J}, \tilde{J}\right\rangle(t)+\left\langle\nabla_{D} \tilde{J}, \nabla_{D} \tilde{J}\right\rangle(t)}=\frac{\left|\nabla_{D} J(0)\right|^{2}}{\left|\nabla_{D} \tilde{J}(0)\right|^{2}}=1
$$

and for the verification of $|\tilde{J}| \leq|J|$ it is sufficient to check $\frac{d}{d t}(h(t) / \tilde{h}(t)) \geq 0$ on $\left.] 0, b\right]$, or equivalently: $h^{\prime} \tilde{h} \geq h \tilde{h}^{\prime}$ on $\left.] 0, b\right]$.

To this end, we fix $\left.\left.t_{0} \in\right] 0, b\right]$ for the rest of the proof and show that

$$
\begin{equation*}
h^{\prime}\left(t_{0}\right) \tilde{h}\left(t_{0}\right) \geq h\left(t_{0}\right) \tilde{h}^{\prime}\left(t_{0}\right) . \tag{2.1.24}
\end{equation*}
$$

Without loss of generality, we may assume $h\left(t_{0}\right)>0$ and $\tilde{h}\left(t_{0}\right)>0$ : For instance, if $h\left(t_{0}\right)=0$, then $h^{\prime}\left(t_{0}\right)=2\left\langle\nabla_{D} J\left(t_{0}\right), J\left(t_{0}\right)\right\rangle=0$ and (2.1.24) holds trivially; analogously for $\tilde{h}\left(t_{0}\right)=0$.
(3) Considering the vector fields $X:=\frac{J}{\left|J\left(t_{0}\right)\right|}$ and $\tilde{X}:=\frac{\tilde{J}}{\left|\tilde{J}\left(t_{0}\right)\right|}$ along $\gamma$, resp. along $\tilde{\gamma}$, we have:

$$
\begin{aligned}
\frac{h^{\prime}\left(t_{0}\right)}{h\left(t_{0}\right)} & =\langle X, X\rangle^{\prime}\left(t_{0}\right)=\int_{0}^{t_{0}}\langle X, X\rangle^{\prime \prime} d t \\
& =2 \int_{0}^{t_{0}}\left\{\left\langle\nabla_{D} X, \nabla_{D} X\right\rangle-\langle R(X, \dot{\gamma}) \dot{\gamma}, X\rangle\right\} d t=2 I_{t_{0}}(X, X)
\end{aligned}
$$

where $I_{t_{0}}(X, X)=I\left(X\left|\left[0, t_{0}\right], X\right|\left[0, t_{0}\right]\right)$; analogously it holds $\tilde{h}^{\prime}\left(t_{0}\right) / \tilde{h}\left(t_{0}\right)=2 I_{t_{0}}(\tilde{X}, \tilde{X})$.
To verify (2.1.24) it is hence sufficient to show $I_{t_{0}}(\tilde{X}, \tilde{X}) \leq I_{t_{0}}(X, X)$.
(4) We choose parallel orthonormal bases $e=\left(e_{1}, \ldots, e_{n}\right)$ and $\tilde{e}=\left(\tilde{e}_{1}, \ldots, \tilde{e}_{n+k}\right)$ (where $n+k=\operatorname{dim} \tilde{M}$ ) along $\gamma$, resp. $\tilde{\gamma}$, such that

$$
e_{1}=\dot{\gamma}, \quad e_{2}\left(t_{0}\right)=X\left(t_{0}\right) \quad \text { and } \quad \tilde{e}_{1}=\dot{\tilde{\gamma}}, \quad \tilde{e}_{2}\left(t_{0}\right)=\tilde{X}\left(t_{0}\right) .
$$

To each vector field $A \in \Gamma\left(\gamma^{*} T M\right)$ we associate a vector field $\iota A \in \Gamma\left(\tilde{\gamma}^{*} T \tilde{M}\right)$ via

$$
A=\sum_{i=1}^{n} a^{i} e_{i} \mapsto \iota A=\sum_{i=1}^{n} a^{i} \tilde{e}_{i} .
$$

Denoting by $\iota_{0}: T_{\gamma(0)} M \rightarrow T_{\tilde{\gamma}(0)} \tilde{M}$ the isometric embedding defined by $e_{i}(0) \mapsto \tilde{e}_{i}(0)$, we have $(\iota A)(t)=\left(/ \widetilde{/}_{0, t} \circ \iota_{0} \circ / /_{t, 0}\right) A(t)=: \iota_{t} A(t)$ with $/ / t, 0$ and $\widetilde{/}_{t, 0}$ the corresponding parallel transports along $\gamma$, resp. along $\tilde{\gamma}$. In particular for $A, B \in \Gamma\left(\gamma^{*} T M\right)$ it holds

$$
\langle\iota A, \iota B\rangle=\langle A, B\rangle \quad \text { and } \quad \nabla_{D}(\iota A)=\iota \nabla_{D} A
$$

By the curvature assumption and the fact that both geodesics are normal, we hence conclude $I_{t_{0}}(\iota X, \iota X) \leq I_{t_{0}}(X, X)$.

On the other hand, $\tilde{X}, \iota X$ are both vector fields along $\tilde{\gamma}$, and $\tilde{X}$ a Jacobi field, hence the assumptions of the Index Lemma (Theorem 2.1.37) are satisfied. In this situation the Index Lemma then gives

$$
I_{t_{0}}(\tilde{X}, \tilde{X}) \leq I_{t_{0}}(\iota X, \iota X) \leq I_{t_{0}}(X, X)
$$

which completes the proof of the Theorem.
Corollary 2.1.40 (Comparison Principle). Let $(M, g),(\tilde{M}, \tilde{g})$ be Riemannian manifolds such that $2 \leq \operatorname{dim} M \leq \operatorname{dim} \tilde{M}$ and let $\gamma:[0, b] \rightarrow M$, resp. $\tilde{\gamma}:[0, b] \rightarrow \tilde{M}$ be normal geodesic curves. If

$$
\operatorname{Riem}^{M}(E) \leq \operatorname{Riem}^{\tilde{M}}(\tilde{E})
$$

for all planes $E \subset T_{\gamma(t)} M$ with $\dot{\gamma}(t) \in E$, resp. $\tilde{E} \subset T_{\tilde{\gamma}(t)} \tilde{M}$ with $\dot{\tilde{\gamma}}(t) \in \tilde{E}$ and all $t \in[0, b]$, then along $\gamma$ the first conjugate point to $\gamma(0)$ does not appear before the first conjugate point to $\tilde{\gamma}(0)$ along $\tilde{\gamma}$.

Proof. We assume that $\tilde{\gamma}$ has no conjugate points $\tilde{\gamma}(0)$ along $\tilde{\gamma}$ on $\left[0, t_{0}\right]$. Let $J$ be a Jacobi field along $\gamma$ with $J(0)=0$, but $J \neq 0$. Then $\nabla_{D} J(0) \neq 0$, and we choose a Jacobi field $\tilde{J}$ along $\tilde{\gamma}$ with $\tilde{J}(0)=0$ such that

$$
\left\langle\nabla_{D} J(0), \dot{\gamma}(0)\right\rangle=\left\langle\nabla_{D} \tilde{J}(0), \dot{\tilde{\gamma}}(0)\right\rangle, \quad\left|\nabla_{D} J(0)\right|=\left|\nabla_{D} \tilde{J}(0)\right|
$$

Then $|J(t)| \geq|\tilde{J}(t)|>0$ for $\left.t \in] 0, t_{0}\right]$ where the fist inequality comes from the Comparison Theorem of Rauch, the second inequality holds according to Theorem 2.1.34. Applying Theorem 2.1.34 one more time then shows that also $\gamma \mid\left[0, t_{0}\right]$ has no points conjugate to $\gamma(0)$ along $\gamma$.

For a given manifold in general there there are topological obstructions for the existence of a Riemannian metric satisfying certain curvature conditions. For instance, negatively curved metrically complete Riemannian manifolds, which in addition are simply connected, are necessarily topologically trivial, as is shown in the next Theorem. We always assume metrically complete Riemannian manifolds to be connected.

THEOREM 2.1.41 (Theorem of Hadamard-Cartan). Any simply connected, metrically complete Riemannian manifold $(M, g)$ of curvature Riem $^{M} \leq 0$ is diffeomorphic to $\mathbb{R}^{n}$. More precisely: If $(M, g)$ is a metrically complete Riemannian manifold with Riem $^{M} \leq 0$, then $\exp _{x}: T_{x} M \rightarrow M$ is a covering for each $x \in M$, and hence a diffeomorphism if $M$ is in addition simply connected.

A differentiable map $f: \tilde{M} \rightarrow M$ between manifolds is said to be a covering, if to each point $x \in M$ there exists an open neighbourhood $U$ such that $f^{-1} U=\bigcup_{i \in I} \tilde{U}_{i}$ for some disjoint family $\left(\tilde{U}_{i}\right)_{i \in I}$ of open sets $\tilde{U}_{i}$ in $\tilde{M}$ with the property that $f \mid \tilde{U}_{i}: \tilde{U}_{i} \xrightarrow{\sim} U$ is a diffeomorphism for each $i \in I$.

Proof of Theorem 2.1.41. (1) Let $(M, g)$ be metrically complete and $x \in M$. According to the Theorem of Hopf-Rinow, $\exp _{x}$ is defined on all of $T_{x} M$ and surjective. If in addition $\operatorname{Riem}^{M} \leq 0$, then $\operatorname{Conj}(x)=\varnothing$ by the Comparison Principle with $\left(\mathbb{R}^{n}\right.$, eucl) as comparison manifold. Hence $\exp _{x}: T_{x} M \rightarrow M$ is a local diffeomorphism und $\left(T_{x} M, \exp _{x}^{*} g\right)$ a metrically complete Riemannian manifold: Metric completeness follows from the Theorem of Hopf-Rinow; geodesic curves emanating from $0 \in T_{x} M$ correspond to the half-rays starting at the origin.
(2) It is hence sufficient to show: Each local isometry $f:(\tilde{M}, \tilde{g}) \rightarrow(M, g)$ between Riemannian manifolds of the same dimension is already a covering in case $(\tilde{M}, \tilde{g})$ is metrically complete. Let now $x \in M$; we show that there exists a connected open neighbourhood $U$ of $x$ in $M$ such that $f$ maps each connected component $\tilde{U}_{i}$ of $f^{-1} U$ diffeomorphically onto $U$. To this end, we choose $r>0$ sufficiently small such that $\exp _{x}$ maps the $r$-ball $V_{r}\left(0_{x}\right)$ about $0_{x} \in T_{x} M$ diffeomorphically onto the geodesic $r$-ball $B_{r}(x)=: U \subset M$ about $x$. If $f^{-1}\{x\}=\left\{\tilde{x}_{i}: i \in I\right\}$ let $\tilde{U}_{i}:=B_{r}\left(\tilde{x}_{i}\right) \subset \tilde{M}$; we claim

$$
\bigcup_{i \in I} \tilde{U}_{i}=f^{-1} U \quad \text { and } \quad f \mid \tilde{U}_{i}: \tilde{U}_{i} \xrightarrow{\sim} U
$$

Firstly, for fixed $i \in I$, we have $\exp _{x} \circ d f_{\tilde{x}_{i}}=f \circ \exp _{\tilde{x}_{i}}$ : indeed if $v \in T_{\tilde{x}_{i}} \tilde{M}$ and if $\gamma$ is the geodesic with $\dot{\gamma}(0)=v$, then $f \circ \gamma$ is a geodesic on $M$, since $f$ maps as local isometry geodesics to geodesics; thus $(f \circ \gamma)(t)=\exp _{x}(t w)$ with $w:=d f_{\tilde{x}_{i}} v \in T_{x} M$ and hence
$\left(f \circ \exp _{\tilde{x}_{i}}\right)(v)=(f \circ \gamma)(1)=\exp _{x}\left(d f \tilde{x}_{i} v\right)$. The diagram

consequently commutes and in particular also $(M, g)$ is metrically complete. Restriction of the maps in (2.1.25) gives

since $\exp _{x} \circ d f_{\tilde{x}_{i}}$ maps $V_{r}\left(0_{\tilde{x}_{i}}\right)$ diffeomorphically onto $B_{r}(x)=U$, hence $\exp _{\tilde{x}_{i}}$ maps $V_{r}\left(0_{\tilde{x}_{i}}\right)$ diffeomorphically to $B_{r}\left(\tilde{x}_{i}\right)$ and consequently $f \mid \tilde{U}_{i}: \tilde{U}_{i} \xrightarrow{\sim} U$ is a diffeomorphism.

Trivially $\bigcup_{i \in I} \tilde{U}_{i} \subset f^{-1} U$; we want to verify $f^{-1} U \subset \bigcup_{i \in I} \tilde{U}_{i}$. To this end, let $\tilde{y} \in f^{-1} U$ and $y:=f(\tilde{y})$. We consider the minimal normal geodesic $c:\left[0, t_{0}\right] \rightarrow M$ connecting $y$ and $x$; it holds $t_{0}=d(x, y)<r$. To $w=\dot{c}(0) \in T_{y} M$ there is a unique tangent vector $v \in T_{\tilde{y}} \tilde{M}$ with $d f_{\tilde{y}} v=w$. By the metric completeness, the geodesic $\tilde{c}(t):=\exp _{\tilde{y}}(t v)$ on $M$ is defined on all of $\mathbb{R}$; by construction $f \circ \tilde{c}=c$. Hence $(f \circ \tilde{c})\left(t_{0}\right)=$ $c\left(t_{0}\right)=x$, and thus $\tilde{c}\left(t_{0}\right)=\tilde{x}_{i}$ for some $i \in I$. But since $d\left(\tilde{x}_{i}, \tilde{y}\right) \leq t_{0}<r$, we have $\tilde{y} \in B_{r}\left(\tilde{x}_{i}\right)=\tilde{U}_{i}$.

It remains to show that $\tilde{U}_{i} \cap \tilde{U}_{j}=\varnothing$ for $i \neq j$. Let $\gamma:\left[0, t_{1}\right] \rightarrow \tilde{M}$ the minimal normal geodesic that connects $\tilde{x}_{i}$ and $\tilde{x}_{j}$. Then $f \circ \gamma$ is a closed geodesic curve on $M$ with $x$ as initial and end point. Hence $f \circ \gamma$ does not lie in the geodesic ball $U=B_{r}(x)$ and must hence have length $>2 r$. But this shows $t_{1}=d\left(\tilde{x}_{i}, \tilde{x}_{j}\right)>2 r$.

Let $(M, g)$ be a metrically complete Riemannian manifold, $x_{0} \in M$ and $\gamma:[0, b] \rightarrow$ $M$ a normal geodesic curve such that $\gamma(0)=x_{0}$. If for $\left.\left.t_{0} \in\right] 0, b\right]$ the point $x:=\gamma\left(t_{0}\right)$ is not conjugate to $x_{0}=\gamma(0)$ along $\gamma$, then $\exp _{x_{0}}: T_{x_{0}} M \rightarrow M$ maps an open neighbourhood of $t_{0} \dot{\gamma}(0)$ diffeomorphically to an open neighbourhood of $x$ in $M$. Hence, if there are no points conjugate to $\gamma(0)$ along $\gamma$, then $\exp _{x_{0}}$ maps an open neighbourhood $V$ of the ray $\{t \dot{\gamma}(0): 0 \leq t \leq b\}$ locally diffeomorphically to an open neighbourhood $U$ of $\gamma([0, b])$ in M.

For instance, if $\gamma:[0, b] \rightarrow M$ is a normal geodesic starting at $x_{0}$ which does not hit the cut locus $C\left(x_{0}\right)$, then by Corollary 2.1.38 there are no conjugate points to $\gamma(0)=x_{0}$ along $\gamma$; in addition we may choose $V$ and $U$ such that $\exp _{x_{0}}$ maps $V$ diffeomorphically to $U$. In particular, $r=|\cdot| \circ\left(\exp _{x_{0}} \mid V\right)^{-1}$ is well-defined on $U$, and coincides there with the distance function $d\left(x_{0}, \cdot\right)$, i.e.

$$
\begin{equation*}
r=d\left(x_{0}, \cdot\right)=|\cdot| \circ\left(\exp _{x_{0}} \mid V\right)^{-1} \tag{2.1.27}
\end{equation*}
$$

consequently $r$ is differentiable on $U \backslash\left\{x_{0}\right\}$.
Assuming that $\gamma:[0, b] \rightarrow M$ does not hit the cut locus of $\gamma(0)$, we fix $x=\gamma\left(t_{0}\right)$ with $\left.\left.t_{0} \in\right] 0, b\right]$ and consider for $u \in T_{x} M$ the geodesic curve $\left.c:\right]-\varepsilon, \varepsilon[\rightarrow U \subset M$ satisfying $c(0)=x$ and $\dot{c}(0)=u$. Through the induced curve $\beta:=\frac{1}{t_{0}}\left(\exp _{x_{0}} \mid V\right)^{-1} \circ c$ in $T_{x} M$ we
obtain a Jacobi variation of $\gamma \mid\left[0, t_{0}\right]$, namely

$$
\left.\alpha:\left[0, t_{0}\right] \times\right]-\varepsilon, \varepsilon\left[\rightarrow M, \quad \alpha(t, s):=\exp _{x_{0}}(t \beta(s))\right.
$$

It holds $\alpha(0, \cdot)=x_{0}$ and $\alpha\left(t_{0}, \cdot\right)=c$; hence $X:=\alpha_{*} D_{2}(\cdot, 0) \in \Gamma\left(\gamma^{*} T M\right)$ is the unique Jacobi field with $X(0)=0$ and $X\left(t_{0}\right)=u$. Furthermore, we have

$$
L(\alpha(\cdot, s))=t_{0}|\beta(s)|=|\cdot| \circ\left(\exp _{x_{0}} \mid V\right)^{-1}(c(s))=(r \circ c)(s)
$$

By the formulas (2.1.2) and (2.1.11) for the first and second variation of length, we then have:

$$
\begin{align*}
(d r)_{x}(u) & =(r \circ c)^{\prime}(0)=\left.\langle X, \dot{\gamma}\rangle\right|_{0} ^{t_{0}}  \tag{2.1.28}\\
(\nabla d r)_{x}(u, u) & =(r \circ c)^{\prime \prime}(0)=I_{t_{0}}\left(X^{\perp}, X^{\perp}\right)
\end{align*}
$$

where $X^{\perp}$ denotes the orthogonal component of the Jacobi fields $X$ along $\gamma \mid\left[0, t_{0}\right]$.
We want to use (2.1.28) to derive comparison theorems for the Hessian $\nabla d r$ depending on curvature relations.

To this end, let $(\tilde{M}, \tilde{g})$ be a further metrically complete Riemannian manifold with $\operatorname{dim} M \leq \operatorname{dim} \tilde{M}$ and $\tilde{\gamma}:[0, b] \rightarrow \tilde{M}$ an additional normal geodesic curve. To put $M$ and $\tilde{M}$ in relation, as in part (4) of the proof to Theorem 2.1.39, we choose an isometric embedding $\iota_{0}: T_{\gamma(0)} M \rightarrow T_{\tilde{\gamma}(0)} \tilde{M}$ with $\iota_{0}(\dot{\gamma}(0))=\dot{\tilde{\gamma}}(0)$ and extend it via parallel transport to isometric embeddings $\iota_{t}: T_{\gamma(t)} M \rightarrow T_{\tilde{\gamma}(t)} \tilde{M}$ :

$$
\begin{align*}
& T_{\gamma(0)} M \stackrel{\iota_{0}}{\longleftrightarrow} T_{\tilde{\gamma}(0)} \tilde{M} \\
& /\left._{0, t}\right|_{\downarrow}  \tag{2.1.29}\\
& T_{\gamma(t)} M \stackrel{\downarrow}{\hookrightarrow--\tilde{/}_{0, t}} \\
& \iota_{t} \\
& T_{\tilde{\gamma}(t)} \tilde{M}
\end{align*}
$$

In this way an isometric bundle embedding $\iota: \gamma^{*} T M \rightarrow \tilde{\gamma}^{*} T \tilde{M}$ over $\mathbb{R}$ with the properties is obtained:

$$
\iota \dot{\gamma}=\dot{\tilde{\gamma}} \quad \text { and } \quad \nabla_{D}(\iota A)=\iota \nabla_{D} A \quad \text { for } A \in \Gamma\left(\gamma^{*} T M\right)
$$

We assume that there is no cut point of $x_{0}=\gamma(0)$ along $\gamma$ and no cut point of $\tilde{x}_{0}=$ $\tilde{\gamma}(0)$ along $\tilde{\gamma}$, and fix for some $\left.\left.t_{0} \in\right] 0, b\right]$ the points $x=\gamma\left(t_{0}\right) \in M$, resp. $\tilde{x}=\tilde{\gamma}\left(t_{0}\right) \in \tilde{M}$. The functions

$$
\begin{equation*}
r=|\cdot| \circ\left(\exp _{x_{0}} \mid V\right)^{-1} \quad \text { and } \quad \tilde{r}=|\cdot| \circ\left(\exp _{\tilde{x}_{0}} \mid \tilde{V}\right)^{-1} \tag{2.1.30}
\end{equation*}
$$

are defined according to (2.1.27). By definition, then $r(x)=\tilde{r}(\tilde{x})=t_{0}$, and for the differentials of $r$ and $\tilde{r}$ at $x$, resp. $\tilde{x}$ the following result holds:

Lemma 2.1.42. Keeping the notions from above, it holds

$$
d(f \circ r)_{x}=\iota_{t_{0}}^{*} d(f \circ \tilde{r})_{\tilde{x}} \equiv d(f \circ \tilde{r})_{\tilde{x}} \circ \iota_{t_{0}}
$$

for each $C^{1}$-function $f:[0, \infty[\rightarrow \mathbb{R}$.
Proof. By the chain rule the claim is reduced to the case $f(t)=t$. Let now $u \in$ $T_{x} M$. Consider along $\gamma \mid\left[0, t_{0}\right]$ the Jacobi field $X$ with $X(0)=0$ and $X\left(t_{0}\right)=u$, and along $\tilde{\gamma} \mid\left[0, t_{0}\right]$ the Jacobi field $\tilde{X}$ with $\tilde{X}(0)=0$ and $\tilde{X}\left(t_{0}\right)=\iota_{t_{0}} u$. Since $\langle X, \dot{\gamma}\rangle=$ $\langle\iota X, \iota \dot{\gamma}\rangle=\langle\tilde{X}, \dot{\tilde{\gamma}}\rangle$, we obtain with the first part of (2.1.28) the claim:

$$
(d r)_{x}(u)=\left.\langle X, \dot{\gamma}\rangle\right|_{0} ^{t_{0}}=\left.\langle\tilde{X}, \dot{\tilde{\gamma}}\rangle\right|_{0} ^{t_{0}}=(d \tilde{r})_{\tilde{x}}\left(\iota_{t_{0}} u\right)
$$

The Hessians are however no longer equal in the sense above, but they can be estimated against each other by means of curvature relations. To this end, we use the following notation:

NOTATION 2.1.43. For symmetric bilinear forms

$$
b \in \Gamma\left(\gamma^{*}\left(T^{*} M \otimes T^{*} M\right)\right), \quad \tilde{b} \in \Gamma\left(\tilde{\gamma}^{*}\left(T^{*} \tilde{M} \otimes T^{*} \tilde{M}\right)\right)
$$

along $\gamma$, resp. along $\tilde{\gamma}$, we write

$$
b \succcurlyeq \tilde{b},
$$

if for each $t \in[0, b]$ and each isometric bundle embedding $\iota: \gamma^{*} T M \rightarrow \tilde{\gamma}^{*} T \tilde{M}$ over $\mathbb{R}$, which as in (2.1.29) is induced by parallel transport from an isometric embedding

$$
\iota_{0}: T_{\gamma(0)} M \rightarrow T_{\tilde{\gamma}(0)} \tilde{M}
$$

with $\iota_{0}(\dot{\gamma}(0))=\dot{\tilde{\gamma}}(0)$, the symmetric bilinear form $b_{\gamma(t)}-\iota_{t}^{*} \tilde{b}_{\tilde{\gamma}(t)}$ on $T_{\gamma(t)} M$ is positive semidefinite. In other words:

$$
\begin{array}{ll}
b \succcurlyeq \tilde{b} \Longleftrightarrow b_{\gamma(t)}(u, u) \geq \tilde{b}_{\tilde{\gamma}(t)}(\tilde{u}, \tilde{u}) \quad \text { for } t \in[0, b], u \in T_{\gamma(t)} M, \tilde{u} \in T_{\tilde{\gamma}(t)} \tilde{M}: \\
& |u|=|\tilde{u}|, \quad\langle u, \dot{\gamma}(t)\rangle=\langle\tilde{u}, \dot{\tilde{\gamma}}(t)\rangle .
\end{array}
$$

THEOREM 2.1.44 (Hessian Comparison Theorem). Let $(M, g)$ and ( $\tilde{M}, \tilde{g})$ be Riemannian manifolds with $2 \leq \operatorname{dim} M \leq \operatorname{dim} \tilde{M}$, and let $\gamma:[0, b] \rightarrow M$, resp., $\tilde{\gamma}:[0, b] \rightarrow$ $\tilde{M}$ be minimal normal geodesic curves. If then the curvature of $M$ along $\gamma$ does not exceed the curvature of $\tilde{M}$ along $\tilde{\gamma}$, in the sense that always

$$
\operatorname{Riem}^{M}(E) \leq \operatorname{Riem}^{\tilde{M}}(\tilde{E})
$$

for $t \in\left[0, b\left[\right.\right.$ and all planes $E \subset T_{\gamma(t)} M$ with $\dot{\gamma}(t) \in E$, resp. $\tilde{E} \subset T_{\tilde{\gamma}(t)} \tilde{M}$ with $\dot{\tilde{\gamma}}(t) \in \tilde{E}$, then for any isotone $C^{2}$-function $f:[0, \infty[\rightarrow \mathbb{R}$ :

$$
\begin{equation*}
\left.\nabla d(f \circ r)_{\gamma(t)} \succcurlyeq \nabla d(f \circ \tilde{r})_{\tilde{\gamma}(t)}, \quad t \in\right] 0, b[, \tag{2.1.31}
\end{equation*}
$$

where $r=d\left(x_{0}, \cdot\right)$ and $\tilde{r}=\tilde{d}\left(\tilde{x}_{0}, \cdot\right)$ denote the distance functions from $x_{0}=\gamma(0)$ in $M$, resp. from $\tilde{x}_{0}=\tilde{\gamma}(0)$ in $\tilde{M}$.

REMARK 2.1.45. (1) The assumed minimality of $\gamma:[0, b] \rightarrow \tilde{M}$ has as consequence that for $t<b$ no $\gamma(t)$ is a cut point of $\gamma(0)$; analogously for $\tilde{\gamma}$. Obviously (2.1.31) also holds for $t=b$, if $r$ and $\tilde{r}$ are differentiable at $\gamma(b)$, resp. at $\tilde{\gamma}(b)$.
(2) It would be sufficient to assume that there are no conjugate points to $\tilde{\gamma}(0)$ along $\tilde{\gamma}$. By the Comparison Theorem of Rauch, along with Theorem 2.1.34, then also no $\gamma(t)$ is conjugate to $\gamma(0)$ along $\gamma$. However $\exp _{x_{0}} \mid V$ and $\exp _{\tilde{x}_{0}} \mid \tilde{V}$ may then be no longer invertible; for fixed $t$ then $\nabla d(f \circ r)_{\gamma(t)}$ needs to be replaced by $\rho=\left(\exp _{x_{0}} \mid V_{\text {loc }}\right)^{-1}$ with $V_{\text {loc }}$ a sufficiently small neighbourhood of $t \dot{\gamma}(0)$, which is mapped by $\exp _{x_{0}}$ diffeomorphically to an open neighbourhood of $\gamma(t)$; the right-hand side of (2.1.31) should then be interpreted correspondingly.

Proof of Theorem 2.1.44. Let $\left.\left.t_{0} \in\right] 0, b\right]$, to this $x=\gamma\left(t_{0}\right)$ and $\tilde{x}=\tilde{\gamma}\left(t_{0}\right)$. Furthermore, let $u \in T_{x} M$ and $\iota: \gamma^{*} T M \rightarrow \tilde{\gamma}^{*} T \tilde{M}$ an isometric bundle embedding over $\mathbb{R}$ constructed according to (2.1.29). We then have to show that

$$
\nabla d(f \circ r)_{x}(u, u) \geq \nabla d(f \circ \tilde{r})_{\tilde{x}}\left(\iota_{t_{0}} u, \iota_{t_{0}} u\right) .
$$

By formula (1.7.2) one obtains at first

$$
\nabla d(f \circ r)_{x}(u, u)=f^{\prime \prime}\left(t_{0}\right)(d r)_{x}(u)+f^{\prime}\left(t_{0}\right)(\nabla d r)_{x}(u, u)
$$

$$
\nabla d(f \circ \tilde{r})_{\tilde{x}}\left(\iota_{t_{0}} u, \iota_{t_{0}} u\right)=f^{\prime \prime}\left(t_{0}\right)(d \tilde{r})_{\tilde{x}}\left(\iota_{t_{0}} u\right)+f^{\prime}\left(t_{0}\right)(\nabla d \tilde{r})_{\tilde{x}}\left(\iota_{t_{0}} u, \iota_{t_{0}} u\right) .
$$

The first two summands are equal by Lemma 2.1.42; since by assumption $f^{\prime}\left(t_{0}\right) \geq 0$, it remains to show that $(\nabla d r)_{x}(u, u) \geq(\nabla d \tilde{r})_{\tilde{x}}\left(\iota_{t_{0}} u, \iota_{t_{0}} u\right)$, by (2.1.28) hence to verify that

$$
I_{t_{0}}\left(X^{\perp}, X^{\perp}\right) \geq I_{t_{0}}\left(\tilde{X}^{\perp}, \tilde{X}^{\perp}\right)
$$

here $X \in \Gamma\left(\gamma^{*} T M\right)$ is the Jacobi field with $X(0)=0, X\left(t_{0}\right)=u$ and $\tilde{X} \in \Gamma\left(\tilde{\gamma}^{*} T \tilde{M}\right)$ the Jacobi field with $\tilde{X}(0)=0, \tilde{X}\left(t_{0}\right)=\iota_{t_{0}} u$. In terms of the vector field $Y:=\iota X^{\perp} \in$ $\Gamma\left(\tilde{\gamma}^{*} T \tilde{M}\right)$ however, it holds

$$
I_{t_{0}}\left(X^{\perp}, X^{\perp}\right) \geq I_{t_{0}}(Y, Y) \geq I_{t_{0}}\left(\tilde{X}^{\perp}, \tilde{X}^{\perp}\right)
$$

and thus the claim: The first inequality follows directly from the Definition of the index form, combined with the observations that $\left|\nabla_{D} X^{\perp}\right|=\left|\nabla_{D} Y\right|$ and that $\left\langle R\left(X^{\perp}, \dot{\gamma}\right) \dot{\gamma}, X^{\perp}\right\rangle \leq$ $\langle R(Y, \dot{\tilde{\gamma}}) \dot{\tilde{\gamma}}, Y\rangle$ by the assumptions on the sectional curvatures; the second inequality is a consequence of the Index Lemma (Theorem 2.1.37), since $Y$ equals $\tilde{X}^{\perp}$ at 0 and $t_{0}$.

Corollary 2.1.46 (Comparison Theorem for the Laplacian: basic version). Let $(M, g),(\tilde{M}, \tilde{g})$ be Riemannian manifolds with $2 \leq \operatorname{dim} M \leq \operatorname{dim} \tilde{M}$ and let $\gamma:[0, b] \rightarrow$ $M$, resp., $\tilde{\gamma}:[0, b] \rightarrow \tilde{M}$ be minimal normal geodesic curves. If then

$$
\operatorname{Riem}^{M}(E) \leq \operatorname{Riem}^{\tilde{M}}(\tilde{E})
$$

for all planes $E \subset T_{\gamma(t)} M$ with $\dot{\gamma}(t) \in E$, resp.. $\tilde{E} \subset T_{\tilde{\gamma}(t)} \tilde{M}$ with $\dot{\tilde{\gamma}}(t) \in \tilde{E}$ and all $t \in[0, b]$, then for each isotone $C^{2}$-function $f:[0, \infty[\rightarrow \mathbb{R}$ the inequality

$$
\Delta(f \circ r)(\gamma(t)) \geq \Delta(f \circ \tilde{r})(\tilde{\gamma}(t)), \quad t \in] 0, b[
$$

holds, where $r=d\left(x_{0}, \cdot\right)$ and $\tilde{r}=\tilde{d}\left(\tilde{x}_{0}, \cdot\right)$ denote the distance functions from $x_{0}=\gamma(0)$ in $M$, resp. from $\tilde{x}_{0}=\tilde{\gamma}(0)$ in $\tilde{M}$.

Proof. The claim follows from Theorem 2.1.44 by taking trace.
Comparison theorems are typically applied by comparing a given Riemannian manifold to simply structured standard manifolds. This procedure obviously depends on the explicit knowledge of suitable comparison manifolds. An important type of model manifolds are covered by the following definition (see [14]).

DEFINITION 2.1.47 (Model, rotationally symmetric manifold). Let $(\mathbb{M}, g)$ be an $n$ dimensional $(n \geq 2)$ Riemannian manifold and $0 \in \mathbb{M}$ be a distinguished point. Then $(\mathbb{M}, g)$ is called a model with center 0 if 0 is a pole for $(\mathbb{M}, g)$ with $\mathbb{M}$ being rotationally symmetric about 0 in the sense that each linear isometry $\varphi: T_{0} \mathbb{M} \rightarrow T_{0} \mathbb{M}$ is he differential of an isometry $\phi: \mathbb{M} \rightarrow \mathbb{M}$, i.e., such that $\phi(0)=0$ and $(d \phi)_{0}=\varphi$.

Before entering the discussion on properties of models, we want to collect some facts about isometries. By Definition 1.5.4, isometries are local isometries with the additional property that they are diffeomorphisms.

REMARK 2.1.48 (on local isometries). Let $(M, g)$ be a metrically complete Riemannian manifold and $\phi: M \rightarrow M$ a local isometry, i.e., $\phi^{*} g=g$. Then:
(i) $\phi$ preserves the length of curves, and hence, if $\phi$ is even an isometry, then also distances, i.e., then it holds: $d(\phi(x), \phi(y))=d(x, y)$ for $x, y \in M$.
(ii) $\phi$ transfers geodesics in geodesics; hence in particular:

$$
\phi \circ \exp _{x}(t v)=\exp _{\phi(x)}\left(t \phi_{*} v\right), \quad x \in M, v \in T_{x} M
$$

(iii) $\phi$ preserves the Levi-Civita connection, i.e.,

$$
d \phi \nabla_{A} B=\nabla_{A}(d \phi B), \quad A, B \in \Gamma(T M)
$$

In particular, for vector fields $X$ along a curve $c$, it holds $d \phi \nabla_{D} X=\nabla_{D}(d \phi X)$, and $X$ is hence parallel along $c$ if and only if $d \phi X$ is parallel along $\phi \circ c$.
(iv) $\phi$ preserves the Riemannian sectional curvature, i.e., if $E=\operatorname{span}\{v, w\} \subset T_{x} M$ and $E^{\prime}=\operatorname{span}\left\{\phi_{*} v, \phi_{*} w\right\} \subset T_{\phi(x)} M$, then $\operatorname{Riem}_{x}^{M}(E)=\operatorname{Riem}_{\phi(x)}^{M}\left(E^{\prime}\right)$.

Proof. (i) is a direct consequence of the definition of the length functional. (ii) follows from (i) since $\phi$ is a local diffeomorphism. (iii) follows from formula (1.7.1) and the observation that $\phi$ is affine, since $\phi$ maps geodesics to geodesics. (iv) finally is a consequence of (iii) and the second Cartan structural equation (see Theorem 1.4.27).

THEOREM 2.1.49. Let $(\mathbb{M}, g)$ be a model and $\gamma:[0, b] \rightarrow \mathbb{M}$ a geodesic curve emanating from the distinguished point $0 \in \mathbb{M}$. Then each each proper Jacobi field along $\gamma$, which vanishes at 0 , is up to a scalar function a parallel vector field along $\gamma$. In particular, two Jacobi fields along $\gamma$, vanishing at 0 , are already orthogonal along $\gamma$ if they are orthogonal at one place.

Proof. Let $J$ be a Jacobi field along $\gamma$ such that $\langle J, \dot{\gamma}\rangle=0$ and $J(0)=0$; denote $v=\dot{\gamma}(0) \in T_{0} \mathbb{M}$ and $w=J^{\prime}(0)$. Since $J$ is a proper Jacobi field, we have $v \perp w$ in $T_{0} \mathbb{M}$. With the identifications $\mathbb{M} \cong T_{0} \mathbb{M}$ via $\exp _{0}$, and correspondingly $T_{\gamma(t)} \mathbb{M} \cong T_{t v} T_{0} \mathbb{M} \cong$ $T_{0} \mathbb{M}$, it holds that $\gamma(t)=\exp _{0}(t v) \equiv t v$ and $J(t)=\left(d \exp _{0}\right)_{t v}(t w) \equiv t w \in T_{t v} \mathbb{M}$. We have to show that the vector field $W$ along $\gamma$ given by $W(t):=w \in T_{t v} \mathbb{M}$, coincides up to multiplication by a scalar function with the parallel transport of $w \in T_{0} \mathbb{M}$ along $\gamma$. To this end, we consider the two-dimensional submanifold

$$
\mathbb{M}_{0}:=\exp _{0}(\mathbb{R} v+\mathbb{R} w) \cong \mathbb{R} v+\mathbb{R} w \subset \mathbb{M}
$$

with the induced Riemannian metric. Now $\gamma$ is also a geodesic in $\mathbb{M}_{0}$ and it holds $W(t) \perp$ $\dot{\gamma}(t)$ in $T_{0} \mathbb{M}_{0}$. By the isometry of the parallel transport with respect to the Levi-Civita connection, then $\bar{W}:=W /|W|$ must be parallel along $\gamma$ in $\mathbb{M}_{0}$. It remains to show that $\bar{W}$ is also parallel along $\gamma$ in $\mathbb{M}$. To this end, it is sufficient to show that $\mathbb{M}_{0}$ as submanifold of $\mathbb{M}$ is totally geodesic, since the inclusion $\iota: \mathbb{M}_{0} \hookrightarrow \mathbb{M}$ is affine and then

$$
\nabla_{D} \iota_{*} \bar{W}=\iota_{*} \nabla_{D} \bar{W}=0
$$

By definition, we have $\mathbb{M}_{0}=\exp _{0}(\mathbb{R} v+\mathbb{R} w)$ with $v \perp w$ in $T_{0} \mathbb{M}$. We choose a linear isometry $\varphi: T_{0} \mathbb{M} \rightarrow T_{0} \mathbb{M}$ with $\varphi(v)=v$ and $\varphi(w)=w$, but $\varphi(u) \neq u$ for any $u \in$ $\{\mathbb{R} v+\mathbb{R} w\}^{\perp}$. Then there is an isometry $\phi: \mathbb{M} \rightarrow \mathbb{M}$ with $\phi(0)=0$ and $d \phi_{0}=\varphi$. According to Remark 2.1.48 (ii), $\mathbb{M}_{0}$ is the fixed point set of $\phi$, i.e., $\mathbb{M}_{0}=\{x \in \mathbb{M}$ : $\phi(x)=x\}$. This already shows the claim since if $c$ is a geodesic in $\mathbb{M}$ with $c(0) \in \mathbb{M}_{0}$, i.e. $c(t)=\exp _{c(0)}(t \dot{c}(0))$, then $c$ lies totally in $\mathbb{M}_{0}$ if and only if $\phi \circ c=c$; because of $(\phi \circ c)(t)=\exp _{c(0)}\left(t \phi_{*} \dot{c}(0)\right)$ this is however the case exactly if $\phi_{*} \dot{c}(0)=\dot{c}(0)$, or in other words, if $\dot{c}(0) \in T_{c(0)} \mathbb{M}_{0}$.

In general, we have for $X \in T_{x} \mathbb{M}$ with $x \in \mathbb{M}_{0}$ that $X \in T_{x} \mathbb{M}_{0}$ if and only if $X=\dot{\beta}(0)$ for a curve $\beta$ in $M_{0}$ with $\beta(0)=x$; indeed, by $\phi \circ \beta=\beta$ this condition implies $\phi_{*} X=X$, conversely from $\phi_{*} X=X$ the existence of an $M_{0}$-valued curve $\beta$ follows with $X=\dot{\beta}(0)$, e.g. $\beta(t)=\exp _{x}(t X)$ according to Remark (2.1.48) (ii).

Lemma 2.1.50. Let $(\mathbb{M}, g)$ be a model and $0 \in \mathbb{M}$ its center. Fix $x, \tilde{x} \in \mathbb{M}$ such that $r=d(x, 0)=d(\tilde{x}, 0)$, and let $\gamma$, resp. $\tilde{\gamma}$, be the normal geodesic curves emanating from 0 with the property that $\gamma(r)=x$ and $\tilde{\gamma}(r)=\tilde{x}$. If then $\left(u_{1}, \ldots, u_{d}\right)$ is an orthonormal
basis for $T_{x} \mathbb{M}$ with $u_{1}=\dot{\gamma}(r)$, and analogously $\left(\tilde{u}_{1}, \ldots, \tilde{u}_{d}\right)$ an orthonormal basis for $T_{\tilde{x}} \mathbb{M}$ with $\tilde{u}_{1}=\dot{\tilde{\gamma}}(r)$, then there is an isometry $\phi: \mathbb{M} \rightarrow \mathbb{M}$ with $\phi(x)=\tilde{x}$, such that

$$
d \phi_{x} u_{i}=\tilde{u}_{i}, \quad i=1, \ldots, d .
$$

Proof. Let $v:=u_{1}=\dot{\gamma}(r)$ and $\tilde{v}:=\tilde{u}_{1}=\dot{\tilde{\gamma}}(r)$. We identify $\mathbb{M} \cong T_{0} \mathbb{M}$, so that $x=\exp _{0}(r \dot{\gamma}(0)) \equiv r \dot{\gamma}(0)$ und

$$
T_{x} \mathbb{M} \cong T_{r \dot{\gamma}(0)} T_{0} \mathbb{M} \cong T_{0} \mathbb{M}
$$

In this sense we understand $u_{1}, \ldots, u_{n}$ as elements of $T_{0} \mathbb{M}$; in particular then $v \equiv \dot{\gamma}(0)$. By the Gauss Lemma, we have $v \perp u_{i}$ for $i=2, \ldots, n$ in $T_{0} \mathbb{M}$, and correspondingly $\tilde{v} \perp \tilde{u}_{i}$ for $i=2, \ldots, n$ in $T_{0} \mathbb{M}$. On the other hand, the Jacobi fields $J_{2}, \ldots, J_{n}$ with

$$
J_{i}(t)=\left(d \exp _{0}\right)_{t v}\left(t u_{i}\right) \equiv t u_{i} \in T_{t v} \mathbb{M}
$$

are pairwise orthogonal along $\gamma$ by Theorem 2.1.49; hence also $\left(u_{2}, \ldots, u_{n}\right)$ is orthogonal in $T_{0} \mathbb{M}$. With the same argument one obtains the orthogonality of $\left(\tilde{u}_{2}, \ldots, \tilde{u}_{n}\right)$ in $T_{0} \mathbb{M}$. Thus we can find a linear isometry $\varphi: T_{0} \mathbb{M} \rightarrow T_{0} \mathbb{M}$ such that $\varphi(v)=\tilde{v}$ and $\varphi\left(u_{i}\right)=\lambda_{i} \tilde{u}_{i}$ for $i=2, \ldots, n$ where $\lambda_{i}>0$. Since $\mathbb{M}$ is a model, there is an isometry $\phi: \mathbb{M} \rightarrow \mathbb{M}$ such that $\phi(0)=0$ and $d \phi_{0}=\varphi$. Because of

$$
\phi \circ \exp _{0}(t v)=\exp _{\phi(0)}\left(t \phi_{*} v\right)=\exp _{\phi(0)}(t \varphi v)=\exp _{\phi(0)}(t \tilde{v})
$$

we have $\phi \circ \gamma=\tilde{\gamma}$, in particular then $\phi(x)=\tilde{x}$ and $d \phi_{x} v=\tilde{v}$.
It remains to verify that $d \phi_{x} u_{i}=\tilde{u}_{i}$ for $i=2, \ldots, n$. To this end, we consider for fixed $i$ the Jacobi field $J$ along $\gamma$ with $J(0)=0, J^{\prime}(0)=u_{i}$, and analogously $\tilde{J}$ the Jacobi field along $\tilde{\gamma}$ with $\tilde{J}(0)=0, \tilde{J}^{\prime}(0)=\tilde{u}_{i}$. According to Example 2.1.32, we have

$$
J(t)=\left.\frac{\partial}{\partial s}\right|_{s=0} \exp _{0}\left(t\left(v+s u_{i}\right)\right)=\left(d \exp _{0}\right)_{t v}\left(t u_{i}\right) \equiv t u_{i} \in T_{t v} \mathbb{M}
$$

and $\tilde{J}(t) \equiv t \tilde{u}_{i} \in T_{t \tilde{v}} \mathbb{M}$. But we have

$$
\phi \circ \exp _{0}\left(t\left(v+s u_{i}\right)\right)=\exp _{0}\left(t \varphi\left(v+s u_{i}\right)\right)=\exp _{0}\left(t\left(\tilde{v}+s \lambda_{i} \tilde{u}_{i}\right)\right)
$$

from where by differentiating with respect to $s$ at $s=0$ the relation

$$
(d \phi)_{t v} J(t)=\lambda_{i} \tilde{J}(t)
$$

is derived, thus $(d \phi)_{t v} t u_{i}=\lambda_{i} t \tilde{u}_{i}$. This shows in particular that $(d \phi)_{r v} u_{i} \equiv d \phi_{x} u_{i}=$ $\lambda_{i} \tilde{u}_{i}$. By the isometry of $d \phi_{x}: T_{x} \mathbb{M} \rightarrow T_{\tilde{x}} \mathbb{M}$, then necessarily $\lambda_{i}=1$.

Let $(M, g)$ be a metrically complete Riemannian manifold and $x_{0} \in M$. The radial vector field $\frac{\partial}{\partial r}$ defined on $M \backslash\left(C\left(x_{0}\right) \cup\left\{x_{0}\right\}\right)$ is given by $\left(\frac{\partial}{\partial r}\right)_{x}=\dot{\gamma}\left(t_{0}\right)$ where $\gamma$ denotes the unique minimal normal geodesic such that $\gamma(0)=x_{0}$ and $\gamma\left(t_{0}\right)=x$. Obviously, it holds $\frac{\partial}{\partial r}=\operatorname{grad} r$ on $M \backslash\left(C\left(x_{0}\right) \cup\left\{x_{0}\right\}\right)$ with $r:=d\left(x_{0}, \cdot\right)$, since $\operatorname{grad} r$ is determined by $\langle\operatorname{grad} r, Y\rangle=Y(r)$ for each vector field $Y$ on $M \backslash\left(C\left(x_{0}\right) \cup\left\{x_{0}\right\}\right)$; on the other hand, one has $Y=Y(r) \frac{\partial}{\partial r}+Y^{\perp}$ with $\left\langle\frac{\partial}{\partial r}, Y^{\perp}\right\rangle=0$ by the Gauss Lemma, so that also $\left\langle\frac{\partial}{\partial r}, Y\right\rangle=Y(r)$ holds.

REMARK 2.1.51. Note that for an arbitrary vector field $X$ on $M \backslash\left(C\left(x_{0}\right) \cup\left\{x_{0}\right\}\right)$ it always holds that

$$
\begin{equation*}
(\nabla d r)\left(X, \frac{\partial}{\partial r}\right)=0 \tag{2.1.32}
\end{equation*}
$$

Indeed one has $(\nabla d r)\left(X, \frac{\partial}{\partial r}\right)=X\left(\frac{\partial}{\partial r} r\right)-\left(\nabla_{X} \frac{\partial}{\partial r}\right) r$, and because of $\frac{\partial}{\partial r} r=1$, then

$$
(\nabla d r)\left(X, \frac{\partial}{\partial r}\right)=-\left\langle\nabla_{X} \frac{\partial}{\partial r}, \operatorname{grad} r\right\rangle=-\left\langle\nabla_{X} \frac{\partial}{\partial r}, \frac{\partial}{\partial r}\right\rangle=-\frac{1}{2} X\left\langle\frac{\partial}{\partial r}, \frac{\partial}{\partial r}\right\rangle=0
$$

Instead of $\frac{\partial}{\partial r}$ we also write occasionally also $\partial$ or $\partial^{M}$. In the particular case that $x_{0}$ is a pole for $(M, g)$, the corresponding radial vector field $\partial^{M}$ is defined on $M \backslash\left\{x_{0}\right\}$.

DEFINITION 2.1.52. Under radial curvature of $(M, g)$ with respect to $x_{0}$ we understand the restriction of the sectional curvature Riem $^{M}$ to radial planes, i.e. planes $E \subset T_{x} M$ such that $\partial_{x}^{M}=\left(\frac{\partial}{\partial r}\right)_{x} \in E$. Planes in $T_{x_{0}} M$ are considered as radial by convention.

REMARK 2.1.53. In a model $(\mathbb{M}, g)$ with with center 0 , the radial curvature at some point $x$ depends only on $r=r_{M}(x)$ where $r_{M}(x)=d(0, x)$.

Proof. This is a consequence of Lemma 2.1.50 and Remark (2.1.48) (iv) which says that isometries preserve the Riemannian sectional curvature.

DEFINITION 2.1.54 (Radial curvature function). Let $(\mathbb{M}, g)$ be a model. The function

$$
k_{\mathbb{M}}: \mathbb{R}_{+} \rightarrow \mathbb{R}, \quad k_{\mathbb{M}}(t):=\text { radial curvature at } x \in \mathbb{M} \text { with } r_{M}(x)=t
$$

is well-defined by Remark 2.1.53 and called radial curvature function of the model $(\mathbb{M}, g)$.
The Comparison Theorem for the Laplacian (Theorem 2.1.46) takes a simpler form in the case of a model as comparison manifold: it is then sufficient to compare the Ricci curvature along normal geodesics.

THEOREM 2.1.55 (Laplacian Comparison Theorem: special version). Let ( $M, g$ ) be a metrically complete Riemannian manifold with $n=\operatorname{dim} M \geq 2$, and $x_{0} \in M$, as well as $\mathbb{M}$ a model of the same dimension with center $0 \in \mathbb{M}$. Let $r_{M}=d\left(x_{0}, \cdot\right)$ and $r_{\mathbb{M}}=d(0, \cdot)$ be the distance functions to $x_{0}$ in $M$, resp. to 0 in $\mathbb{M}$, and $\partial^{M}$ resp. $\partial^{\mathbb{M}}$ the corresponding radial vector fields. Suppose that for some $R>0$,

$$
\operatorname{Ric}^{M}\left(\partial^{M}, \partial^{M}\right)(x) \geq \operatorname{Ric}^{\mathbb{M}}\left(\partial^{\mathbb{M}}, \partial^{\mathbb{M}}\right)(\tilde{x})\left(\equiv(n-1) k_{\mathbb{M}}(r)\right)
$$

for all $x \in M \backslash\left(C\left(x_{0}\right) \cup\left\{x_{0}\right\}\right)$ and $\tilde{x} \in \mathbb{M} \backslash\{0\}$ with $r=r_{M}(x)=r_{\mathbb{M}}(\tilde{x})<R$. Then for each isotone $C^{2}$-function $f:[0, R[\rightarrow \mathbb{R}$ and all $x, \tilde{x}$ as above, it holds

$$
\Delta\left(f \circ r_{M}\right)(x) \leq \Delta\left(f \circ r_{\mathbb{M}}\right)(\tilde{x})
$$

Proof. We follow the proof of Theorem 2.1.44. At first we remark that it is again sufficient to consider the case $f(r)=r$, since from

$$
\nabla d\left(f \circ r_{M}\right)_{x}(u, u)=\left(f^{\prime \prime} \circ r_{M}\right)\left(d r_{M}\right)_{x}(u)+\left(f^{\prime} \circ r_{M}\right)\left(\nabla d r_{M}\right)_{x}(u, u)
$$

for $u \in T_{x} M$ one concludes immediately $\Delta\left(f \circ r_{M}\right)=f^{\prime \prime} \circ r_{M}+\left(f^{\prime} \circ r_{M}\right) \Delta r_{M}$, and analogously $\Delta\left(f \circ r_{\mathbb{M}}\right)=f^{\prime \prime} \circ r_{\mathbb{M}}+\left(f^{\prime} \circ r_{\mathbb{M}}\right) \Delta r_{\mathbb{M}}$.

Let now $x \in M \backslash\left(C\left(x_{0}\right) \cup\left\{x_{0}\right\}\right)$ and $\tilde{x} \in \mathbb{M} \backslash\{0\}$ such that $r=r_{M}(x)=r_{\mathbb{M}}(\tilde{x})<$ $R$; to this let $\gamma:[0, r] \rightarrow M$ be the geodesic emanating from $x_{0}$ with $\gamma(r)=x$, and correspondingly $\tilde{\gamma}:[0, r] \rightarrow \mathbb{M}$ the geodesic emanating from 0 with $\tilde{\gamma}(r)=\tilde{x}$. We fix orthonormal bases $\left(u_{1}, \ldots, u_{n}\right)$ for $T_{x} M$ where $u_{1}=\partial_{x}^{M}$ and $\left(\tilde{u}_{1}, \ldots, \tilde{u}_{n}\right)$ for $T_{\tilde{x}} \mathbb{M}$ where $\tilde{u}_{1}=\partial_{\tilde{x}}^{\mathbb{M}}$. For $i=2, \ldots, n$ let $X_{i}$ be unique (proper) Jacobi field along $\gamma$ such that $X_{i}(0)=0, X_{i}(r)=u_{i}$, and analogously $\tilde{X}_{i}$ the corresponding Jacobi field along $\tilde{\gamma}$, such that $\tilde{X}_{i}(0)=0$, as well as $\tilde{X}_{i}(r)=\tilde{u}_{i}$. Taking (2.1.32) into account, one obtains

$$
\Delta r_{M}(x)=\sum_{i=2}^{n}\left(\nabla d r_{M}\right)\left(u_{i}, u_{i}\right)=\sum_{i=2}^{n} I\left(X_{i}, X_{i}\right)
$$

respectively,

$$
\Delta r_{\mathbb{M}}(\tilde{x})=\sum_{i=2}^{n}\left(\nabla d r_{\mathbb{M}}\right)\left(\tilde{u}_{i}, \tilde{u}_{i}\right)=\sum_{i=2}^{n} I\left(\tilde{X}_{i}, \tilde{X}_{i}\right) .
$$

But $\mathbb{M}$ is a model so that for any $0<t \leq r$ the vectors $\tilde{X}_{i}(t)$ pairwise orthogonal and in addition $\left|\tilde{X}_{i}(t)\right|=\left|\tilde{X}_{j}(t)\right|$ : indeed, the first by Theorem 2.1.49 and the second claim, since each $\tilde{X}_{i}(t)$ can be transfered to $\tilde{X}_{j}(t)$ by the differential of an isometry which lets $\tilde{\gamma}$ invariant. Hence we have

$$
\begin{aligned}
\Delta r_{\mathbb{M}}(\tilde{x})=\sum_{i=2}^{n} I\left(\tilde{X}_{i}, \tilde{X}_{i}\right) & =\sum_{i=2}^{n} \int_{0}^{r}\left\{\left|\nabla_{D} \tilde{X}_{i}\right|^{2}-\left\langle R\left(\tilde{X}_{i}, \partial^{\mathbb{M}}\right) \partial^{\mathbb{M}}, \tilde{X}_{i}\right\rangle\right\} d t \\
& =\int_{0}^{r}\left\{\sum_{i=2}^{n}\left|\nabla_{D} \tilde{X}_{i}\right|^{2}-\left|\tilde{X}_{2}\right|^{2} \operatorname{Ric}^{\mathbb{M}}\left(\partial^{\mathbb{M}}, \partial^{\mathbb{M}}\right)\right\} d t .
\end{aligned}
$$

If now $\iota_{r}: T_{\tilde{x}} \mathbb{M} \rightarrow T_{x} M$ is the linear isometry such that $\iota_{r}\left(\tilde{u}_{i}\right)=u_{i}$ for $i=1, \ldots, n$ and if one extends it canonically via parallel transport according to

$$
\begin{aligned}
& T_{\tilde{\gamma}(r)} \mathbb{M} \stackrel{\iota_{r}}{\longrightarrow} T_{\gamma(r)} M \\
& / /_{t, r} \uparrow \downarrow \quad \imath \uparrow / /_{t, r} \\
& T_{\tilde{\gamma}(t)} \mathbb{M} \stackrel{\iota_{t}}{\iota_{-}} T_{\gamma(t)} M,
\end{aligned}
$$

then one obtains an isometric bundle embedding $\iota: \tilde{\gamma}^{*} T \mathbb{M} \rightarrow \gamma^{*} T M$ over $\mathbb{R}$, which commutes with the covariant derivative $\nabla_{D}$ of vector fields and transfers the radial vector field along $\tilde{\gamma}$ into the radial vector field along $\gamma$. One applies now again the index lemma (Theorem 2.1.37):

$$
\begin{aligned}
\Delta r_{M}(x) & =\sum_{i=2}^{n} I\left(X_{i}, X_{i}\right) \leq \sum_{i=2}^{n} I\left(\iota \tilde{X}_{i}, \iota \tilde{X}_{i}\right) \\
& =\sum_{i=2}^{n} \int_{0}^{r}\left\{\left|\nabla_{D} \iota \tilde{X}_{i}\right|^{2}-\left\langle R\left(\iota \tilde{X}_{i}, \partial^{M}\right) \partial^{M}, \iota \tilde{X}_{i}\right\rangle\right\} d t \\
& =\int_{0}^{r}\left\{\sum_{i=2}^{n}\left|\iota \nabla_{D} \tilde{X}_{i}\right|^{2}-\sum_{i=2}^{n}\left|\iota \tilde{X}_{i}\right|^{2} \operatorname{Ric}^{M}\left(\partial^{M}, \partial^{M}\right)\right\} d t \\
& =\int_{0}^{r}\left\{\sum_{i=2}^{n}\left|\nabla_{D} \tilde{X}_{i}\right|^{2}-\left|\tilde{X}_{2}\right|^{2} \operatorname{Ric}^{M}\left(\partial^{M}, \partial^{M}\right)\right\} d t \\
& \leq \int_{0}^{r}\left\{\sum_{i=2}^{n}\left|\nabla_{D} \tilde{X}_{i}\right|^{2}-\left|\tilde{X}_{2}\right|^{2} \operatorname{Ric}^{\mathbb{M}}\left(\partial^{\mathbb{M}}, \partial^{\mathbb{M}}\right)\right\} d t=\Delta r_{\mathbb{M}}(\tilde{x})
\end{aligned}
$$

where the last inequality comes from the assumption on the Ricci curvature.
An important tool for the explicit description and construction of models are Euclidean spheres, even if themselves they are not covered by the class of models. For $a>0$ let

$$
\mathbb{S}_{a}^{n}:=\left\{x \in \mathbb{R}^{n+1}:|x|=a\right\} \Longleftrightarrow \mathbb{R}^{n+1}
$$

be the sphere of radius $a$, equipped with the Riemannian metric $g=\iota^{*}$ eucl induced from $\left(\mathbb{R}^{n+1}\right.$, eucl) (see Example 1.5.12), where occasionally for historical reasons the metric is written $g=d \vartheta^{2}$.

For $x \in \mathbb{S}_{a}^{n}$ we identify canonically

$$
T_{x} \mathbb{S}_{a}^{n} \cong\{x\}^{\perp} \subset \mathbb{R}^{n+1}
$$

where $\{x\}^{\perp}$ is the orthogonal complement of $\mathbb{R} x$ in $\mathbb{R}^{n+1}$. Each orthogonal transformation $A \in \mathrm{O}(n+1)$ defines by restriction

$$
A \mid \mathbb{S}_{a}^{n}: \mathbb{S}_{a}^{n} \rightarrow \mathbb{S}_{a}^{n}
$$

an isometry of $\mathbb{S}_{a}^{n}$. In particular, $\left(\mathbb{S}_{a}^{n}, d \vartheta^{2}\right)$ has constant sectional curvature, since for given $x, y \in \mathbb{S}_{a}^{n}$ and orthonormal vectors $u_{1}, u_{2} \in T_{x} \mathbb{S}_{a}^{n}$, resp. $v_{1}, v_{2} \in T_{y} \mathbb{S}_{a}^{n}$, there is an orthogonal transformation $A \in \mathrm{O}(n+1)$ such that $A x=y$ and $A_{*} u_{i}=v_{i}$ for $i=1,2$. Modulo multiplication by a constant, $d \vartheta^{2}$ is however the only Riemannian metric on $\mathbb{S}_{a}^{n}$ invariant under the full orthogonal group $\mathrm{O}(n+1)$. For $a=1$ we write simply $\mathbb{S}^{n}$ instead of $\mathbb{S}_{1}^{n}$.

REMARK 2.1.56. As a compact manifold $\left(\mathbb{S}_{a}^{n}, d \vartheta^{2}\right)$ is metrically complete and the maximal geodesics coincide with the great circles on $\mathbb{S}_{a}^{n}$. For instance, fix $x, y \in \mathbb{S}_{a}^{n}, x \neq y$ and $r=d(x, y)$ sufficiently small such that there is exactly one geodesic $\gamma:[0, r] \rightarrow \mathbb{S}_{a}^{n}$ with $\gamma(0)=x$ and $\gamma(r)=y$. To the plane $E=\mathbb{R} x+\mathbb{R} y$ consider now an orthogonal transformation $A \in \mathrm{O}(n+1)$ which has $E$ as fixed point set, e.g. the mirror map at $E$. Then also $A \circ \gamma$ is a minimal geodesic connecting $x$ and $y$, hence $A \circ \gamma=\gamma$ and $\gamma$ lies on the great circle $E \cap \mathbb{S}_{a}^{n}$.

Let now $(\mathbb{M}, g)$ again be a model and $0 \in \mathbb{M}$ its center. Then

$$
\exp _{0}:\left(T_{0} \mathbb{M}, \exp _{0}^{*} g\right) \rightarrow(\mathbb{M}, g)
$$

defines an isometry of Riemannian manifolds. Without restrictions, we may identify $\mathbb{M} \cong$ $\mathbb{R}^{n}$, where the center $0 \in \mathbb{M}$ corresponds to the origin in $\mathbb{R}^{n}$ and where we identify $T_{0} \mathbb{M}$ and $\mathbb{R}^{n}$ isometically as Euclidean $\mathbb{R}$-vector spaces. The metric $\exp _{0}^{*} g$ restricted to $\mathbb{R}^{n} \backslash\{0\}$ takes under pull-back with $] 0, \infty\left[\times \mathbb{S}^{n-1} \xrightarrow{\sim} \mathbb{R}^{n} \backslash\{0\},(r, v) \mapsto r v\right.$ by the Gauss Lemma the form $d r \otimes d r+h_{r}$ where $h_{r}$ denotes the metric on the $(n-1)$-dimensional unit sphere $\mathbb{S}^{n-1}$ induced by $\exp _{0}^{*} g$ on $\mathbb{S}_{r}^{n-1}$.

Recall that $\mathbb{M}$ is a model and that the metric $h_{r}$ is hence invariant under the full $n$ dimensional orthogonal group: thus $h_{r}$ coincides up to a positive constant (depending on $r$ ) with the standard metric $d \vartheta^{2}$ on $\mathbb{S}^{n-1}$. We write $h_{r}=f(r)^{2} d \vartheta^{2}$. With the positive function $f:] 0, \infty[\rightarrow \mathbb{R}$ defined in this way, each $n$-dimensional model $(\mathbb{M}, g)$ takes the form

$$
\left(\mathbb{R}^{n}, d r \otimes d r+f(r)^{2} d \vartheta^{2}\right)
$$

THEOREM 2.1.57 (Elementary properties of models). Let ( $\mathbb{M}, g$ ) be a model with $\mathbb{M} \cong \mathbb{R}^{n}$ and $g=d r \otimes d r+f(r)^{2} d \vartheta^{2}$ on $\mathbb{R}^{n} \backslash\{0\}$, as well as $k=k_{\mathbb{M}}$ the radial curvature function of $\mathbb{M}$. Then the following items hold:
(i) (Jacobi equation) $f^{\prime \prime}(t)+k(t) f(t) \equiv 0$ with $f(0)=0$ and $f^{\prime}(0)=1$.
(ii) $\nabla d r=\left(\left(f^{\prime} / f\right) \circ r\right)(g-d r \otimes d r)$ with $r=d(0, \cdot)$ the radial function of the model. In particular, it holds that

$$
\Delta r=(n-1)\left(f^{\prime} / f\right) \circ r .
$$

The statement $f(0)=0$ and $f^{\prime}(0)=1$ are to read as $f(0+)=0$ and $f^{\prime}(0+)=1$.
Proof. (i) At first let $v, w \in T_{0} \mathbb{M}$ such that $|v|=|w|=1$ and $v \perp w$. We consider the Jacobi field $J$ along the geodesic $\gamma$ where $\gamma(t)=t v$ such that $J(0)=0$ and $J^{\prime}(0)=w$. By Theorem 2.1.49 $J$ is up to a scalar function a parallel vector field along $\gamma$, hence $J(t)=c(t) W(t)$ with $\nabla_{D} W=0$ and without restriction $|W| \equiv 1$ as well as $c(t)>0$
for $t>0$; on the other hand $J(t)=\left(d \exp _{0}\right)_{t v}(t w) \equiv t w \in T_{t v} \mathbb{M}$. Thus $c(t)^{2}=$ $|c(t) W(t)|^{2}=|t w|^{2}=f(t)^{2}$, hence $J(t)=f(t) W(t)$ and in particular $f(0)=0$. In addition, it holds $\left(\nabla_{D} J\right)(t)=f^{\prime}(t) W(t)$; because of $\left(\nabla_{D} J\right)(0)=w=W(0)$ (see the proof of Theorem 2.1.49), hence $f^{\prime}(0)=1$.

Now $J=f W$ is a Jacobi field along $\gamma$, and by (2.1.16) hence

$$
\left\langle\nabla_{D} \nabla_{D} J, J\right\rangle=-\langle R(J, \dot{\gamma}) \dot{\gamma}, J\rangle
$$

This means $\left\langle f^{\prime \prime} W, f W\right\rangle=-\langle R(f W, \dot{\gamma}) \dot{\gamma}, f W\rangle=-f^{2}\langle R(W, \dot{\gamma}) \dot{\gamma}, W\rangle$, from where by

$$
\langle R(W, \dot{\gamma}) \dot{\gamma}, W\rangle(t)=k(t)|W(t)|^{2}=k(t)
$$

the relation $f^{\prime \prime}(t) f(t)=-f^{2}(t) k(t)$ (or equivalently $f^{\prime \prime}(t)+k(t) f(t)=0$ ) follows.
(ii) According to (2.1.32) we have $\nabla d r\left(\partial^{\mathbb{M}}, X\right)=0$ for each vector field $X$ on $\mathbb{M} \backslash\{0\}$, so that in particular $\nabla d r\left(\partial^{\mathbb{M}}, \partial^{\mathbb{M}}\right)=0$. It is hence sufficient to show that

$$
(\nabla d r)_{x}(u, u)=\left(f^{\prime} / f\right)(r(x)) \quad \text { for } u \in T_{x} \mathbb{M}, x \neq 0,|u|=1 \text { and } u \perp \partial_{x}^{\mathbb{M}}
$$

Let $\gamma:[0, b] \rightarrow \mathbb{M}$ be the normal geodesic with $\gamma(0)=0$ and $\gamma(b)=x$, and $J$ the unique Jacobi field along $\gamma$ with $J(0)=0$ and $J(b)=u$. Then, according to (2.1.28) along with the Jacobi equation for $J$, it holds that

$$
\begin{aligned}
(\nabla d r)_{x}(u, u) & =\int_{0}^{b}\left\{\left|\nabla_{D} J\right|^{2}-\langle R(J, \dot{\gamma}) \dot{\gamma}, J\rangle\right\} d t \\
& =\int_{0}^{b}\left\{\left|\nabla_{D} J\right|^{2}+\left\langle\nabla_{D} \nabla_{D} J, J\right\rangle\right\} d t=\left\langle\nabla_{D} J, J\right\rangle(b)=\left\langle\left(\nabla_{D} J\right)(b), u\right\rangle
\end{aligned}
$$

On the other hand, by Theorem 2.1.49 and part (i), we have $J=f W$ with $W$ a parallel vector field along $\gamma$; in particular then $\nabla_{D} J=f^{\prime} W$ and $f(b) W(b)=u$. This shows

$$
(\nabla d r)_{x}(u, u)=\left\langle\left(\nabla_{D} J\right)(b), u\right\rangle=\left\langle f^{\prime}(b) W(b), u\right\rangle=\left\langle f^{\prime}(b) / f(b) u, u\right\rangle=f^{\prime}(b) / f(b)
$$

and hence the claim.
On the other hand, Theorem 2.1.57 (i) opens a simple strategy for the construction of models: Starting with a differentiable function $k:[0, \infty[\rightarrow \mathbb{R}$, one determines $f$ as solution to the equation

$$
\begin{equation*}
f^{\prime \prime}(t)=-k(t) f(t), \quad f(0)=0, f^{\prime}(0)=1 \tag{2.1.33}
\end{equation*}
$$

If then $f>0$ on $] 0, \infty\left[\right.$, then $d r \otimes d r+f(r)^{2} d \vartheta^{2}$ defines a Riemannian metric on $\mathbb{R}^{n} \backslash\{0\}$, and one shows that because of $f(0)=0$ and $f^{\prime}(0)=1$ this metric allows a differentiable continuation to $\mathbb{R}^{n}$, in other words, there exists a Riemannian metric $g$ on $\mathbb{R}^{n}$ which restricted to $\mathbb{R}^{n} \backslash\{0\}$ coincides with $d r \otimes d r+f(r)^{2} d \vartheta^{2}$ (see [14] p. 60). Obviously ( $\mathbb{R}^{n}, g$ ) is then a model with $k_{\mathbb{M}}=k$ as radial curvature function. The problem which functions $k: \mathbb{R}_{+} \rightarrow \mathbb{R}$ can serve as radial curvature function of a model thus reduces to the question whether the corresponding solution $f$ to (2.1.33) stays positive on all of $] 0, \infty[$.

LEMMA 2.1.58. Let $n \geq 2$ and $k:\left[0, \infty\left[\rightarrow \mathbb{R}\right.\right.$ be a $C^{\infty}$-function such that either
(a) $k \leq 0$, or
(b) $k \geq 0$ and $\int_{0}^{\infty} s k(s) d s \leq 1$.

Then, up to isometry, there exists a unique model $\left(\mathbb{R}^{n}, g\right)$ with radial curvature function $k$.

Proof. We show that in both cases $f^{\prime}>0$ on $[0, \infty[$ must hold; since $f(0)=0$ then also $f>0$ on $] 0, \infty\left[\right.$ holds. Assume that $r:=\inf \left\{t>0: f^{\prime}(t)=0\right\}$ is finite. Because of $f^{\prime}(0)=1$ one has $r>0$, and then $f>0$ and $f^{\prime}>0$ on $] 0, r[$.
(a) Assume that $k \leq 0$. Then $f^{\prime \prime} \geq 0$ on $] 0, r\left[\right.$ and then $f^{\prime}(r)-f^{\prime}(0)=\int_{0}^{r} f^{\prime \prime}(s) d s \geq$ 0 , in contradiction to the definition of $r$.
(b) Assume now that $k \geq 0$. Then $f^{\prime \prime} \leq 0$ on $] 0, r\left[\right.$ and hence $f^{\prime} \leq 1$ on $] 0, r[$. This implies $f(s) \leq s$ for $s \in[0, \bar{r}]$ and hence $\int_{0}^{r} f(s) k(s) d s \leq \int_{0}^{r} s k(s) \bar{d} s$ with equality if and only if $f(s)=s$ for each $s \in[0, r]$, which however would imply $f^{\prime}(r)=1$ and is excluded by the definition of $r$. But then we have

$$
-1=f^{\prime}(r)-f^{\prime}(0)=\int_{0}^{r} f^{\prime \prime}(s) d s=-\int_{0}^{r} f(s) k(s) d s>-\int_{0}^{r} s k(s) d s \geq-1
$$

which is a contradiction.
The case of constant radial curvature is of particular interest. Let $c>0$ be a constant and suppose that $f^{\prime \prime}(t)=-k(t) f(t)$ with $f(0)=0$ and $f^{\prime}(0)=1$. Then:
(i) $f(t)=t$ for $t \in[0, \infty[$ if $k \equiv 0$ on $[0, \infty[$.
(ii) $f(t)=(1 / c) \sin c t$ for $t \in[0, r]$ if $k \equiv c^{2}$ on [0,r] with $r<\pi / c$.
(iii) $f(t)=(1 / c) \sinh c t$ for $t \in\left[0, \infty\left[\right.\right.$ if $k \equiv-c^{2}$ on $[0, \infty[$.

We want to investigate the different cases and to give descriptions of the spaces of constant curvature.
A. (Euclidean space) Let $(\mathbb{M}, g)=\left(\mathbb{R}^{n}\right.$, eucl) be the Euclidean space $\mathbb{R}^{n}$ with the standard metric. Obviously $\left(\mathbb{R}^{n}\right.$, eucl) is a model: it holds

$$
\left.\exp _{0} \mid\left(\mathbb{R}^{n} \backslash\{0\}\right):\right] 0, \infty\left[\times \mathbb{S}^{n-1} \rightarrow \mathbb{R}^{n}, \quad(t, v) \mapsto t v\right.
$$

and hence $\exp _{0}^{*} g=d r \otimes d r+r^{2} d \vartheta^{2}$. This corresponds to case (i) and gives up to isometry the unique model with radial curvature $k \equiv 0$; in addition the sectional curvature of $\left(\mathbb{R}^{n}\right.$, eucl) vanishes as well.
B. (Sphere) Let $\mathbb{S}_{a}^{n}=\left\{x \in \mathbb{R}^{n+1}:|x|=a\right\} \subset \mathbb{R}^{n+1}$ be the sphere of radius $a>0$, equipped with Riemannian metric $g$ induced from $\mathbb{R}^{n+1}$. Geodesics stay on great circles; fixing an arbitrary point on $\mathbb{S}_{a}^{n}$, for simplicity the north pole $n$, and $v \in T_{n} \mathbb{S}_{a}^{n} \cong\{n\}^{\perp}$ with $|v|=1$, we have

$$
\exp _{n}(t v)=\cos (t / a) n+a \sin (t / a) v
$$

Therefore it holds $\exp _{n}^{*} g=d r \otimes d r+a^{2} \sin ^{2}(r / a) d \vartheta^{2}$ on $] 0, a \pi\left[\times \mathbb{S}^{n-1}\right.$. Recall that $\left(\mathbb{S}_{a}^{n}, g\right)$ has constant sectional curvature, as already deduced from symmetry arguments. The representation of the metric in polar coordinates locally about the north pole as $d r \otimes$ $d r+f(r)^{2} d \vartheta^{2}$ with $f(r)=a \sin (r / a)$ gives for radial planes $E \subset T_{x} \mathbb{S}_{a}^{n}$ as value of the sectional curvature

$$
k(r)=-f^{\prime \prime}(r) / f(r)=1 / a^{2}
$$

where $r=d(n, x)<a \pi$; hence the sectional curvature of $\left(\mathbb{S}_{a}^{n}, g\right)$ is constant and equal to $1 / a^{2}$. However there is no model with positive radial curvature.
C. (Hyperbolic space) For $a>0$ let $M=\mathbb{B}_{a}^{n}=\left\{x \in \mathbb{R}^{n}:|x|<a\right\}$ be the open unit ball endowed with the Riemannian metric $g$ given by

$$
g_{x}(u, v):=\frac{4\langle u, v\rangle}{\left(1-|x|^{2} / a^{2}\right)^{2}}, \quad u, v \in T_{x} M \cong \mathbb{R}^{n}
$$

The normal geodesics $\gamma$ emanating from 0 with $\dot{\gamma}(0)=v \in T_{0} M$ obviously take the form $\gamma(t)=\kappa(t) v$ with $\kappa$ a scalar function such that $\kappa(0)=0$ and $\dot{\kappa}(0)=1$; from $|\dot{\gamma}(t)|=1$
we conclude that $\dot{\kappa}(t)=1-\kappa(t)^{2} /\left(4 a^{2}\right)$ and hence $\kappa(t)=2 a \tanh (t / 2 a)$. This gives

$$
\begin{array}{cc}
\mathbb{R}^{n} \backslash\{0\} & T_{0} M \\
\| & \| \\
\left.\exp _{0} \mid\left(\mathbb{R}^{n} \backslash\{0\}\right):\right] 0, \infty\left[\times \mathbb{S}^{n-1} \longrightarrow\right] 0, \infty\left[\times S_{1}(0) \longrightarrow M\right. \\
(t, v) \longmapsto(2 t, v / 2) \longmapsto \kappa(2 t) v / 2=a \tanh (t / a) v
\end{array}
$$

and $\exp _{0}^{*} g=d r \otimes d r+f(r)^{2} d \vartheta^{2}$ with $f(r)=a \sinh (r / a)$. Hence $\left(\mathbb{B}_{a}^{n}, g\right)$ with the origin as distinguished point is the up to isometry unique model with radial curvature $k \equiv-1 / a^{2}$. From invariance properties of the metric $g$ one deduces that $\left(\mathbb{B}_{a}^{n}, g\right)$ has constant sectional curvature $-1 / a^{2}$. One calls $\left(\mathbb{B}_{a}^{n}, g\right)$ the $n$-dimensional hyperbolic space with curvature $-1 / a^{2}$; in the case of constant negative curvature -1 one calls it simply the ( $n$-dimensional) hyperbolic space and writes $\mathbb{B}^{n}$ instead of $\mathbb{B}_{1}^{n}$.

There are other classical realizations of the hyperbolic space; $\left(\mathbb{B}_{a}^{n}, g\right)$ is usually called the ball model of hyperbolic geometry. We sketch two equivalent models, where we restrict ourselves to the case of curvature -1 :
(a) Let $\left(\mathbb{H}^{n}, h\right)$ be the upper half space $\mathbb{H}^{n}=\left\{x \in \mathbb{R}^{n}: x^{n}>0\right\}$ with the metric

$$
h_{x}(u, v)=\langle u, v\rangle /\left(x^{n}\right)^{2}, \quad u, v \in T_{x} \mathbb{H}^{n} \cong \mathbb{R}^{n}
$$

$\left(\mathbb{H}^{n}, h\right)$ is called Poincaré model of the $n$-dimensional hyperbolic space.
(b) Let $\mathbb{R}^{n+1}=\mathbb{R} \times \mathbb{R}^{n}$ equipped with the "Lorentz metric" $\langle x \mid y\rangle=-x^{0} y^{0}+$ $\sum_{i=1}^{n} x^{i} y^{i}$ and

$$
\mathbb{L}^{n}:=\left\{x \in \mathbb{R}^{n+1}:\langle x \mid x\rangle=-1, x^{0}>0\right\} .
$$

As inverse image of a regular value under a differentiable function, $\mathbb{L}^{n}$ is an $n$-dimensional submanifold of $\mathbb{R}^{n+1}$ : the sheet determined by the positive sign of $x^{0}$ of the two-sheet hyperboloid $\left\{x \in \mathbb{R}^{n+1}:\langle x \mid x\rangle=-1\right\}$. For $x \in \mathbb{L}^{n}$ one identifies

$$
T_{x} \mathbb{L}^{n} \cong\left\{y \in \mathbb{R}^{n+1}:\langle x \mid y\rangle=0\right\}=:\{x\}^{\perp} .
$$

Because of $\langle x \mid x\rangle=-1$, for $x \in \mathbb{L}^{n}$ the restriction of $\langle\cdot \mid \cdot\rangle$ to $\{x\}^{\perp}$ is positive definite, namely
$\langle y \mid y\rangle=\sum_{i=1}^{n}\left(y^{i}\right)^{2}-\left(\frac{1}{x^{0}} \sum_{i=1}^{n} x^{i} y^{i}\right)^{2} \geq \sum_{i=1}^{n}\left(y^{i}\right)^{2}-\frac{1}{\left(x^{0}\right)^{2}} \sum_{i=1}^{n}\left(x^{i}\right)^{2} \sum_{i=1}^{n}\left(y^{i}\right)^{2}=\sum_{i=1}^{n}\left(\frac{y^{i}}{x^{0}}\right)^{2} \geq 0$, and $\langle\cdot \mid \cdot\rangle$ defines canonically a Riemannian metric $k$ on $\mathbb{L}^{n}$. We call $\left(\mathbb{L}^{n}, k\right)$ hyperboloid model of the $n$-dimensional hyperbolic space.

THEOREM 2.1.59. $\left(\mathbb{B}^{n}, g\right),\left(\mathbb{H}^{n}, h\right)$ and $\left(\mathbb{L}^{n}, k\right)$ are isometric models of hyperbolic geometry.

Proof. The map $f(x):=\frac{\left(x^{1}, \ldots, x^{n}\right)}{1+x^{0}}$ defines an isometry $f: \mathbb{L}^{n} \rightarrow \mathbb{B}^{n}$, i.e., $f$ is a diffeomorphism and it holds:

$$
g_{f(x)}\left(d f_{x} u, d f_{x} v\right)=k_{x}(u, v), \quad u, v \in T_{x} \mathbb{L}^{n}, x \in \mathbb{L}^{n}
$$

Likewise, an isometry $\mathbb{H}^{n} \rightarrow \mathbb{B}^{n}$ is given by $\phi \circ \sigma: \mathbb{H}^{n} \rightarrow \mathbb{B}^{n}$ where $\sigma$ denotes the reflection at the plane $x^{n}=0$ and

$$
\phi(x):=e_{d}+\frac{2\left(x-e_{d}\right)}{\left|x-e_{d}\right|^{2}}
$$

where $e_{d}=(0, \ldots, 0,1) \in \mathbb{R}^{n}$. The verification of these properties is left to the reader.

### 2.2. Brownian Motion and Curvature

After these differential geometric preparations we continue again with probabilistic questions and start with the description of the distance process $d_{M}(o, X)$ of an $M$-valued Brownian motion $X$ to a given point $o \in M$. To this end, we refer to some elementary facts about one-dimensional diffusion processes, which are put together in Appendix A.1.

THEOREM 2.2.1. Let $(M, g)$ be a metrically complete Riemannian manifold with $n=$ $\operatorname{dim} M \geq 2$ and let $o \in M$ be a fixed point. Let $r(\cdot)=d_{M}(o, \cdot)$ denote the distance function to $o$ and $X$ be a Brownian motion on $(M, g)$ such that $X_{0}=x_{0} \neq o$ a.s. If

$$
\tau=\inf \left\{t \geq 0: X_{t} \in \operatorname{cut}(o)\right\}
$$

denotes the first hitting time of $X$ of the cut locus $\operatorname{cut}(o)$ of $M$ with respect to o, then there is a one-dimensional Brownian motion $\hat{W}$ (on some possibly enlarged filtered probability space) such that on $[0, \tau[$ it holds:

$$
\begin{equation*}
d(r(X))=d \hat{W}+\frac{1}{2}(\Delta r \circ X) d t \tag{2.2.1}
\end{equation*}
$$

Proof. Consider first the stopping time $\tau^{\prime}=\inf \left\{t \geq 0: X_{t} \in \operatorname{cut}(o) \cup\{o\}\right\}$. Since $r=d_{M}(o, \cdot)$ is differentiable on $M \backslash(\operatorname{cut}(o) \cup\{o\})$, the Geometric Itô formula can be applied and one obtains on $\left[0, \tau^{\prime}[\right.$

$$
d(r(X))=(d r)(X)(U d W)+\frac{1}{2}(\Delta r \circ X) d t .
$$

Letting $d \hat{W}:=(d r)(X)(U d W) \equiv \sum_{i}(d r)(X)\left(U e_{i}\right) d W^{i}$ with $\hat{W}_{0}=0$, we obtain
$d[\hat{W}, \hat{W}]=\sum_{i}\left((d r)(X) U e_{i}\right)^{2} d t=\sum_{i}\left\langle(\operatorname{grad} r) \circ X, U e_{i}\right\rangle^{2} d t=|(\operatorname{grad} r)(X)|^{2}=d t$.
Hence $\hat{W}$ defines a (stopped) Brownian motion which can be extended to all of $\mathbb{R}_{+}$with the usual methods.

It remains to show that $\tau^{\prime}=\tau$ a.s. To this end, we have to verify that $X$ does not hit the point $o$ a.s. We fix $\varepsilon>0$ with $B_{2 \varepsilon}(o) \cap \operatorname{cut}(o)=\varnothing$ such that $d_{M}\left(o, x_{0}\right)>\varepsilon$, and consider for $R:=d_{M}(o, X)$ inductively the following stopping times

$$
\begin{aligned}
& \sigma_{0}=\tau_{0}=0, \quad \text { and } \\
& \sigma_{n}=\inf \left\{t \geq \tau_{n-1}: R_{t}=\varepsilon\right\} \wedge \tau^{\prime}, \quad \tau_{n}=\inf \left\{t \geq \sigma_{n}: R_{t}=2 \varepsilon\right\} \wedge \tau^{\prime}, \quad n \geq 1
\end{aligned}
$$

It is obviously sufficient to show that the process $R \mid\left[\sigma_{n}, \tau_{n}[\right.$ does not hit 0 a.s. for any $n$. Without restrictions let $\sigma_{n}<\infty$ a.s. Now the Riemannian sectional curvature on $B_{2 \varepsilon}(o)$ is bounded, i.e. $\operatorname{Riem}^{M} \mid B_{2 \varepsilon}(o) \leq c^{2}$ for some $c>0$. After possibly diminishing $\varepsilon$ we may assume that $\varepsilon<\pi / 2 c$. Comparison with the sphere $S_{1 / c}^{n}$ combined with Theorems 2.1.55 and 2.1.57 (ii) gives

$$
(\Delta r)(x) \geq(n-1) c \cot c t, \quad t=d(x, o)<2 \varepsilon, \quad x \in B_{2 \varepsilon}(o) \subset M
$$

Using the abbreviations $\tilde{X}_{t}=X_{\left(\sigma_{n}+t\right) \wedge \tau_{n}}$ and $\tilde{R}_{t}=R_{\left(\sigma_{n}+t\right) \wedge \tau_{n}}$, we get in terms of the Brownian motion $\tilde{W}_{t}=\hat{W}_{\sigma_{n}+t}-\hat{W}_{\sigma_{n}}$ starting anew at $\sigma_{n}$ and the stopping time $\tilde{\tau}=\inf \left\{t \geq 0: \tilde{R}_{t}=2 \varepsilon\right.$ or $\left.\tilde{R}_{0}=0\right\}$ (both with respect to the transformed filtration) on the interval interval $[0, \tilde{\tau}[$ the equation

$$
\begin{equation*}
d \tilde{R}=d \tilde{W}+\frac{1}{2} \Delta r(\tilde{X}) d t, \quad \tilde{R}_{0}=\varepsilon \tag{2.2.2}
\end{equation*}
$$

We compare $\tilde{R}$ with the solution to the SDE

$$
\begin{equation*}
d Y=d \tilde{W}+\frac{n-1}{2} c \cot (c Y) d t, \quad Y_{0}=\varepsilon \tag{2.2.3}
\end{equation*}
$$

on the real interval $] 0,2 \varepsilon[$. At first we conclude from Theorem A.1.9 (ii) for the SDE (2.2.3) that 0 is a non-reachable boundary point of $Y$. It is hence sufficient to show that that $\tilde{R} \geq Y$ on $[0, \tilde{\tau}[$ which implies that also $\tilde{R}$ does not hit the point 0 a.s. To this end we can conclude as in Comparison Theorem A.1.8: If $\left.\left[a_{n}, b_{n}\right] \uparrow\right] 0,2 \varepsilon$ [ denotes a compact exhaustion with $a_{n}<\varepsilon<b_{n}$, then one has first for

$$
t \leq \inf \left\{s \geq 0: \tilde{R}_{s} \notin\left[a_{n}, b_{n}\right] \text { or } Y_{s} \notin\left[a_{n}, b_{n}\right]\right\}
$$

the pathwise inequalities

$$
\begin{aligned}
\left(Y_{t}-\tilde{R}_{t}\right)_{+} & =\int_{0}^{t} 1_{\left\{Y_{s}>\tilde{R}_{s}\right\}} d(Y-\tilde{R})_{s} \\
& =\int_{0}^{t} 1_{\left\{Y_{s}>\tilde{R}_{s}\right\}} \frac{1}{2}\left[(n-1) c \cot \left(c Y_{s}\right)-\left(\Delta r \circ \tilde{X}_{s}\right)\right] d s \\
& \leq \int_{0}^{t} 1_{\left\{Y_{s}>\tilde{R}_{s}\right\}} \frac{1}{2}(n-1) c\left[\cot \left(c Y_{s}\right)-\cot \left(c \tilde{R}_{s}\right)\right] d s \\
& \leq C_{n} \int_{0}^{t} 1_{\left\{Y_{s}>\tilde{R}_{s}\right\}}\left|Y_{s}-\tilde{R}_{s}\right| d s=C_{n} \int_{0}^{t}\left(Y_{s}-\tilde{R}_{s}\right)_{+} d s
\end{aligned}
$$

with a real constant $C_{n}$. By the Gronwall lemma it follows $\left(Y_{t}-\tilde{R}_{t}\right)_{+}=0$, hence $Y_{t} \leq \tilde{R}_{t}$, and then the claim as $n \rightarrow \infty$.

Theorem 2.2.1 indicates the general procedure: the distance process $r(X)=d_{M}(o, X)$ of an $M$-valued Brownian motion $X$ to a fixed reference point $o$ is (at least up to the first entrance in the cut locus $\operatorname{cut}(o)$ of $M$ with respect to $o$ ) of the form

$$
\begin{equation*}
r\left(X_{t}\right)=r\left(X_{0}\right)+\hat{W}_{t}+\frac{1}{2} \int_{0}^{t}\left(\Delta r \circ X_{s}\right) d s \tag{2.2.4}
\end{equation*}
$$

with a one-dimensional Brownian motion $\hat{W}$, where the drift part in (2.2.4) is controlled by curvature bounds according to Theorem 2.1.55.

THEOREM 2.2.2 (Comparison Theorem for Brownian motion). Let $(M, g)$ be a metrically complete Riemannian manifold of dimension $n=\operatorname{dim} M \geq 2$ and let $B_{\rho}(o)$ be an open geodesic ball of radius $\rho>0$ about a fixed point $o \in M$ which does not intersect the cut locus $\operatorname{cut}(o)$ of $M$ with respect to o. To this, suppose that there is a model $\mathbb{M}$ of the same dimension with center 0 and radial curvature function $k_{\mathbb{M}}$ such that for any $x \in M \backslash\{o\}$ with $0<d_{M}(o, x)=r<\rho$ it holds:

$$
\operatorname{Ric}_{x}^{M}\left(\partial^{M}, \partial^{M}\right) \geq(n-1) k_{\mathbb{M}}(r)
$$

respectively,

$$
\left.\operatorname{Riem}_{x}^{M}(E) \leq k_{\mathbb{M}}(r) \quad \text { for any radial plane } E \text { in } T_{x} M\right]
$$

Let $X$ be a Brownian motion on $(M, g)$, starting from a point $x_{0} \in B_{\rho}(o)$, and $\tau_{\rho}$ its exit time from $B_{\rho}(o)$. Correspondingly let $\tilde{X}$ be a Brownian motion on $\mathbb{M}$, starting from
$\tilde{x}_{0} \in \mathbb{M}$ with $d_{\mathbb{M}}\left(0, \tilde{x}_{0}\right)=d_{M}\left(o, x_{0}\right)$, and $\tilde{\tau}_{\rho}$ the exit time of $\tilde{X}$ from the open geodesic $\rho$-ball about 0 . Then for any antitone function $\varphi:[0, \rho[\rightarrow \mathbb{R}$,

$$
\begin{equation*}
E\left[\left(\varphi \circ d_{M}\left(o, X_{t}\right)\right) 1_{\left\{t<\tau_{\rho}\right\}}\right] \underset{[\leq]}{\geq} E\left[\left(\varphi \circ d_{\mathbb{M}}\left(0, \tilde{X}_{t}\right)\right) 1_{\left\{t<\tilde{\tau}_{\rho}\right\}}\right] \tag{2.2.5}
\end{equation*}
$$

In particular, for $0<\rho^{\prime}<\rho$ the following inequalities hold:

$$
\mathbb{P}\left\{d_{M}\left(o, X_{t}\right)<\rho^{\prime} \text { und } t<\tau_{\rho}\right\} \underset{[\leq]}{\geq} \mathbb{P}\left\{d_{\mathbb{M}}\left(0, \tilde{X}_{t}\right)<\rho^{\prime} \text { und } t<\tilde{\tau}_{\rho}\right\}
$$

Proof. Denote by $r_{M}(\cdot)=d_{M}(o, \cdot)$ and $r_{\mathbb{M}}(\cdot)=d_{\mathbb{M}}(0, \cdot)$ the distance processes to the distinguished points $o \in M, 0 \in \mathbb{M}$ and let $r_{0}:=r_{M}\left(x_{0}\right)=r_{\mathbb{M}}\left(\tilde{x}_{0}\right)$. Then, for $t<\tau_{\rho}$, respectively $t<\tilde{\tau}_{\rho}$,

$$
\begin{align*}
& r_{M}\left(X_{t}\right)=r_{0}+\hat{W}_{t}+\frac{1}{2} \int_{0}^{t} \Delta r_{M}\left(X_{s}\right) d s  \tag{2.2.6}\\
& r_{\mathbb{M}}\left(\tilde{X}_{t}\right)=r_{0}+\tilde{W}_{t}+\frac{1}{2} \int_{0}^{t} \Delta r_{\mathbb{M}}\left(\tilde{X}_{s}\right) d s \tag{2.2.7}
\end{align*}
$$

Since $\mathbb{M}$ is a model, we have $\Delta r_{\mathbb{M}}=(n-1)\left(f^{\prime} / f\right) \circ r_{\mathbb{M}}=: a \circ r_{\mathbb{M}}$ where $f$ denotes the radial function of the model. If the curvature of $M$ can be estimated from below in the way indicated, then as in the proof of Theorem 2.2.1, by Theorem 2.1.55 (Laplacian Comparison Theorem) the radial process $r_{M}(X)$ may be compared to the solution of the SDE

$$
\begin{equation*}
d Y=d \hat{W}+\frac{1}{2} a(Y) d t, \quad Y_{0}=r_{0} \tag{2.2.8}
\end{equation*}
$$

and one obtains $r_{M}\left(X_{t}\right) \leq Y_{t}$ for $t<\tau_{\rho}$ a.s. Hence if $\varphi$ is antitone, i.e. monotonically decreasing, then $\varphi \circ r_{M} \circ X_{t} \geq \varphi \circ Y_{t}$ for $t<\tau_{\rho}$ a.s. By the uniqueness in law for solutions of (2.2.8) we then have, as claimed,

$$
\mathbb{E}\left[\left(\varphi \circ d_{M}\left(o, X_{t}\right)\right) 1_{\left\{t<\tau_{\rho}\right\}}\right] \geq \mathbb{E}\left[\left(\varphi \circ d_{\mathbb{M}}\left(0, \tilde{X}_{t}\right)\right) 1_{\left\{t<\tilde{\tau}_{\rho}\right\}}\right]
$$

The case of upper curvature bounds for $M$ can be treated completely analogously by means of Corollary 2.1.46.

COROLLARY 2.2.3. Keeping the assumptions and notation of Theorem 2.2.2, we have in addition

$$
\mathbb{P}\left\{\tau_{\rho} \leq t\right\}_{[\geq]}^{\leq} \mathbb{P}\left\{\tilde{\tau}_{\rho} \leq t\right\}
$$

for any $t \geq 0$, and since $\mathbb{E}\left[\tau_{\rho}\right]=\int_{0}^{\infty} \mathbb{P}\left\{\tau_{\rho}>t\right\} d t$, then in particular

$$
\mathbb{E}\left[\tau_{\rho}\right] \geq \mathbb{E}\left[\tilde{\tau}_{\rho}\right]
$$

EXAMPLE 2.2.4. Let $(M, g)$ be a simply connected, metrically complete Riemannian manifold with $n=\operatorname{dim} M \geq 2$ and $o \in M$ be a fixed point. Suppose that $\operatorname{Riem}^{M} \leq 0$. If then $X$ is a Brownian motion on $(M, g)$ with $X_{0}=x_{0} \in M$, we have

$$
\mathbb{P}\left\{d_{M}\left(o, X_{t}\right)<\rho\right\} \leq \mathbb{P}\left\{R_{t}<\rho\right\}
$$

for any $\rho>0$ and $t>0$ where $R$ denotes a weak solution of the SDE

$$
d R=d W+(n-1) /(2 R) d t, \quad R_{0}=d_{M}\left(o, x_{0}\right)
$$

with $W$ representing a one-dimensional Brownian motion.

Proof. By Theorem 2.1.41 (Cartan-Hadamard), the cut locus of $M$ with respect to any point is empty, and the claim follows from Theorem 2.2.2 by comparison with ( $\mathbb{R}^{n}$, eucl).

Before continuing the discussion of the radial part of $M$-valued Brownian motions, we want to note some general properties of Brownian motions on Riemannian manifolds.

REMARK 2.2.5 (Strong Markov property of Brownian motion). Let $(M, g)$ be a Riemannian manifold. For $x \in M$, let $X^{x}$ denote a Brownian motion on $(M, g)$, starting at $x$, which we extend to a continuous process defined on $\mathbb{R}_{+}$and taking values in the one-point-compactification $\hat{M}$ of $M$. If then $H: C\left(\mathbb{R}_{+} ; \hat{M}\right) \rightarrow \mathbb{R}_{+}$is a bounded measurable function, then for any Brownian motion $X$ on $(M, g)$ and each stopping time $\tau$, it holds

$$
\begin{equation*}
\mathbb{E}^{\mathscr{F} \tau}\left[H\left(X_{\tau+\bullet}\right)\right]=\left.\mathbb{E}\left[H\left(X_{\cdot}^{y}\right)\right]\right|_{y=X_{\tau}} \quad \text { a.s. on }\{\tau<\infty\} . \tag{2.2.9}
\end{equation*}
$$

Proof. Taking into account the specific construction of Brownian motions as solutions of SDEs on the orthonormal frame bundle, the claim reduces to the strong Markov property of maximal solutions of SDEs with locally Lipschitz-continuous coefficients.

In general, the cut locus cut $(x)$ on a metrically complete Riemannian manifold $(M, g)$ with respect to a point $x$ is not a polar set; Brownian motions may hit the cut locus with positive probability, as can be seen from simple examples. However, for almost all paths of an $M$-valued Brownian motion, the occupation time on the cut locus equals zero, which comes from the fact that the cut locus is a nullset of the canonical Riemannian volume measure. We want briefly discuss this point.

DEFINITION 2.2.6 (Riemannian volume measure). On a Riemannian manifold ( $M, g$ ) there is exactly one measure vol on the Borel $\sigma$-algebra $\mathscr{B}(M)$ with the property that for each measurable function $f: M \rightarrow \mathbb{R}_{+}$with $\operatorname{supp}(f)$ in the domain of a chart $(\varphi, U)$ for $M$, it holds that

$$
\begin{equation*}
\int f d \mathrm{vol}=\int_{\varphi(U)}\left(f \sqrt{g^{(\varphi)}}\right) \circ \varphi^{-1} d x \tag{2.2.10}
\end{equation*}
$$

where $g^{(\varphi)}=\operatorname{det} G^{(\varphi)}>0$ with $G_{i j}^{(\varphi)}=g\left(\partial_{i}, \partial_{j}\right) \in C^{\infty}(U)$ and $\partial_{i}=\partial / \partial \varphi^{i}$. The measure vol is called Riemannian volume measure on $(M, g)$.

Remark 2.2.7. Note that if $(\psi, V)$ is another chart, then

$$
\sqrt{g^{(\varphi)}}=\sqrt{g^{(\psi)}}\left|\operatorname{det} J\left(\psi \circ \varphi^{-1}\right)\right| \circ \varphi \quad \text { on } U \cap V .
$$

On the other hand, if $\phi: D_{1} \rightarrow D_{2}$ is a diffeomorphism between two domains of $\mathbb{R}^{n}$, then by the transformation formula, for any non-negative measurable function $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$,

$$
\int_{D_{1}}(f \circ \phi)|\operatorname{det} J(\phi)| d x=\int_{D_{2}} f d x
$$

with $J(\phi)$ the Jacobian of $\phi$. Both observations together show that (2.2.10) is independent of the choice of the chart. Indeed, through (2.2.10), vol is first well-defined on Borel sets contained in the domain of a chart, and then also on all of $\mathscr{B}(M)$. The Riemannian volume measure on $\left(\mathbb{R}^{n}\right.$, eucl $)$ is obviously the $n$-dimensional Lebesgue measure.

On a complete Riemannian manifold $(M, g)$ the Laplacian generates a canonical semigroup of operators on the space $B(M)$ of bounded measurable functions on $M$ in the sense of a family of linear operators

$$
\begin{equation*}
P_{t}: B(M) \rightarrow B(M), \quad t \geq 0 \tag{2.2.11}
\end{equation*}
$$

with the properties:
(a) $P_{s} P_{t} f=P_{s+t} f$ for $f \in B(M)$.
(b) $P_{t} f \geq 0$ for $0 \leq f \in B(M)$, as well as $P_{t} 1 \leq 1$.
(c) $\left(P_{t} f\right)(x)-f(x)=\frac{1}{2} \int_{0}^{t}\left(P_{s} \Delta f\right)(x) d s$ for any test function $f \in C_{c}^{\infty}(M)$.
(d) $\left(P_{t}\right)_{t \geq 0}$ is minimal, i.e., for any other family $\left(Q_{t}\right)_{t \geq 0}$ of positive linear operators on $B(M)$ satisfying (a), (b), (c), it holds

$$
P_{t} f \leq Q_{t} f, \quad 0 \leq f \in B(M), t \geq 0
$$

In addition, $\left(P_{t}\right)_{t \geq 0}$ possesses a $C^{\infty}$-kernel $p \in C^{\infty}(] 0, \infty[\times M \times M)$ such that

$$
\begin{equation*}
\left(P_{t} f\right)(x)=\int p(t, x, y) f(y) \operatorname{vol}(d y), \quad f \in B(M), t>0 \tag{2.2.12}
\end{equation*}
$$

and $u(t, x):=\left(P_{t} f\right)(x)$ defines a classical solution of the heat equation

$$
\left\{\begin{array}{l}
\frac{\partial}{\partial t} u-\frac{1}{2} \Delta u=0  \tag{2.2.13}\\
\left.u\right|_{t=0}=f .
\end{array}\right.
$$

These are well-known facts from Spectral Theory of the heat kernel (see for instance [4], p. 187 ff .). We want briefly sketch the relation to Brownian motion.

THEOREM 2.2.8. Let $(M, g)$ be a metrically complete Riemannian manifold and $\left(P_{t}\right)_{t \geq 0}$ the minimal semigroup (2.2.11) generated by $\frac{1}{2} \Delta$. Then

$$
\begin{equation*}
\left(P_{t} f\right)(x)=\mathbb{E}\left[f\left(X_{t}^{x}\right) 1_{\left\{t<\zeta^{x}\right\}}\right], \quad f \in B(M) \tag{2.2.14}
\end{equation*}
$$

where $X^{x}$ denotes a Brownian motion with lifetime $\zeta^{x}$, starting in $x$. In particular, $(M, g)$ is BM -complete if and only if $P_{t} 1=1$.

Proof. Let $\left(Q_{t} f\right)(x):=\mathbb{E}\left[f\left(X_{t}^{x}\right) 1_{\left\{t<\zeta^{x}\right\}}\right]$. We fix a non-negative function $f \in$ $B(M)$, as well as $t \geq 0$. Since $u(t, x):=\left(P_{t} f\right)(x)$ solves the heat equation (2.2.13), it follows from Itô's formula that

$$
\left(Y_{s}\right)_{0 \leq s<t \wedge \zeta^{x}}, \quad Y_{s}:=\left(P_{t-s} f\right)\left(X_{s}^{x}\right)
$$

defines a non-negative local martingale. Hence there exists a localizing sequence of stopping times $\left(\zeta_{n}^{x}\right)_{n \in \mathbb{N}}$ with $\zeta_{n}^{x} \uparrow \zeta^{x}$ such that

$$
\left(P_{t} f\right)(x)=Y_{0}=\mathbb{E}\left[Y_{t \wedge \zeta_{n}^{x}}\right] \geq \mathbb{E}\left[\liminf _{n \rightarrow \infty} Y_{t \wedge \zeta_{n}^{x}}\right] \geq \mathbb{E}\left[Y_{t} 1_{\left\{t<\zeta^{x}\right\}}\right]=\left(Q_{t} f\right)(x)
$$

Now also $\left(Q_{t}\right)_{t \geq 0}$ satisfies the conditions (a), (b), (c) from above, where for instance (a) follows from the strong Markov property (Remark 2.2.5). We then conclude from the minimality of $\left(P_{t}\right)_{t \geq 0}$ that $P_{t}=Q_{t}$.

COROLLARY 2.2.9. For a metrically complete Riemannian manifold $(M, g)$ are equivalent:
(i) Bounded solutions $u$ of the heat equation $\frac{\partial}{\partial t} u-\frac{1}{2} \Delta u=0$ are uniquely determined by the initial condition $u(0, \cdot)$.
(ii) $(M, g)$ is BM -complete.

Proof. (i) $\Rightarrow$ (ii): For $x \in M$ let $X^{x}$ be again a Brownian motion starting at $x$ with lifetime $\zeta^{x}$. Then $u(t, x):=\left(P_{t} 1\right)(x)=\mathbb{P}\left\{\zeta^{x}>t\right\}$ solves the heat equation to the initial condition $u(0, \cdot) \equiv 1$. By means of the unique solvability we have $\mathbb{P}\left\{\zeta^{x}>t\right\}=1$ for any $t \geq 0$ and hence $\mathbb{P}\left\{\zeta^{x}=\infty\right\}=1$.
(ii) $\Rightarrow$ (i): Conversely, let $u$ be a bounded solution of the heat equation with initial condition $f=u(0, \cdot)$. For fixed $t>0$ and $x \in M$ then

$$
\left(Y_{s}\right)_{0 \leq s<t}, \quad Y_{s}:=u\left(t-s, X_{s}^{x}\right)
$$

defines a bounded martingale; hence $u(t, x)=\mathbb{E}\left[Y_{0}\right]=\mathbb{E}\left[Y_{t}\right]=\mathbb{E}\left[f \circ X_{t}^{x}\right]$ which gives the claim.

THEOREM 2.2.10. On a metrically complete Riemannian manifold $(M, g)$ the cut locus $\operatorname{cut}(x)$ of $M$ with respect to any point $x$ is a nullset of the Riemannian volume measure.

Proof. For $x \in M$ we have $\operatorname{cut}(x)=\exp _{x}\left(C_{x}\right)$ (according to Definition 2.1.10) with

$$
C_{x}=\left\{s(v) v: v \in T_{x} M,|v|=1, s(v)<\infty\right\}
$$

and the strictly positive continuous function $s:\left\{v \in T_{x} M:|v|=1\right\} \rightarrow \overline{\mathbb{R}}_{+}$defined by $s(v)=\sup \left\{t \geq 0: d\left(x, \exp _{x}(t v)\right)=t\right\}$. Now $C_{x} \subset T_{x} M$ is a Lebesgue nullset, as graph in polar coordinates of a (continuous) function. Then also cut $(x) \subset M$, as image of the differentiable map $\exp _{x}$ defined on $T_{x} M$, is a nullset with respect to the Riemannian volume measure, which is an immediate consequence of the definition of the Riemannian volume measure and the fact that Lebesgue nullsets are preserved under differentiable transformations of $\mathbb{R}^{n}$.

COROLLARY 2.2.11. The occupation time of a Brownian motion on the cut locus $\operatorname{cut}(x)$ of a metrically complete Riemannian manifold $(M, g)$ with respect to any point $x \in M$ is zero, i.e., for each Brownian motion $X$ on $(M, g)$ with lifetime $\zeta$ it holds:

$$
\int_{0}^{\zeta} 1_{\left\{t: X_{t} \in \operatorname{cut}(x)\right\}} d t=0 \quad \text { a.s. }
$$

Proof. Let $X$ be a Brownian motion on $(M, g)$; by the Markov property 2.2.9 (with $\tau=0$ ) without restriction with deterministic starting point. Then

$$
\mathbb{E}\left[\int_{0}^{\zeta} 1_{\left\{t: X_{t} \in \operatorname{cut}(x)\right\}} d t\right]=\int_{0}^{\infty} \mathbb{E}\left[1_{\operatorname{cut}(x)}\left(X_{t}\right) 1_{\{t<\zeta\}}\right] d t=0
$$

because of (2.2.12) and (2.2.14) the last equality is a consequence of Theorem 2.2.10.
We want to investigate now for Brownian motions $X$ on a Riemannian manifold $(M, g)$ properties of the distance process $d_{M}(o, X)$ (with respect to a given point $o \in M$ ) beyond the first entrance time of the Brownian motion into the cut locus cut $(o)$ of $M$ with respect to $o$. The main difficulty hereby, namely the distance function $d_{M}(o, \cdot)$ being differentiable only on $M \backslash(\operatorname{cut}(o) \cup\{o\})$, with the consequence that it is not even clear whether $d_{M}(o, X)$ represents a globally defined semimartingale, can be approached in different ways. On one hand, it is well-known that estimates for $\Delta d_{M}(o, \cdot)$ on $M \backslash(\operatorname{cut}(o) \cup\{o\})$ extend globally to all of $M$ if interpreted in the distributional sense (see [47], p. 669-70). We follow in contrast the approach of W.S. Kendall [23] and use the observation above that in general Brownian motions may in fact hit the cut locus but "spend no time on it" (see Corollary 2.2.11).

THEOREM 2.2.12. Let $(M, g)$ be a metrically complete Riemannian manifold with $n=\operatorname{dim} M \geq 2$ and let $r=d_{M}(o, \cdot)$ be the distance function to a given point $o \in M$. Let $X$ be a Brownian motion on $(M, g)$ with starting point $x_{0} \in M$ and lifetime $\zeta$, as well as $U$ a horizontal lift of $X$ to $\mathrm{O}(T M)$ and $W$ the $\mathbb{R}^{n}$-valued BM given as anti-development of $X$ (with respect to the initial basis $U_{0}$ ). Then, for $t<\zeta$,

$$
\begin{equation*}
r\left(X_{t}\right)-r\left(x_{0}\right)=\sum_{i=1}^{n} \int_{0}^{t} d r(X)\left(U e_{i}\right) d W^{i}+\frac{1}{2} \int_{0}^{t} \Delta r(X) d t-L_{t}^{(o)} \tag{2.2.15}
\end{equation*}
$$

where $d r$ and $\Delta r$ are set zero on $\operatorname{Cut}(o) \cup\{o\}$; here $L^{(o)}$ is an adapted isotone process which increases only when $X$ hits the cut locus $\operatorname{cut}(o)$, i.e.,

$$
\begin{equation*}
\int_{0}^{\zeta} 1_{\left\{t: X_{t} \notin \operatorname{cut}(o)\right\}} d L_{t}^{(o)}=0 \quad \text { a.s. } \tag{2.2.16}
\end{equation*}
$$

Note that the convention $d r=0$ and $\Delta r=0$ at places where $r$ is not differentiable, is inessential by Corollary 2.2.11.

Theorem 2.2.12 generalizes the Geometric Itô formula (Theorem 1.6.45) for the radial part of a Brownian motion $X$ to its whole life interval including the hitting times of the cut locus. The necessary subtraction of a "correction term" in the form of an isotone process $L^{(o)}$ which grows only when $X$ hits the cut locus, can be interpreted as local time of the Brownian motion on the cut locus cut $(o)$, see [5] for a detailed analysis of the geometric and stochastic nature of $L^{(o)}$.

Proof of Theorem 2.2.12. (see [23]) (1) The lifetime of $X$ (considered as continuous process taking values in the one-point-compactification of $M$ ) is given by

$$
\zeta=\sup \{t>0: r(X) \text { is bounded on }[0, t]\}
$$

It is hence sufficient to verify (2.2.15) up to the first exit from a geodesic ball, that is, up to the first time $r(X)$ exceeds a certain value. The claim to verify is then only concerns a sufficiently large geodesic ball $B$. An elementary consideration thus shows that $(M, g)$ may be modified outside of $B$ to a compact Riemannian manifold. For simplicity we may hence assume without restriction of generality $M$ to be already compact; in particular then Riem ${ }^{M} \geq-c^{2}$ for some $c>0$, and the injectivity radius of $M$ being strictly positive, i.e., $\varrho=\inf \{d(x, \operatorname{cut}(x)): x \in B\}>0$.
(2) We verify first that $r(X)$ defines a semimartingale. To this end, we show that for a suitable function $V$ on $M$ the process

$$
\begin{equation*}
r\left(X_{t}\right)-r\left(x_{0}\right)-\int_{0}^{t} V\left(X_{s}\right) d s, \quad t \geq 0 \tag{2.2.17}
\end{equation*}
$$

is a supermartingale. This is sufficient since by the general Doob-Meyer decomposition (e.g., [28], section 3.7) each supermartingale is in particular a semimartingale. By (2.2.17) then trivially also $r(X)$ is a semimartingale and can be decomposed as

$$
\begin{equation*}
r(X)=r\left(x_{0}\right)+N+A, \quad N \in \mathscr{M}_{0}, A \in \mathscr{A}_{0} \tag{2.2.18}
\end{equation*}
$$

We consider

$$
V: M \backslash\{o\} \rightarrow \mathbb{R}, \quad V(x):= \begin{cases}\frac{n-1}{2} c \operatorname{coth} c r(x) & \text { for } r(x) \leq \varrho / 3, \\ \frac{n-1}{2} c \operatorname{coth} c \varrho / 3 & \text { for } r(x) \geq \varrho / 3 .\end{cases}
$$

Comparison with the $n$-dimensional hyperbolic space of constant curvature $-c^{2}$, i.e. the model with radial function $f(t)=(1 / c) \sinh c t$, gives by Theorem 2.1.55 (Laplacian Comparison Theorem) and Theorem 2.1.57,

$$
\begin{equation*}
\frac{1}{2}(\Delta r)(x) \leq V(x) \quad \text { for } x \notin \operatorname{cut}(o) \cup\{o\} \tag{2.2.19}
\end{equation*}
$$

As already verified in the proof of Theorem 2.2.1, the Brownian motion $X=\left(X_{t}\right)_{t \geq 0}$ does not hit the given point $o$ for $t>0$ a.s. For arbitrary $0<t_{1} \leq t_{2}$, we have to show that

$$
\begin{equation*}
\mathbb{E}^{\mathscr{F}_{t_{1}}}\left[r\left(X_{t_{2}}\right)-r\left(X_{t_{1}}\right)-\int_{t_{1}}^{t_{2}} V\left(X_{s}\right) d s\right] \leq 0 \tag{2.2.20}
\end{equation*}
$$

By the strong Markov property of Brownian motion (Remark 2.2.5) it is then sufficient to show that for each Brownian motion $X$ on $(M, g)$ with deterministic starting point (different to $o$ ) in $M$

$$
\begin{equation*}
\mathbb{E}\left[r\left(X_{t}\right)-r\left(X_{0}\right)-\int_{0}^{t} V\left(X_{s}\right) d s\right] \leq 0, \quad t \geq 0 \tag{2.2.21}
\end{equation*}
$$

We divide the proof into several steps.
(3) Let $x_{0} \in \operatorname{cut}(o)$ be an arbitrary point, and

$$
\gamma_{v}(t)=\exp _{o}(t v), \quad 0 \leq t \leq s(v), v \in T_{o} M,|v|=1
$$

be a minimal geodesic from $o$ to $x_{0}$. Then $\gamma_{v}(\varrho / 3) \notin \operatorname{cut}\left(x_{0}\right)$, or equivalently $x_{0} \notin$ $\operatorname{cut}\left(\gamma_{v}(\varrho / 3)\right)$. Hence the following two subsets of $M \times M$,

$$
\begin{aligned}
& \text { cut }:=\{(x, y) \in M \times M: y \in \operatorname{cut}(x)\} \\
& C:=\left\{\left(\gamma_{v}(s(v)), \gamma_{v}(\varrho / 3)\right): v \in T_{o} M,|v|=1\right\}
\end{aligned}
$$

are disjoint and have positive distance with respect to the product metric on $M \times M$. Hence there exists $\delta>0$ with the following property: If $x_{0}=\gamma_{v}(s(v)) \in \operatorname{cut}(o)$, then $x \notin \operatorname{cut}\left(\gamma_{v}(\varrho / 3)\right)$ for each $x \in M$ such that $d\left(x, x_{0}\right)<\delta$. We choose such a $\delta>0$ such that in addition $\delta<\varrho / 3$. This leads to the following

Claim: If $X$ is a Brownian motion with $X_{0}=x_{0} \in \operatorname{cut}(o)$ and $\tau:=\inf \{t \geq 0:$ $\left.d\left(X_{0}, X_{t}\right)=\delta\right\}$, then

$$
\begin{equation*}
\mathbb{E}\left[r\left(X_{t \wedge \tau}\right)-r\left(X_{0}\right)-\int_{0}^{t \wedge \tau} V\left(X_{s}\right) d s\right] \leq 0, \quad t \geq 0 \tag{2.2.22}
\end{equation*}
$$

Indeed, fixing a minimal geodesic $\gamma_{v}(t)=\exp _{o}(t v)$ from $o$ to $x_{0}$, then with $\hat{o}:=\gamma_{v}(\varrho / 3)$, according to the choice of $\delta$, die function

$$
\hat{r}(x):=d(x, \hat{o})
$$

is differentiable on the geodesic ball $B_{\delta}\left(x_{0}\right)$ about $x_{0}$ of radius $\delta$, and by Theorem 2.2.1 we have

$$
\begin{equation*}
\mathbb{E}\left[\hat{r}\left(X_{t \wedge \tau}\right)-\hat{r}\left(X_{0}\right)-\frac{1}{2} \int_{0}^{t \wedge \tau}(\Delta \hat{r})\left(X_{s}\right) d s\right]=0, \quad t \geq 0 \tag{2.2.23}
\end{equation*}
$$

The same comparison argument leading to (2.2.19) now gives

$$
\begin{equation*}
\frac{1}{2}(\Delta \hat{r})(x) \leq \hat{V}(x) \quad \text { for } x \notin \operatorname{cut}(\hat{o}) \cup\{\hat{o}\} \tag{2.2.24}
\end{equation*}
$$

with the modified function

$$
\hat{V}: M \backslash\{\hat{o}\} \rightarrow \mathbb{R}, \quad \hat{V}(x):= \begin{cases}\frac{n-1}{2} c \operatorname{coth} c \hat{r}(x) & \text { for } \hat{r}(x) \leq \varrho / 3 \\ \frac{n-1}{2} c \operatorname{coth} c \varrho / 3 & \text { for } \hat{r}(x) \geq \varrho / 3\end{cases}
$$

where for $x \in B_{\delta}\left(x_{0}\right)$ by definition $\hat{V}(x)=V(x)$ holds according to $\delta<\varrho / 3$. Considering finally the function

$$
r_{+}(x):=\hat{r}(x)+\varrho / 3
$$

we observe that $r_{+}(x) \geq r(x)$ by the triangle inequality where $r_{+}\left(x_{0}\right)=r\left(x_{0}\right)$. Hence it holds

$$
\frac{1}{2}\left(\Delta r_{+}\right)(x)=\frac{1}{2}(\Delta \hat{r})(x) \leq \hat{V}(x)=V(x), \quad x \in B_{\delta}\left(x_{0}\right)
$$

and (2.2.22) follows:

$$
\begin{aligned}
& \mathbb{E}\left[r\left(X_{t \wedge \tau}\right)-r\left(X_{0}\right)-\int_{0}^{t \wedge \tau} V\left(X_{s}\right) d s\right] \\
& \quad \leq \mathbb{E}\left[r_{+}\left(X_{t \wedge \tau}\right)-r_{+}\left(X_{0}\right)-\int_{0}^{t \wedge \tau} V\left(X_{s}\right) d s\right] \\
& \quad=\mathbb{E}\left[\hat{r}\left(X_{t \wedge \tau}\right)-\hat{r}\left(X_{0}\right)-\int_{0}^{t \wedge \tau} \hat{V}\left(X_{s}\right) d s\right] \leq 0
\end{aligned}
$$

the last inequality holds by (2.2.23) und (2.2.24). This completes the proof of the Claim.
(4) Assertion (2.2.21) can now be verified by means of the Claim in part (3): For each Brownian motion $X$ on $(M, g)$ with $X_{0}=x_{0} \neq o$, it holds

$$
\begin{equation*}
\mathbb{E}\left[r\left(X_{t}\right)-r\left(X_{0}\right)-\int_{0}^{t} V\left(X_{s}\right) d s\right] \leq 0, \quad t \geq 0 \tag{2.2.25}
\end{equation*}
$$

Indeed, defining inductively sequences of stopping times $\left(\tau_{n}\right)_{n \geq 0}$ and $\left(\sigma_{n}\right)_{n \geq 1}$ by $\tau_{0}=0$ and

$$
\begin{align*}
\sigma_{n} & =\inf \left\{t \geq \tau_{n-1}: X_{t} \in \operatorname{cut}(o)\right\} \\
\tau_{n} & =\inf \left\{t \geq \sigma_{n}: d\left(X_{t}, X_{\sigma_{n}}\right)=\delta\right\}, \quad n \geq 1 \tag{2.2.26}
\end{align*}
$$

one obtains by using the strong Markov property (Remark (2.2.5)), for any $n \in \mathbb{N}$,

$$
\begin{align*}
\mathbb{E}^{\mathscr{F} \tau_{n-1}} & {\left[r\left(X_{t \wedge \sigma_{n}}\right)-r\left(X_{t \wedge \tau_{n-1}}\right)-\int_{t \wedge \tau_{n-1}}^{t \wedge \sigma_{n}} V\left(X_{s}\right) d s\right] \leq 0 }  \tag{2.2.27}\\
& \mathbb{E}^{\mathscr{F}_{\sigma_{n}}}\left[r\left(X_{t \wedge \tau_{n}}\right)-r\left(X_{t \wedge \sigma_{n}}\right)-\int_{t \wedge \sigma_{n}}^{t \wedge \tau_{n}} V\left(X_{s}\right) d s\right] \leq 0 \tag{2.2.28}
\end{align*}
$$

Here (2.2.27) is a consequence of Theorem 2.2.1 (on the radial part of a Brownian motion) and estimate (2.2.19), whereas (2.2.28) reduces to the Claim by means of the strong Markov property. To complete the proof of (2.2.25) only $\tau_{n} \uparrow \infty$ a.s. needs to be verified.

To this end, consider to a fixed $\varepsilon>0$ the independent sequence of events

$$
A_{n}:=\left\{\tau_{n}-\sigma_{n} \geq \varepsilon\right\}, \quad n \in \mathbb{N}
$$

By the Lemma of Borel-Cantelli it is sufficient to show $\sum_{n=1}^{\infty} \mathbb{P}\left(A_{n}\right)=\infty$. We may compare with the $n$-dimensional hyperbolic space $\mathbb{H}^{n}\left(-c^{2}\right)$ of constant curvature $-c^{2}$ : If $B$ denotes the geodesic ball in $\mathbb{H}^{n}\left(-c^{2}\right)$ about 0 of Radius $\delta$, then we get by Theorem 2.2.2 (Comparison Theorem for Brownian motion), using again the strong Markov property of the Brownian motion $X$ on $M$,

$$
\begin{aligned}
& \mathbb{E}^{\mathscr{F}_{\sigma_{n}}}\left[1_{\left\{\tau_{n}-\sigma_{n} \geq \varepsilon\right\}}\right] \\
& \geq \mathbb{P}\left\{\text { exit time of } \mathrm{BM}\left(\mathbb{H}^{n}\left(-c^{2}\right)\right) \text { from } B \text { when starting in } 0 \text { is at least } \varepsilon\right\} \\
& \geq 1 / 2 \quad \text { for } \varepsilon>0 \text { sufficiently small. }
\end{aligned}
$$

This shows that $r(X)$ can be written as sum of a supermartingale and an isotone process; hence, in particular, $r(X)$ is a semimartingale. We want to continue by giving a more detailed description of the terms in (2.2.18).
(5) Adopting the convention $d r=0$ on $\operatorname{cut}(o) \cup\{o\}$, the process

$$
\sum_{i=1}^{n} \int_{0}^{t} d r(X)\left(U e_{i}\right) d W^{i} \equiv \int_{0}^{t} d r(X) U d W, \quad t \geq 0
$$

is seen to be the martingale part of $r(X)$.
Indeed, denoting by $r(X)=r\left(x_{0}\right)+N+A$ the decomposition of $r(X)$ as semimartingale, the martingale part $N$ allows an integral representation of the form

$$
N_{t}=\int_{0}^{t} F d W \equiv \int_{0}^{t} \tilde{F} U d W
$$

with a uniquely determined predictable $\mathbb{R}^{n}$-valued process $F$, respectively $\tilde{F}:=F U^{-1}$ the corresponding $T^{*} M$-valued process over $X$. Considering the difference

$$
\bar{N}_{t}:=\int_{0}^{t} \tilde{F} U d W-\int_{0}^{t} d r(X) U d W
$$

one observes that the local martingale $\bar{N}$ is constant on each stochastic interval, on which $X$ doesn't hit the cut locus cut $(o)$; since $X$ avoids the point $o$ almost surely, we have on such an interval by the geometric Itô formula

$$
d(r(X))=(d r)(X)(U d W)+\frac{1}{2}(\Delta r \circ X) d t
$$

For $\delta>0$ sufficiently small, we consider again the stopping times (2.2.26) and set

$$
\left.I_{\delta}:=\bigcup_{n \in \mathbb{N}}\right] \tau_{n-1}, \sigma_{n}\left[\uparrow I_{*} \quad \text { for } \delta \searrow 0\right.
$$

Obviously, it holds $I_{*}=\left\{(t, \omega): X_{t}(\omega) \notin \operatorname{cut}(o)\right\}$. As already noted, $[\bar{N}, \bar{N}]$ is constant on each $I_{\delta}$, and hence $\int_{I_{\delta}} d[\bar{N}, \bar{N}]=0$, which implies $\int_{I_{*}} d[\bar{N}, \bar{N}]=0$ almost surely. In addition also $\int_{0}^{\infty} 1_{\left\{X_{t} \in \operatorname{cut}(o)\right\}} d[\bar{N}, \bar{N}]$ almost surely, since $[\bar{N}, \bar{N}]$ is absolutely continuous with respect to the Lebesgue measure and since $\int_{0}^{\infty} 1_{\left\{t: X_{t} \in \mathrm{cut}(0)\right\}} d t=0$ holds almost surely by Corollary 2.2.11. Together it shows $[\bar{N}, \bar{N}]=0$ almost surely and hence $\bar{N}=0$ modulo indistinguishability. This gives

$$
N=\int d r(X) U d W
$$

as wanted.
(6) Following the convention $\Delta r=0$ on $\operatorname{cut}(o) \cup\{o\}$, the $L^{(o)}$,

$$
L_{t}^{(o)}=\int_{0}^{t} d r(X) U d W+\frac{1}{2} \int_{0}^{t} \Delta r(X) d s-\left(r\left(X_{t}\right)-r\left(x_{0}\right)\right)
$$

is an isotone process with the property $\int_{0}^{\infty} 1_{\left\{t: X_{t} \notin \operatorname{cut}(o)\right\}} d L_{t}^{(o)}=0$ almost surely.
Let $I_{\delta}$ be as in (5) with $\delta=1 / n, n \in \mathbb{N}$. For sufficiently large $n$,

$$
L_{t}^{(o, n)}:=\int_{0}^{t} d r(X) U d W+\frac{1}{2} \int_{[0, t] \cap I_{1 / n}} \Delta r(X) d s+\int_{[0, t] \backslash I_{1 / n}} V(X) d s-\left(r\left(X_{t}\right)-r\left(x_{0}\right)\right)
$$

determines an isotone process $L^{(o, n)}$. By (2.2.19) it holds $L_{t}^{(o, n)} \geq L_{t}^{(o, n+1)}$, and hence

$$
L_{t}^{(o, \infty)}:=\lim _{n \rightarrow \infty} L_{t}^{(o, n)}
$$

defines an isotone process $L^{(o, \infty)}$ such that

$$
\begin{aligned}
L_{t}^{(o, \infty)} & =\int_{0}^{t} d r(X) U d W+\frac{1}{2} \int_{[0, t] \cap I_{*}} \Delta r(X) d s+\int_{[0, t] \backslash I_{*}} V(X) d s-\left(r\left(X_{t}\right)-r\left(x_{0}\right)\right) \\
& =\int_{0}^{t} d r(X) U d W+\frac{1}{2} \int_{0}^{t} \Delta r(X) d s-\left(r\left(X_{t}\right)-r\left(x_{0}\right)\right)
\end{aligned}
$$

for the last equality we used $I_{*}=\left\{(t, \omega): X_{t}(\omega) \notin \operatorname{cut}(o)\right\}$ together with Corollary 2.2.10. This shows $L^{(o, \infty)}=L^{(o)}$.

The still missing property $\int_{0}^{\infty} 1_{\left\{t: X_{t} \notin \operatorname{cut}(o)\right\}} d L_{t}^{(o)}=0$ a.s. comes from the equation $\int_{0}^{\infty} 1_{I_{\delta}}(t, \cdot) d L_{t}^{(o)}=0$ a.s. which holds for each $\delta>0$.

Theorem 2.2.12 allows to sharpen the Comparison Theorem for Brownian motion (Theorem 2.2.2) in the case of lower curvature bounds: in this case one may consider arbitrary geodesic balls, also balls which intersect the cut locus.

THEOREM 2.2.13 (Comparison Theorem for Brownian motion; strong version). Let $(M, g)$ be a metrically complete Riemannian manifold of dimension $n \geq 2$ and let $B_{\rho}(o)$ be the open geodesic ball of radius $\rho>0$ about some given point $o \in M$. Suppose that there exists a model $\mathbb{M}$ of same dimension with center 0 and radial curvature function $k_{\mathbb{M}}$ such that for any $x \in M \backslash(\operatorname{cut}(o) \cup\{o\})$ with $0<d_{M}(o, x)=r<\rho$ it holds:

$$
\operatorname{Ric}_{x}^{M}\left(\partial^{M}, \partial^{M}\right) \geq(n-1) k_{\mathbb{M}}(r)
$$

Let $X$ be a Brownian motion on $(M, g)$, starting in a point $x_{0} \in B_{\rho}(o)$, and $\tau_{\rho}$ be its exit time from $B_{\rho}(o)$. Accordingly, let $\tilde{X}$ be a Brownian motion on $\mathbb{M}$, starting in $\tilde{x}_{0} \in \mathbb{M}$ with $d_{\mathbb{M}}\left(0, \tilde{x}_{0}\right)=d_{M}\left(o, x_{0}\right)$, and $\tilde{\tau}_{\rho}$ the exit time of $\tilde{X}$ from the open geodesic $\rho$-ball about 0 . Then, for any antitone function $\varphi:[0, \rho[\rightarrow \mathbb{R}$,

$$
\mathbb{E}\left[\left(\varphi \circ d_{M}\left(o, X_{t}\right)\right) 1_{\left\{t<\tau_{\rho}\right\}}\right] \geq \mathbb{E}\left[\left(\varphi \circ d_{\mathbb{M}}\left(0, \tilde{X}_{t}\right)\right) 1_{\left\{t<\tilde{\tau}_{\rho}\right\}}\right]
$$

In particular, for $0<\rho^{\prime}<\rho$, one has the inequalities:

$$
\mathbb{P}\left\{d_{M}\left(o, X_{t}\right)<\rho^{\prime} \text { and } t<\tau_{\rho}\right\} \geq \mathbb{P}\left\{d_{\mathbb{M}}\left(0, \tilde{X}_{t}\right)<\rho^{\prime} \text { and } t<\tilde{\tau}_{\rho}\right\}
$$

Proof. According to (2.2.15) we have for $t<\tau_{\rho}$,

$$
r_{M}\left(X_{t}\right) \leq r_{M}\left(x_{0}\right)+\hat{W}_{t}+\frac{1}{2} \int_{0}^{t} \Delta r_{M}(X) d t
$$

where $\hat{W}_{t}:=\sum_{i=1}^{n} \int_{0}^{t} d r_{M}(X)\left(U e_{i}\right) d W^{i}$ is a one-dimensional Brownian motion (stopped at $\zeta$ ). The remainder of the proof of Theorem 2.2.2 then carries over verbatim.

Before continuing the discussion on further asymptotic properties of Brownian motions, we want to note some fundamental facts about harmonic functions.

Lemma 2.2.14. Let $(M, g)$ be a Riemannian manifold and $h: M \rightarrow \mathbb{R}$ a bounded measurable function. The following conditions are equivalent:
(i) $h$ is harmonic (i.e., $h \in C^{\infty}(M)$ and $\Delta h=0$ ).
(ii) $h(x)=\mathbb{E}\left[h \circ X_{\tau}^{x}\right]$ for any $x \in M$ and any stopping time $\tau$ such that $0 \leq \tau<\zeta$ a.s.
(iii) $h$ has the mean-value property, i.e., for any $x_{0} \in M$ and any sufficiently small geodesic $\varepsilon$-ball $B_{\varepsilon}\left(x_{0}\right) \subset M$ about $x_{0}$,

$$
h(x)=\mathbb{E}\left[h \circ X_{\tau^{x}}^{x}\right], \quad x \in B_{\varepsilon}\left(x_{0}\right),
$$

where $\tau^{x}=\inf \left\{t \geq 0: X_{t}^{x} \notin B_{\varepsilon}\left(x_{0}\right)\right\}$ is the first exit time from $B_{\varepsilon}\left(x_{0}\right)$.
PROOF. (i) $\Rightarrow$ (ii) is a direct consequence of Itô's formula combined with the Optional Sampling Theorem.
(ii) $\Rightarrow$ (iii) is a weakening; the almost sure finiteness of the exit time of Brownian motions from small geodesic balls follows immediately from Theorem 2.2.1.
(iii) $\Rightarrow$ (i): We exploit the solvability of the Dirichlet problem for small geodesic balls in the following sense (see e.g. [2]): To each $x_{0} \in M$ and sufficiently small $\varepsilon>0$
there exists a family $\left(k_{\varepsilon}(x, d y)\right)_{x \in B_{\varepsilon}\left(x_{0}\right)}$ of "harmonic" measures $k_{\varepsilon}(x, d y)$ on $S_{\varepsilon}\left(x_{0}\right)=$ $\partial B_{\varepsilon}\left(x_{0}\right)$, namely $k_{\varepsilon}(x, d y)=\mathbb{P} \circ\left(X_{\tau^{x}}^{x}\right)^{-1}(d y)$, such that for each bounded measurable boundary function $f: S_{\varepsilon}\left(x_{0}\right) \rightarrow \mathbb{R}$ a harmonic function $\varphi_{f}$ is defined on $B_{\varepsilon}\left(x_{0}\right)$ by

$$
\varphi_{f}(x)=\int k_{\varepsilon}(x, d y) f(y) \equiv \mathbb{E}\left[f \circ X_{\tau^{x}}^{x}\right]
$$

For $f \in C\left(S_{\varepsilon}\left(x_{0}\right) ; \mathbb{R}\right)$ the function $\varphi_{f}$ is the unique harmonic continuation of $f$ to $B_{\varepsilon}\left(x_{0}\right)$. The support of $\mathbb{P} \circ\left(X_{\tau^{x}}^{x}\right)^{-1}$ is $S_{\varepsilon}\left(x_{0}\right)$, i.e., $\mathbb{P}\left\{X_{\tau^{x}}^{x} \in U\right\}>0$ for each non-empty open subset $U \subset S_{\varepsilon}\left(x_{0}\right)$.

Applied to our situation this means that $h=\varphi_{h}$ on $B_{\varepsilon}\left(x_{0}\right)$ for each $x_{0} \in M$ and each sufficiently small $\varepsilon>0$; in particular $h$ is harmonic.

Note that the equivalence (i) $\Leftrightarrow$ (iii) in Lemma 2.2.14 also holds for not necessarily bounded functions.

COROLLARY 2.2.15 (Maximum principle). Let $(M, g)$ be a Riemannian manifold, $h: M \rightarrow \mathbb{R}$ a harmonic function, $m=\sup _{x \in M} h(x) \in \overline{\mathbb{R}}$. If $h\left(x_{0}\right)=m$ for some $x_{0} \in M$, then $h$ is constant.

Proof. The set $M_{0}:=\{x \in M: h(x)=m\}$ is open in $M$ as a consequence of the mean value property; trivially, $M_{0}$ is closed by the continuity of $h$. Since all manifolds are assumed to be connected, the claim follows.

The next Theorem shows how on a Riemannian manifold $(M, g)$ asymptotic properties of $\mathrm{BM}(M, g)$ and richness of harmonic functions on $M$ correspond to each other.

Theorem 2.2.16. For a Riemannian manifold $(M, g)$ the following two items are equivalent:
(i) $\mathrm{BM}(M, g)$ has only trivial exit sets, i.e., if $X$ is a Brownian motion on $(M, g)$ starting from a deterministic initial point and $U \subset \hat{M}$ an open subset of the one-pointcompactification $\hat{M}$ of $M$, then

$$
\mathbb{P}\left\{X_{t} \in U \text { eventually }\right\} \in\{0,1\}
$$

(ii) $(M, g)$ is a Liouville manifold, i.e., all bounded harmonic functions on $M$ are constant.

For a Brownian motion $X$ on $(M, g)$ with lifetime $\zeta$, we use again the convention $X_{t}(\omega)=\infty$ in $\hat{M}$ for $t>\zeta(\omega)$. If $X_{0}=x \in M$, we write $X=X^{x}$ and denote by $\zeta^{x}$ the corresponding lifetime.

Proof of Theorem 2.2.16. For $U \subset \hat{M}$ let

$$
\mathscr{H}_{U}:=\left\{\alpha \in C\left(\mathbb{R}_{+} ; \hat{M}\right): \alpha(t) \in U \text { eventually }\right\}
$$

so that

$$
\begin{aligned}
X_{\bullet}^{-1}\left(\mathscr{H}_{U}\right) & =\left\{X_{t} \in U \text { eventually }\right\} \\
& =\left\{\omega \in \Omega: \exists t_{0}(\omega)>0 \text { such that } X_{t}(\omega) \in U \text { for all } t \geq t_{0}(\omega)\right\}
\end{aligned}
$$

We first note that for an open $U \subset \hat{M}$ the function $h_{U}: M \rightarrow \mathbb{R}$,

$$
h_{U}(x):=\mathbb{P}\left\{X_{t}^{x} \in U \text { eventually }\right\}=\mathbb{E}\left[1_{\mathscr{H}_{U}} \circ X_{.}^{x}\right]
$$

is harmonic. Indeed, by Remark 2.2.5 (strong Markov property of Brownian motion) it holds for each stopping time $\tau$ with $0 \leq \tau<\zeta^{x}$ that

$$
h_{U}\left(X_{\tau}^{x}\right)=\left.\mathbb{E}\left[1_{\mathscr{H}_{U}}\left(X_{\bullet}^{y}\right)\right]\right|_{y=X_{\tau}^{x}}=\mathbb{E}^{\mathscr{F}_{\tau}}\left[1_{\mathscr{H}_{U}}\left(X_{\tau+\cdot}^{x}\right)\right]=\mathbb{E}^{\mathscr{F}_{\tau}}\left[1_{\mathscr{H}_{U}}\left(X_{\cdot}^{x}\right)\right] \quad \text { a.s. }
$$

and hence $h_{U}(x)=\mathbb{E}\left[h_{U}\left(X_{\tau}^{x}\right)\right]$. Thus, by Lemma 2.2.14 (ii), the function $h_{U}$ is harmonic. In particular, if in addition $h_{U}(x) \in\{0,1\}$ for some $x \in M$, then already $h_{U} \equiv 0$ or $h_{U} \equiv 1$ according to the maximum principle.
(i) $\Rightarrow$ (ii): Let $h$ be a bounded harmonic function on $M$. Then $h\left(X^{x}\right)$ is a bounded and hence almost surely convergent martingale; let $\xi^{x}:=\lim _{t \uparrow \zeta^{x}} h\left(X_{t}^{x}\right)$. To $\alpha \in \mathbb{R}$ we consider the open set $U_{\alpha}:=\{h>\alpha\}$. By assumption and the maximum principle, for each of the harmonic functions $h_{U_{\alpha}}$ on $M$,

$$
h_{U_{\alpha}}(x)=\mathbb{P}\left\{X_{t}^{x} \in U_{\alpha} \text { eventually }\right\}=\mathbb{P}\left\{h \circ X_{t}^{x}>\alpha \text { eventually }\right\}
$$

it follows that either $h_{U_{\alpha}} \equiv 0$ or $h_{U_{\alpha}} \equiv 1$. For any real $\alpha$, hence $\mathbb{P}\left\{\xi^{x} \leq \alpha\right\} \in\{0,1\}$, independently of $x$. This shows that $\xi^{x} \equiv \lambda$ a.s. with a constant $\lambda$ independent of $x$. Hence $h$ is constant, namely $h(x)=\mathbb{E}\left[\xi^{x}\right] \equiv \lambda$.
(ii) $\Rightarrow$ (i): Since by assumption bounded harmonic functions on $M$ are constant, it holds in particular for any open subset $U \subset \hat{M}$ that

$$
h_{U}(x)=\mathbb{P}\left\{X_{t}^{x} \in U \text { eventually }\right\} \equiv \lambda \in[0,1] .
$$

We need to show that $h_{U}(M) \subset\{0,1\}$. But we have

$$
\lambda \equiv h_{U}\left(X_{t}^{x}\right)=\mathbb{E}^{\mathscr{F}_{t}}\left[1_{\mathscr{H}_{U}}\left(X_{\cdot}^{x}\right)\right] \rightarrow 1_{\mathscr{H}_{U}} \circ X_{\cdot}^{x} \quad \text { a.s. as } t \rightarrow \infty
$$

and hence $\lambda \in\{0,1\}$.
THEOREM 2.2.17. $\mathrm{BM}(M, g)$ is either recurrent or transient, i.e., for any Brownian motion $X$ on a metrically complete Riemannian manifold $(M, g)$ the following dichotomy holds: Either it holds
(i) $\underset{t \uparrow \zeta}{\liminf } d\left(X_{0}, X_{t}\right)=0$ a.s. or (ii) $\liminf _{t \uparrow \zeta} d\left(X_{0}, X_{t}\right)=\infty$ a.s.

Proof. For $x \in M$ let $X^{x}$ be a Brownian motion on $(M, g)$ starting from $x$ and

$$
A_{x}:=\left\{\liminf _{t \uparrow \zeta^{x}} d\left(X_{0}^{x}, X_{t}^{x}\right)=0\right\} \subset\left\{\zeta^{x}=\infty\right\}
$$

The function $h_{1}$ on $M$ defined by $h_{1}(x):=\mathbb{P}\left(A_{x}\right)$ is independent of the choice of the Brownian motion starting in $x$, and as consequence of the strong Markov property (Remark 2.2.5) harmonic on $M$ by Lemma 2.2.14 (ii). From

$$
\begin{equation*}
\mathbb{P}\left\{\liminf _{t \uparrow \zeta^{x}} d\left(X_{0}^{x}, X_{t}^{x}\right)=0\right\}>0 \tag{2.2.29}
\end{equation*}
$$

for one $x \in M$, hence (2.2.29) already follows for all $x \in M$. In addition, given a nonempty open subset $U \subset M$, then for any $x \in M$,

$$
\mathbb{P}\left\{X_{t}^{x} \in U \text { infinitely often }\right\}>0
$$

since also $h_{2}(x):=\mathbb{P}\left\{X_{t}^{x} \in U\right.$ infinitely often $\}$ is harmonic and $h_{2} \mid U>0$ by (2.2.29). But then, for each non-empty open subset $U \subset M$, it must already hold that

$$
\mathbb{P}\left\{X_{t}^{x} \in U \text { infinitely often }\right\}=1
$$

for any $x \in M$, since $h_{2}(X)$ is an almost surely convergent martingale and $X^{x}$ enters with positive probability every non-empty open subset infinitely often, from where it follows that $h_{2}$ is constant and hence $h_{2} \equiv 1$.

The alternative to the condition $\liminf _{t \uparrow \zeta^{x}} d\left(X_{0}^{x}, X_{t}^{x}\right)=0$ a.s. for each $x \in M$ is $\liminf _{t \uparrow \zeta^{x}} d\left(X_{0}^{x}, X_{t}^{x}\right)>0$ a.s. for one (and then each) $x \in M$. By the strong Markov property of the Brownian motion however the condition $\liminf _{t \uparrow \zeta^{x}} d\left(X_{0}^{x}, X_{t}^{x}\right)>0$ a.s. implies $X_{t}^{x} \rightarrow \infty$ a.s. as $t \uparrow \zeta^{x}$, and since $(M, g)$ is metrically complete by assumption, this means $\liminf _{t \uparrow \zeta^{x}} d\left(X_{0}^{x}, X_{t}^{x}\right)=\infty$ a.s.

The generalization of the dichotomy above to Brownian motions with not necessarily deterministic starting values is again a consequence of the strong Markov property.

Recurrence implies that Brownian motions return to any fixed non-empty open set infinitely often; their lifetime is hence infinite. In contrast, transience means that Brownian motions eventually leave every compact set. Historically, Riemannian manifolds with transient Brownian motion are called hyperbolic, whereas Riemannian manifolds with recurrent Brownian motion are called parabolic.

We now turn the discussion to the asymptotics of Brownian motions on model manifolds which is a type spaces typically used as comparison manifolds. From a probabilistic point of view, they have the property that in polar coordinates their radial and angular process can be decoupled by a time transformation and completely characterized (see also $[18,34])$. In this connection the angular behaviour is described by a martingale on the sphere, and hence the martingale convergence theorem allows to decide whether Brownian motion takes eventually an asymptotic direction or leaves every angular sector infinitely often. The same problem, namely existence of an asymptotic angle of Brownian motion, has been studied also for simply connected negatively curved Riemannian manifolds such that $-a^{2} \leq$ Riem $^{M} \leq-b^{2}<0$ in [44], [1]; see also [30].

Let $(M, g)$ be an $n$-dimensional Riemannian manifold with $n \geq 2$ and $x \in M$ such that the exponential map at $x$ defines a diffeomorphism; we may pull back the metric to $T_{x} M$ and identify $T_{x} M \cong \mathbb{R}^{n}$ by choosing an orthonormal basis in $T_{x} M$ :

$$
\begin{equation*}
\left(\mathbb{R}^{n}, \exp _{x}^{*} g\right) \cong\left(T_{x} M, \exp _{x}^{*} g\right) \underset{\exp _{x}}{\widetilde{ }}(M, g) \tag{2.2.30}
\end{equation*}
$$

In this way we identify $M$ and $\mathbb{R}^{n}$ also as sets. In geodesic polar coordinates on $M \backslash\{0\}=$ $] 0, \infty\left[\times \mathbb{S}^{n-1}\right.$ about $0 \in M$ we then have $g=d r \otimes d r+h_{r}$ with a Riemannian metric $h_{r}$ on $\mathbb{S}^{n-1}$ depending on $r$, where we consider the following special cases:
(a) $g=d r \otimes d r+f^{2}(r, \cdot) h$ where $\left.f:\right] 0, \infty\left[\times \mathbb{S}^{n-1} \rightarrow\right] 0, \infty[$ is a scalar function and $h$ a Riemannian metric on $\mathbb{S}^{n-1}$ which is independent of $r$.
(b) $g=d r \otimes d r+f^{2}(r) d \vartheta^{2}$ where $\left.f:\right] 0, \infty[\rightarrow] 0, \infty\left[\right.$ is a scalar function and $d \vartheta^{2}$ the standard metric on $\mathbb{S}^{n-1}$.

The situation (b) corresponds to the already treated model manifolds, whereas in (a) the induced metric on $\mathbb{S}^{n-1}$ is allowed to vary with the angle via the function $\left.f:\right] 0, \infty\left[\times \mathbb{S}^{n-1} \rightarrow\right.$ $] 0, \infty[$. In order to study the angular behaviour of $\operatorname{BM}(M, g)$ on such manifolds relative to 0 , we first investigate geometric properties of the angular map

$$
\begin{equation*}
q: M \backslash\{0\} \rightarrow \mathbb{S}^{n-1}, \quad(r, \vartheta) \mapsto \vartheta \tag{2.2.31}
\end{equation*}
$$

induced by (2.2.30).
Lemma 2.2.18. Let $q: M \backslash\{0\} \rightarrow \mathbb{S}^{n-1}$ be the angular map defined in (2.2.31).
(i) In situation (a) the map $q:(M \backslash\{0\}, g) \rightarrow\left(\mathbb{S}^{n-1}, h\right)$ is harmonic if and only if

$$
(n-3) \operatorname{grad} f_{r}=0
$$

where grad $f_{r}$ denotes the gradient vector field of $f_{r}=f(r, \cdot)$ on $\left(S^{n-1}, h\right)$.
(ii) In situation (b) the map $q:(M \backslash\{0\}, g) \rightarrow\left(\mathbb{S}^{n-1}, d \vartheta^{2}\right)$ is affine and in addition a harmonic morphism.

Proof. We want to calculate the second fundamental form of $q:(M \backslash\{0\}, g) \rightarrow$ $\left(\mathbb{S}^{n-1}, h\right)$ with respect to the fixed Riemannian metric on $\mathbb{S}^{n-1}$. As can be seen from
formula (1.7.2), in charts $(\varphi, U)$ for $M \backslash\{0\}$ and $(\psi, V)$ for $\mathbb{S}^{n-1}$ such that $\varphi(U) \subset V$ the following general representation in coordinates holds:

$$
(\nabla d q)_{i j}^{k}=\partial_{i} \partial_{j} q^{k}-\sum_{\alpha} \Gamma_{i j}^{\alpha}\left(\partial_{\alpha} q^{k}\right)+\sum_{\alpha, \beta} \bar{\Gamma}_{\alpha \beta}^{k}\left(\partial_{i} q^{\alpha}\right)\left(\partial_{j} q^{\beta}\right)
$$

with indices $1 \leq i, j \leq n, 1 \leq k \leq n-1$ and $\Gamma$ resp. $\bar{\Gamma}$ the Christoffel symbols with respect to the Levi-Civita connection on $(M \backslash\{0\}, g)$, resp. on $\left(\mathbb{S}^{n-1}, h\right)$. Thus choosing coordinates of the form $\varphi=\left(\theta^{1}, \ldots, \theta^{n-1}, r\right)$ for $M \backslash\{0\}$ with $r(\cdot)=d(0, \cdot)$ and $\psi=$ $\left(\theta^{1}, \ldots, \theta^{n-1}\right)$ for $\mathbb{S}^{n-1}$, one obtains

$$
(\nabla d q)_{i j}^{k}= \begin{cases}-\Gamma_{i j}^{k} & i=n \text { or } j=n \\ -\Gamma_{i j}^{k}+\bar{\Gamma}_{i j}^{k} & 1 \leq i, j \leq n-1\end{cases}
$$

As by (1.5.2) the Christoffel symbols of the Levi-Civita connection can be expressed via the Riemannian metric through

$$
\Gamma_{i j}^{k}=\frac{1}{2} \sum_{\ell} g^{k \ell}\left\{\partial_{i} g_{\ell j}+\partial_{j} g_{i \ell}-\partial_{\ell} g_{i j}\right\}
$$

we observe that $(\nabla d q)_{i j}^{k}=0$ for $i=n$ or $j=n$, as well as

$$
(\nabla d q)_{i j}^{k}=(1 / f) \sum_{\ell=1}^{n-1} h^{k \ell}\left\{h_{\ell j} \partial_{i} f+h_{i \ell} \partial_{j} f-h_{i j} \partial_{\ell} f\right\}, \quad 1 \leq i, j \leq n-1
$$

This shows at one hand that $q$ is affine in case (b), whereas in case (a)

$$
\tau(q)^{k}=\sum_{i, j=1}^{n} g^{i j}(\nabla d q)_{i j}^{k}=-(n-3) f^{-3} \sum_{i=1}^{n-1} h^{k i} \partial_{i} f
$$

holds, and hence $\tau(q)(r, \cdot)=-(n-3) f_{r}^{-3} \operatorname{grad} f_{r}$. It remains to verify that $q$ in case (b) defines in addition an harmonic morphism. Denoting by $\Delta$ the Laplacian on $(M, g)$ and accordingly by $\bar{\Delta}$ the Laplacian on $\left(\mathbb{S}^{n-1}, h\right)$, it is immediate to check that

$$
f^{2} \Delta(\varphi \circ q)=(\bar{\Delta} \varphi) \circ q
$$

for every differentiable function $\varphi \in C^{\infty}\left(\mathbb{S}^{n-1}\right)$, which according to Theorem 1.7.19 implies that $q$ is an harmonic morphism with Dilatation $f^{-1}$.

From the probabilistic point of view, Lemma 2.2.18 shows in particular that the angular process of a Brownian motion on a model is an $\mathbb{S}^{n-1}$-valued martingale with respect to the standard metric. This observation enables a complete description of the asymptotics of the angular behaviour.

THEOREM 2.2.19 (Brownian motion on models). Let $M$ be a $n$-dimensional model with center $0 \in M$ and Riemannian metric $g=d r \otimes d r+f^{2}(r) d \vartheta^{2}$. Let $X$ be a Brownian motion on $(M, g)$ with $X_{0}=x_{0} \neq 0$, decomposed according to $\left.M \backslash\{0\}=\right] 0, \infty\left[\times \mathbb{S}^{n-1}\right.$ in its radial and angular part $X=(R, \Theta)$.
(i) For the radial process, $R_{t} \rightarrow \infty$ almost surely (i.e., $X$ is transient) if and only if

$$
\int_{1}^{\infty} f^{1-n}(r) d r<\infty
$$

(ii) The lifetime $\zeta$ of $X$ is either a.s. finite or a.s. infinite, and a.s. finite if and only if

$$
\int_{1}^{\infty} f^{n-1}(r)\left\{\int_{r}^{\infty} f^{1-n}(\rho) d \rho\right\} d r<\infty
$$

(iii) The angular process $\Theta$ converges on $\mathbb{S}^{n-1}$ for $t \uparrow \zeta$ a.s., if and only if

$$
\int_{1}^{\infty} f^{n-3}(r)\left\{\int_{r}^{\infty} f^{1-n}(\rho) d \rho\right\} d r<\infty
$$

The latter is equivalent to $M$ being no Liouville manifold.
Note that in case $d=2$, a model $(M, g)$ is a Liouville manifold if and only if $\mathrm{BM}(M, g)$ is recurrent.

Proof. (1) By Theorem 2.2.1 it holds that $d R=d W+\frac{1}{2} \Delta r(X) d t$ where $W$ is a one-dimensional Brownian motion. Since $\Delta r(X)=(n-1)\left(f^{\prime} / f\right)(R)$, the radial process $R$ satisfies the SDE

$$
d R=d W+\frac{1}{2}(n-1)\left(f^{\prime} / f\right)(R) d t
$$

Thus $R$ is a one-dimensional diffusion with infinitesimal generator

$$
\frac{1}{2}\left\{D^{2}+(n-1)\left(f^{\prime} / f\right) D\right\}
$$

We want to calculate the Riemannian quadratic variation of the martingales $\Theta$. To this end, we first note that in each chart $(\varphi, U)$ of the form $\varphi=(r, \theta)$ and $U=] 0, \infty\left[\times U^{\prime}\right.$ where $\left(\theta, U^{\prime}\right)$ is a chart for $\mathbb{S}^{n-1}$ obviously

$$
\begin{aligned}
1_{\{X \in U\}} d[X, X] & =1_{\{X \in U\}} \sum_{i, j=1}^{n} g_{i j}(X) d\left[X^{i}, X^{j}\right] \\
& =1_{\{X \in U\}}\left\{d[R, R]+f^{2}(R) d[\Theta, \Theta]\right\}
\end{aligned}
$$

holds, from where we conclude that

$$
d[X, X]=d[R, R]+f^{2}(R) d[\Theta, \Theta] .
$$

Taking into account that $d[X, X]=(\operatorname{dim} M) d t=n d t$ and $d[R, R]=d t$, we finally obtain

$$
\begin{equation*}
d[\Theta, \Theta]=(n-1) f^{-2}(R) d t \tag{2.2.32}
\end{equation*}
$$

By Theorem 1.8.8 (convergence theorem of Darling-Zheng), hence $\Theta$ converges on $\mathbb{S}^{n-1}$ for $t \uparrow \zeta$ almost surely if and only if $\int_{0}^{\zeta} f^{-2}\left(R_{t}\right) d t<\infty$ almost surely. We let

$$
T(t)=\int_{0}^{t} f^{-2}\left(R_{s}\right) d s, \quad t<\zeta
$$

and consider for $t<T_{\zeta}$ the continuous time change $\left(\tau_{t}\right)$ where

$$
\tau_{t}:=T^{-1}(t) \equiv \inf \left\{s \in \mathbb{R}_{+}: T(s) \geq t\right\}
$$

Since $T_{\zeta}$ is obviously the maximal lifetime of the time transformed radial process $\tilde{R}_{t}:=$ $R_{\tau_{t}}$, we get as consequence that $\Theta$ converges for $t \uparrow \zeta$ if and only the lifetime of $\tilde{R}$ is finite, almost surely.

By Lemma 2.2.18 the map $q$ is an harmonic morphism with dilatation $f^{-1}$; consequently $X$ decomposes as $X_{t}=\left(R_{t}, B_{T(t)}\right)$ with $B$ a Brownian motion on $\left(\mathbb{S}^{n-1}, d \vartheta^{2}\right)$. By the time-change $\left(\tau_{t}\right)$ the radial and angular process decompose as $X_{\tau_{t}}=\left(R_{\tau_{t}}, B_{t}\right)$ for $t<T_{\zeta}$; in the new clock the angular component is described by a $\operatorname{BM}\left(\mathbb{S}^{n-1}, d \vartheta^{2}\right)$ which runs up to time $T_{\zeta}$; hence it converges if and only if $T_{\zeta}$ is almost surely finite.
(2) By Eq. (2.2.1) and Theorem 2.1.57 (ii) the radial process $R$ solves the one-dimensional SDE

$$
d R=d \hat{W}+\frac{1}{2}(n-1)\left(f^{\prime} / f\right)(R) d t
$$

hence $R$ is a diffusion on $] 0, \infty[$ with infinitesimal generator

$$
\frac{1}{2}\left(D^{2}+(n-1)\left(f^{\prime} / f\right) D\right)
$$

Correspondingly the time-changed process $\tilde{R}$ is a diffusion on $] 0, \infty[$ with generator

$$
\frac{1}{2} f^{2}\left(D^{2}+(n-1)\left(f^{\prime} / f\right) D\right)
$$

The questions that interest us here concerning transience and lifetime of $R$ resp. $\tilde{R}$ can hence be answered by means of Theorem A.1.9. Set $c_{1}=0, c_{2}=\infty$ and without restrictions $c=1$. Consider

$$
H(r)=\exp \left(-(n-1) \int_{1}^{r} \frac{f^{\prime}}{f}(s) d s\right)=f^{1-n}(r) f^{n-1}(1), \quad 0<r<\infty
$$

As seen in the proof of Theorem 2.2.1, almost surely, Brownian motion does not hit any point fixed in advance; hence 0 is both for $R$ and $\tilde{R}$ a non-accessible boundary point, and it is thus sufficient to investigate the respective behaviour at the right-hand boundary point $c_{2}=\infty$. In detail we find:

$$
\begin{aligned}
R_{t} \rightarrow \infty \text { a.s. (transient) } & \Longleftrightarrow \int_{1}^{\infty} f^{1-n}(r) d r<\infty ; \\
R \text { has finite lifetime a. s. } & \Longleftrightarrow \int_{1}^{\infty} H(r)\left[\int_{1}^{r} \frac{d \rho}{H(\rho)}\right] d r<\infty \\
& \Longleftrightarrow \int_{1}^{\infty} \frac{1}{H(r)}\left[\int_{r}^{\infty} H(\rho) d \rho\right] d r<\infty \\
& \Longleftrightarrow \int_{1}^{\infty} f^{n-1}(r)\left[\int_{r}^{\infty} f^{1-n}(\rho) d \rho\right] d r<\infty ; \\
\tilde{R} \text { has finite lifetime a. s. } & \Longleftrightarrow \int_{1}^{\infty} f^{n-3}(r)\left[\int_{r}^{\infty} f^{1-n}(\rho) d \rho\right] d r<\infty
\end{aligned}
$$

(3) It remains to show that $M$ supports con-constant bounded harmonic functions exactly if the angular process $\Theta$ converges almost surely on the sphere $\mathbb{S}^{n-1}$ as $t \uparrow \zeta$.

We consider first the case that $M$ supports non-constant bounded harmonic functions. Then $(M, g)$ is not a Liouville manifold, and by Theorem 2.2.16 there exist non-trivial exit set for $\mathrm{BM}(M, g)$. It is easy to see that among them some must be of the form $\mathbb{R}_{+} \times V$ where $V \subset \mathbb{S}^{n-1}$ is an open subset. This excludes the possibility that $T_{\zeta} \equiv \infty$ almost surely, as a consequence of the recurrence of $\operatorname{BM}\left(\mathbb{S}^{n-1}, d \vartheta^{2}\right)$. Recall that $T_{\zeta}$ is the lifetime of the time-changed radial process $\tilde{R}$ and hence almost surely infinite or almost surely finite. Thus $T_{\zeta}$ almost surely finite must hold, which, as already shown, is equivalent to convergence of the angular process $\Theta$ on $\mathbb{S}^{n-1}$.

Conversely, suppose that $\Theta$ converges almost surely on the sphere $\mathbb{S}^{n-1}$ as $t \uparrow \zeta$. Write $\Theta^{x}$ for the angular part of a Brownian motion $X^{x}$ starting at $x \in M$. For each function $\varphi \in C\left(\mathbb{S}^{n-1}\right)$ then

$$
u(x):=\mathbb{E}\left[\varphi\left(\Theta_{\zeta}^{x}\right)\right]
$$

defines a bounded harmonic function $u$ on $M$. Since

$$
u\left(X_{t \wedge \zeta}^{x}\right) \rightarrow \varphi\left(\Theta_{\zeta}^{x}\right) \quad \text { almost surely as } t \rightarrow \infty
$$

the considered function $u$ is non-constant if and only if $\varphi\left(\Theta_{\zeta}^{x}\right)$ is non-degenerate on $\mathbb{S}^{n-1}$, i.e. not almost surely constant. If however the "exit measure" $\mathbb{P} \circ\left(\Theta_{\zeta}^{x}\right)^{-1}$ equals the Dirac measure $\delta$ of a point on $\mathbb{S}^{n-1}$, then as an easy application of the maximum principle
shows, the measures $\mathbb{P} \circ\left(\Theta_{\zeta}^{x}\right)^{-1}$ for $x \in M$ must be identical, hence $\mathbb{P} \circ\left(\Theta_{\zeta}^{x}\right)^{-1}=\delta$, independently of $x$. This is in contradiction to the rotational invariance of $\mathbb{P} \circ\left(\Theta_{\zeta}^{0}\right)^{-1}$ on $\mathbb{S}^{n-1}$ which comes from the fact that $(M, g)$ is a model.

It is straight-forward to translate the integral conditions in Theorem 2.2.19 into curvature bounds. This then enables a geometric characterization of the Liouville property.

THEOREM 2.2.20. Let $(M, g)$ be an $n$-dimensional model with center $0 \in M$ and Riemannian metric $g=d r \otimes d r+f^{2}(r) d \vartheta^{2}$. For the radial curvature function $k_{M}(r)=$ $-f^{\prime \prime}(r) / f(r)$ of $(M, g)$ assume that $k_{M}(\cdot) \leq 0$. Furthermore let $c=1$ in case $n=2$, respectively $c=1 / 2$ in case $n \geq 3$. Then there exist non-constant harmonic functions on $M$, if $k_{M}(r) \leq-\frac{(c+\varepsilon)}{r^{2} \log r}$ for some $\varepsilon>0$ and sufficiently large $r$. In contrary, if $k_{M}(r) \geq-\frac{(c-\varepsilon)}{r^{2} \log r}$ for some $\varepsilon>0$ and sufficiently large $r$, then $M$ is a Liouville manifold.

Note that constant negative curvature outside a compact set is not sufficient for the existence of non-constant bounded harmonic functions. The condition is not even enough for transience of Brownian motion, as the following example shows. Let $M=\mathbb{R}^{2}$ be a twodimensional rotationally symmetric manifold, for instance, with radial function $f(r)=$ $\exp (-r)$ for $r>1$ and a differentiable interpolation for $0 \leq r \leq 1$, such that $f(0)=0$ and $f^{\prime}(0)=1$ holds. Then $(M, g)$ has constant negative curvature outside the unit disk, but according to Theorem 2.2.19, Brownian motion on $(M, g)$ is recurrent; $M$ is hence a Liouville manifold.

We proceed with an elementary Lemma before turning to the proof of Theorem 2.2.20.
LEMMA 2.2.21. Let $n \geq 2$ and $(M, g)$ be an $n$-dimensional model with Riemannian metric $g=d r \otimes d r+f^{2}(r) d \vartheta^{2}$. Let $k=-f^{\prime \prime} / f$ be the radial curvature function of $(M, g)$ and

$$
I(f)=\int_{1}^{\infty} f^{n-3}(r)\left[\int_{r}^{\infty} f^{1-n}(\rho) d \rho\right] d r
$$

Furthermore, let $(\tilde{M}, \tilde{g})$ another $n$-dimensional model with metric $\tilde{g}=d r \otimes d r+\tilde{f}^{2}(r) d \vartheta^{2}$, and define $\tilde{k}=-\tilde{f}^{\prime \prime} / \tilde{f}$ and $I(\tilde{f})$ correspondingly.
(i) If $\tilde{k} \leq k$ on $] 0, \infty\left[\right.$ then $f \leq \tilde{f}$ and $f^{\prime} / f \leq \tilde{f}^{\prime} / \tilde{f}$ on $] 0, \infty[$.
(ii) If $\tilde{k} \leq k$ on $] \rho_{0}, \infty\left[\right.$ for some $\rho_{0}>0$, then $\left(f^{\prime} / f\right)\left(\rho_{0}\right) \leq\left(\tilde{f}^{\prime} / \tilde{f}\right)\left(\rho_{0}\right)$ implies $f^{\prime} / f \leq$ $\tilde{f}^{\prime} / \tilde{f}$ on the interval $\left[\rho_{0}, \infty[\right.$.
(iii) If $k, \tilde{k} \leq 0$ and $\tilde{k} \leq k$ on $] \rho_{0}, \infty\left[\right.$ for some $\rho_{0}>0$, then there exists a constant $c>0$ such that $f \leq c \tilde{f}$ on $[0, \infty[$,
(iv) If $k, \tilde{k} \leq 0$ and $\tilde{k} \leq k$ on $] \rho_{0}, \infty\left[\right.$ for some $\rho_{0}>0$, then with $I(f)<\infty$ also $I(\tilde{f})<\infty$.

Proof. (1) The claims in (i) and (ii) follow immediately from

$$
\begin{equation*}
\left(f^{2}(\tilde{f} / f)^{\prime}\right)^{\prime}=\left(f \tilde{f}^{\prime}-\tilde{f} f^{\prime}\right)^{\prime}=\tilde{f} f\left(\tilde{f}^{\prime \prime} / \tilde{f}-f^{\prime \prime} / f\right)=\tilde{f} f(k-\tilde{k}) \tag{2.2.33}
\end{equation*}
$$

(2) Let now $k, \tilde{k} \leq 0$. From $k \leq 0$ we deduce by (i) that $r \leq f(r)$ and $1 / r \leq$ $f^{\prime}(r) / f(r)$ for $0<r<\infty$; correspondingly for $\tilde{f}$. In particular, $f^{\prime}$ and $\tilde{f}^{\prime}$ are strictly positive on $] 0, \infty\left[\right.$. We set $k_{+}:=k \vee \tilde{k}, k_{-}:=k \wedge \tilde{k}$ and consider the $C^{2}$ solutions $u_{ \pm}$of

$$
u^{\prime \prime}+k_{ \pm} u=0 \quad \text { with } u(0)=0, u^{\prime}(0)=1
$$

As above, we conclude that $u_{ \pm}$and $u_{ \pm}^{\prime}$ are strictly positive on $] 0, \infty[$; in particular

$$
\begin{aligned}
& c_{1} u_{+}\left(\rho_{0}\right) \leq f\left(\rho_{0}\right) \leq c_{2} u_{+}\left(\rho_{0}\right) \\
& c_{1} u_{+}^{\prime}\left(\rho_{0}\right) \leq f^{\prime}\left(\rho_{0}\right) \leq c_{2} u_{+}^{\prime}\left(\rho_{0}\right)
\end{aligned}
$$

with appropriate constants $c_{1}, c_{2}>0$. But $u_{+}$and $f$ satisfy identical differential equations on $] \rho_{0}, \infty[$ which implies

$$
\begin{equation*}
c_{1} u_{+}(r) \leq f(r) \leq c_{2} u_{+}(r), \quad r \geq \rho_{0} . \tag{2.2.34}
\end{equation*}
$$

Analogously, for suitable constants $\tilde{c}_{1}, \tilde{c}_{2}>0$, we obtain the inequality

$$
\begin{equation*}
\tilde{c}_{1} u_{-}(r) \leq \tilde{f}(r) \leq \tilde{c}_{2} u_{-}(r), \quad r \geq \rho_{0} . \tag{2.2.35}
\end{equation*}
$$

By (2.2.33), using $u_{+}, u_{-}$instead of $f, \tilde{f}$, we conclude $u_{+} \leq u_{-}$on $[0, \infty[$; the combination of (2.2.34) and (2.2.35) then gives $f \leq c \tilde{f}$ for some constant $c>0$; at first on $\left[\rho_{0}, \infty[\right.$ and after eventually enlarging $c$ then on all of $[0, \infty[$.
(3) The claim in part (iv) can be easily reduced to the case $\rho_{0}=0$ by following the arguments used in (2): $\operatorname{By}(2.2 .34)$ the condition $I(f)<\infty$ is equivalent $I\left(u_{+}\right)<\infty$; thus one may replace $f$ by $u_{+}$with the consequence that then $\tilde{k} \leq k_{+}$holds on all of $] 0, \infty[$. Without restrictions we may thus assume that $\rho_{0}=0$. We let

$$
I_{s}(f)=\int_{1}^{s} f^{n-3}(r)\left[\int_{r}^{\infty} f^{1-n}(\rho) d \rho\right] d r, \quad s \geq 1
$$

and will show that under the condition $\tilde{k} \leq k \leq 0$ it holds that

$$
I_{s}(\tilde{f}) \leq I_{s}(f), \quad s \geq 1
$$

To this end we may assume that $\int_{1}^{\infty} f^{1-n}(\rho) d \rho<\infty$ and $n \geq 3$. At first we then have

$$
\frac{d}{d s} I_{s}(f)=f^{n-3}(s)\left[\int_{s}^{\infty} f^{1-n}(\rho) d \rho\right] \geq 0
$$

and hence

$$
\begin{equation*}
f^{3-n}(s) \frac{d}{d s} I_{s}(f)=\int_{s}^{\infty} f^{1-n}(\rho) d \rho \searrow 0 \quad \text { as } s \uparrow \infty . \tag{2.2.36}
\end{equation*}
$$

Differentiation of (2.2.36) gives

$$
\frac{d}{d s}\left[f^{3-n}(s) \frac{d}{d s} I_{s}(f)\right]=-f^{1-n}(s)
$$

and then

$$
\frac{d^{2}}{d s^{2}} I_{s}(f)-(n-3) \frac{f^{\prime}(s)}{f(s)} \frac{d}{d s} I_{s}(f)+\frac{1}{f(s)^{2}}=0
$$

Using the assumptions along with part (i), we get

$$
\frac{d^{2}}{d s^{2}} I_{s}(f)-(n-3) \frac{\tilde{f}^{\prime}(s)}{\tilde{f}(s)} \frac{d}{d s} I_{s}(f)+\frac{1}{\tilde{f}(s)^{2}} \leq 0
$$

and hence
$-\tilde{f}^{1-n}(s) \geq \tilde{f}^{3-n}(s) \frac{d^{2}}{d s^{2}} I_{s}(f)-(n-3) \frac{\tilde{f}^{\prime}(s)}{\tilde{f} n-2(s)} \frac{d}{d s} I_{s}(f)=\frac{d}{d s}\left[\tilde{f}^{3-n}(s) \frac{d}{d s} I_{s}(f)\right]$.
Integration then gives

$$
-\int_{r}^{\infty} \tilde{f}^{1-n}(\rho) d \rho \geq \lim _{u \rightarrow \infty}\left[\tilde{f}^{3-n}(s) \frac{d}{d s} I_{s}(f)\right]_{s=r}^{s=u}
$$

and thus

$$
\tilde{f}^{n-3}(r) \int_{r}^{\infty} \tilde{f}^{1-n}(\rho) d \rho \leq \frac{d}{d r} I_{r}(f)
$$

Integrating again finally gives for $s \geq 1$

$$
I_{s}(\tilde{f})=\int_{1}^{s} \tilde{f}^{n-3}(r)\left[\int_{r}^{\infty} \tilde{f}^{1-n}(\rho) d \rho\right] d r \leq I_{s}(f)-I_{1}(f)=I_{s}(f)
$$

which is the claim.
Proof of Theorem 2.2.20. For $\alpha>0$ and $\varrho>1$ let $\Phi(r):=r(\log r)^{\alpha}$ and $F(r):=(\Phi(r+\varrho)-\Phi(\varrho)) / \Phi^{\prime}(\varrho)$. Then it holds that $F(0)=0, F^{\prime}(0)=1$ and $-\left(F^{\prime \prime} / F\right)(r) \leq 0$ for $r>0$, as well as $-\left(F^{\prime \prime} / F\right)(r) \sim-\alpha /\left(r^{2} \log r\right)$ as $r \rightarrow \infty$. Using that

$$
\int_{2}^{\infty} \frac{d r}{r(\log r)^{a}}<\infty \Longleftrightarrow a>1
$$

we obtain in case $n=2$ that $I(F)<\infty$ holds if and only if $\alpha>1$. Analogously one sees in case $n \geq 3$, by using

$$
\int_{r}^{\infty} F(\rho)^{1-n} d \rho \sim\left[(n-2) r^{n-2}(\log r)^{(n-1) \alpha}\right]^{-1} \quad \text { as } r \rightarrow \infty
$$

that $I(F)<\infty$ holds if and only if $\alpha>1 / 2$.
Through the dependence on curvature of the exit time of Brownian motions from geodesic balls, in connection with the explicit knowledge of the asymptotics of Brownian motions on models, comparison theorems play an important role for applications (see [19, 20]).

THEOREM 2.2.22 (Comparison criterion for BM-completeness and transience).
A. Let $(M, g)$ be a metrically complete Riemannian manifold of dimension $n \geq 2$. Suppose there is a point $o \in M$ and a model $(\mathbb{M}, \tilde{g})$ of equal dimension with center 0 such that

$$
\operatorname{Ric}_{x}^{M}\left(\partial^{M}, \partial^{M}\right) \geq(n-1) k_{\mathbb{M}}(r)
$$

for $x \in M \backslash \operatorname{cut}(o)$ and $0<r=d_{\mathbb{M}}(0, x)$. Then if $(\mathbb{M}, \tilde{g})$ is $\mathbf{B M}$-complete also $(M, g)$ is BM -complete; if $\mathrm{BM}(\mathbb{M}, \tilde{g})$ is recurrent, then also $\operatorname{BM}(M, g)$.
B. Let $(M, g)$ be a metrically complete Riemannian manifold of dimension $n \geq 2$. Suppose there is a point $o \in M$ and a model $(\mathbb{M}, \tilde{g})$ of equal dimension with center such that

$$
\operatorname{Riem}_{x}^{M}(E) \leq k_{\mathbb{M}}(r)
$$

for each radial plane $E$ in $T_{x} M$ with $x \in M \backslash \operatorname{cut}(o)$ and $0<r=d_{\mathbb{M}}(0, x)$. If $M$ is simply connected, then with $(\mathbb{M}, \tilde{g})$ also $(M, g)$ is not BM -complete and $\mathrm{BM}(M, g)$ has finite lifetime almost surely; in this case $\mathrm{BM}(M, g)$ is transient if $\mathrm{BM}(\mathbb{M}, \tilde{g})$ is transient.

Proof. Part A follows from Theorem 2.2.13; Part B follows from Theorem 2.2.2, along with the observation that $\operatorname{cut}(o)=\varnothing$ under the given assumptions. Indeed, by the given curvature assumptions in case B , along normal minimal geodesic curves $\gamma:[0, a] \rightarrow$ $M$ with $\gamma(0)=o$ one has

$$
\operatorname{Riem}_{\gamma(r)}^{M}(E) \leq k_{\mathbb{M}}(r)
$$

for each plane $E \subset T_{\gamma(r)} M$ such that $\dot{\gamma}(r) \in E$. The Comparison principle (Corollary 2.1.40) then gives $\operatorname{Conj}(o)=\varnothing$. Thus $\exp _{o}:\left(T_{o} M, \exp _{o}^{*} g\right) \rightarrow(M, g)$ is a local isometry and hence a covering (as already justified in the proof of Theorem 2.1.41). For simply connected $M$ hence $\exp _{o}: T_{o} M \rightarrow M$ defines a diffeomorphism which in particular shows
$\operatorname{cut}(o)=\varnothing$. The claims then follow from the Comparison Theorem for Brownian motion (Theorem 2.2.2). Indeed, if for instance ( $\mathbb{M}, \tilde{g}$ ) is not BM-complete, then Brownian motions on $(\mathbb{M}, \tilde{g})$ have only a finite lifetime almost surely, and one concludes by Theorem 2.2.2 (and the notation there) that also the lifetime $\zeta \equiv \sup _{\rho} \tau_{\rho}$ for each Brownian motion on $(M, g)$ must be finite almost surely; almost sure explosion follows first for Brownian motions with deterministic starting point, but then also in general by means of the Markov property 2.2.9.

COROLLARY 2.2.23 (Test for BM-completeness). Let $(M, g)$ be a metrically complete Riemannian manifold of dimension $n \geq 2$. For a given point $o \in M$ suppose that the Ricci curvature in the radial direction is bounded below by quadratic function of the distance $r_{M}=d_{M}(o, \cdot)$ to o, i.e.

$$
\operatorname{Ric}_{x}^{M}\left(\partial^{M}, \partial^{M}\right) \geq-c_{1}-c_{2} r_{M}^{2}(x), \quad c_{1}, c_{2}>0
$$

Then $(M, g)$ is BM-complete, i.e., Brownian motions on $(M, g)$ have infinite lifetime.
Proof. We compare $(M, g)$ to the model $(\mathbb{M}, \tilde{g})=\left(\mathbb{R}^{n}, d r \otimes d r+f^{2}(r) d \vartheta^{2}\right)$ where $f(r)=r \exp \left(c r^{2}\right)$ with $c>0$. Then it holds that $-\left(f^{\prime \prime} / f\right)(r)=-\left(6 c+4 c^{2} r^{2}\right)$, and then

$$
\operatorname{Ric}_{x}^{M}\left(\partial^{M}, \partial^{M}\right) \geq-(n-1)\left(f^{\prime \prime} / f\right)\left(r_{M}(x)\right)
$$

by choosing the constant $c$ appropriately. The claim then follows from

$$
\int_{1}^{\infty} f^{n-1}(r)\left\{\int_{r}^{\infty} f^{1-n}(\rho) d \rho\right\} d r=\infty
$$

(which is easy to verify) by Theorem 2.2 .22 A .

### 2.3. Heat Equation on Sections of a Vector Bundle

In this section we use martingale arguments to derive stochastic formulae for the differential of diffusion semigroups on functions and the codifferential on 1-forms. In particular, for a semigroup generated by some elliptic operator, we aim at probabilistic representations for $d P_{t} f$ and $P_{t}(V f)$, not involving derivatives of $f$, where $V$ is a vector field. Such formulae are typically referred to under the name of Bismut type formulae [3]. For nonsymmetric generators, such formulae for $d P_{t} f$ and $P_{t}(V f)$ are related to the derivative of the heat kernel in the backward, resp. forward variable.
2.3.1. Notations. Let $(M, g)$ be an $n$-dimensional Riemannian manifold, $\nabla$ be the Levi-Civita connection on $M$ and $\pi: \mathrm{O}(T M) \rightarrow M$ be the orthonormal frame bundle over $M$. Let $E \rightarrow M$ be an associated vector bundle with fiber $V$ and structure group $G=\mathrm{O}(n)$. The induced covariant derivative

$$
\nabla: \Gamma(E) \rightarrow \Gamma\left(T^{*} M \otimes E\right)
$$

determines the so-called connection Laplacian (or rough Laplacian) $\square$ on $\Gamma(E)$,

$$
\square a=\operatorname{trace} \nabla^{2} a
$$

Note that $\nabla^{2} a \in \Gamma\left(T^{*} M \otimes T^{*} M \otimes E\right)$ and hence $(\square a)_{x}=\sum_{i} \nabla^{2} a\left(v_{i}, v_{i}\right) \in E_{x}$ where $v_{i}$ runs through an orthonormal basis of $T_{x} M$. For $a, b \in \Gamma(E)$ of compact support it is immediate to check that

$$
\langle\square a, b\rangle_{L^{2}(E)}=-\langle\nabla a, \nabla b\rangle_{L^{2}\left(T^{*} M \otimes E\right)} .
$$

In this sense we have $\square=-\nabla^{*} \nabla$.

Let $H$ be the horizontal subbundle of the $G$-invariant splitting of $T \mathrm{O}(T M)$ and

$$
h: \pi^{*} T M \xrightarrow{\sim} H \longleftrightarrow T \mathrm{O}(T M)
$$

be the horizontal lift of the $G$-connection; fiberwise this bundle isomorphism reads as reads as

$$
h_{u}: T_{\pi(u)} M \xrightarrow{\sim} H_{u}, \quad u \in \mathrm{O}(T M)
$$

In terms of the standard-horizontal vector fields $H_{1}, \ldots, H_{n}$ on $\mathrm{O}(T M)$,

$$
H_{i}(u):=h_{u}\left(u e_{i}\right), \quad u \in \mathrm{O}(T M)
$$

Bochner's horizontal Laplacian $\Delta^{\text {hor }}$, acting on smooth functions on $\mathrm{O}(T M)$, is given as

$$
\Delta^{\mathrm{hor}}=\sum_{i=1}^{n} H_{i}^{2}
$$

To formulate the relation between $\square$ and $\Delta^{\text {hor }}$, it is convenient to write sections $a \in$ $\Gamma(E)$ as "equivariant functions" $F_{a}: \mathrm{O}(T M) \rightarrow V$ via

$$
F_{a}(u)=u^{-1} a_{\pi(u)}
$$

where we read $u \in \mathrm{O}(T M)$ as isomorphism $u: V \xrightarrow{\sim} E_{\pi(u)}$. Equivariance means that

$$
F_{a}(u g)=g^{-1} F_{a}(u), \quad u \in \mathrm{O}(T M), g \in G=\mathrm{O}(n)
$$

LEMMA 2.3.1 (see [26], p. 115). For $a \in \Gamma(E)$ and $F_{a}$ the corresponding equivariant function on $\mathrm{O}(T M)$, we have

$$
\left(H_{i} F_{a}\right)(u)=F_{\nabla_{u e_{i}} a}(u), \quad u \in \mathrm{O}(T M)
$$

Hence

$$
\Delta^{\mathrm{hor}} F_{a}=F_{\square a} \text {, }
$$

where as above

$$
\square: \Gamma(E) \xrightarrow{\nabla} \Gamma\left(T^{*} M \otimes E\right) \xrightarrow{\nabla} \Gamma\left(T^{*} M \otimes T^{*} M \otimes E\right) \xrightarrow{\text { trace }} \Gamma(E) .
$$

Proof. Fix $u \in \mathbf{O}(T M)$ and choose a curve $\gamma$ in $M$ such that $\gamma(0)=\pi(u)$ and $\dot{\gamma}=u e_{i}$. Let $t \mapsto u(t)$ be the horizontal lift of $\gamma$ to $\mathrm{O}(T M)$ such that $u(0)=u$. Note that $\dot{u}(t)=h_{u(t)}(\dot{\gamma}(t))$, and in particular $\dot{u}(0)=h_{u}\left(u e_{i}\right)=H_{i}(u)$. Hence, denoting the parallel transport along $\gamma$ by $/ \|_{\varepsilon}=u(\varepsilon) u(0)^{-1}$, we get

$$
\begin{aligned}
F_{\nabla_{u e_{i}}} a(u) & =u^{-1}\left(\nabla_{u e_{i}} a\right)_{\pi(u)} \\
& =u^{-1} \lim _{\varepsilon \downarrow 0} \frac{/ /_{\varepsilon}^{-1} a_{\gamma(\varepsilon)}-a_{\gamma(0)}}{\varepsilon} \\
& =\lim _{\varepsilon \downarrow 0} \frac{u(\varepsilon)^{-1} a_{\gamma(\varepsilon)}-u(0)^{-1} a_{\gamma(0)}}{\varepsilon} \\
& =\lim _{\varepsilon \downarrow 0} \frac{F_{a}(u(\varepsilon))-F_{a}(u(0))}{\varepsilon} \\
& =\left(H_{i}\right)_{u} F_{a}=\left(H_{i} F_{a}\right)(u) .
\end{aligned}
$$

2.3.2. Preliminaries. Consider diffusion processes $X_{t}$ on $M$ generated by the operator

$$
\begin{equation*}
L=\Delta+Z \tag{2.3.1}
\end{equation*}
$$

where $Z \in \Gamma(T M)$ is a smooth vector field on $M$. We assume that $X_{t}$ is non-explosive. Such diffusions on $M$ may be constructed from the corresponding horizontal diffusions on $\mathrm{O}(T M)$ generated by

$$
\Delta^{\mathrm{hor}}+\bar{Z}
$$

where the vector field $\bar{Z}$ is the horizontal lift of $Z$ to $\mathrm{O}(T M)$, i.e. $\bar{Z}_{u}=h_{u}\left(Z_{\pi(u)}\right), u \in$ $\mathrm{O}(T M)$. More precisely, we start from the Stratonovich stochastic differential equation on $\mathrm{O}(T M)$,

$$
d U=\sum_{i=1}^{n} H_{i}(U) \circ d B^{i}+\bar{Z}(U) d t, \quad U_{0}=u \in \mathbf{O}(T M)
$$

Then for $X_{t}=\pi\left(U_{t}\right)$, the following equation holds:

$$
d X=\sum_{i=1}^{n} U e_{i} \circ d B^{i}+Z(X) d t, \quad X_{0}=x:=\pi u
$$

Here $B$ is a Brownian motion on $\mathbb{R}^{n}$, speeded up by the factor 2, i.e., $d B^{i} d B^{j}=2 \delta_{i j} d t$; recall that $B$ is the martingale part of the anti-development $\int_{U} \vartheta$ of $X$ where $\vartheta$ denotes the canonical 1-form $\vartheta$ on $\mathrm{O}(T M)$, i.e. $\vartheta_{u}(e)=u^{-1} e_{\pi(u)}, e \in T_{u} \mathrm{O}(T M)$.

In particular, for $F \in C^{\infty}(\mathrm{O}(T M))$, resp. $f \in C^{\infty}(M)$, we get

$$
\begin{align*}
d(F \circ U) & =\sum_{i=1}^{n}\left(H_{i} F\right)(U) \circ d B^{i}+(\bar{Z} F)(U) d t \\
& =\sum_{i=1}^{n}\left(H_{i} F\right)(U) d B^{i}+\left(\Delta^{\mathrm{hor}}+\bar{Z}\right)(F)(U) d t \tag{2.3.2}
\end{align*}
$$

resp.

$$
\begin{align*}
d(f \circ X) & =\sum_{i=1}^{n} d f\left(U e_{i}\right) \circ d B^{i}+(Z f)(X) d t \\
& =\sum_{i=1}^{n} d f\left(U e_{i}\right) d B^{i}+(\Delta+Z)(f)(X) d t \tag{2.3.3}
\end{align*}
$$

Proposition 2.3.2. Let $/{ }_{t}: E_{X_{0}} \rightarrow E_{X_{t}}$ be parallel transport in $E$ along $X$, induced by the parallel transport on $M$,

$$
/ /_{t}=U_{t} U_{0}^{-1}: T_{X_{0}} M \rightarrow T_{X_{t}} M
$$

Then, for $a \in \Gamma(E)$, we have

$$
\begin{equation*}
d\left(/ / t^{-1} a\left(X_{t}\right)\right)=\sum_{i=1}^{n} / /_{t}^{-1}\left(\nabla_{U_{t} e_{i}} a\right) \circ d B^{i}+/ / t^{-1}\left(\nabla_{Z} a\right)\left(X_{t}\right) d t \tag{2.3.4}
\end{equation*}
$$

resp. in Itô form,

$$
\begin{equation*}
d\left(/ /\left.\right|_{t} ^{-1} a\left(X_{t}\right)\right)=\sum_{i=1}^{n} / /_{t}^{-1}\left(\nabla_{U_{t} e_{i}} a\right) d B^{i}+/ /_{t}^{-1}\left(\square a+\nabla_{Z} a\right)\left(X_{t}\right) d t \tag{2.3.5}
\end{equation*}
$$

In short terms, the last two equations may be written as

$$
d\left(/ / t^{-1} a\left(X_{t}\right)\right)=/ /_{t}^{-1} \nabla_{\circ d X} a
$$

resp.

$$
d\left(/ /{ }_{t}^{-1} a\left(X_{t}\right)\right)=/ /_{t}^{-1} \nabla_{d X} a+/ /_{t}^{-1}(\square a)\left(X_{t}\right) d t
$$

Proof. We have $/{ }_{t}^{-1} a\left(X_{t}\right)=U_{0} U_{t}^{-1} a\left(X_{t}\right)=U_{0} F_{a}\left(U_{t}\right)$. It is easily checked that

$$
\bar{Z} F_{a}=F_{\nabla_{z} a} .
$$

Thus, we obtain from Eq. (2.3.2)

$$
\begin{aligned}
d F_{a}(U) & =\sum_{i=1}^{n}\left(H_{i} F_{a}\right)(U) d B^{i}+\left(\Delta^{\mathrm{hor}} F_{a}+\bar{Z} F_{a}\right)(U) d t \\
& =\sum_{i=1}^{n}\left(F_{\nabla_{U e_{i}} a}\right)(U) d B^{i}+\left(F_{\square a}+F_{\nabla_{Z} a}\right)(U) d t \\
& =\sum_{i=1}^{n} U^{-1}\left(\nabla_{U e_{i}} a\right)(X) d B^{i}+U^{-1}\left(\square a+\nabla_{Z} a\right)(X) d t
\end{aligned}
$$

Corollary 2.3.3. Fix $T>0$ and let $a_{t} \in \Gamma(E)$ solve the equation

$$
\frac{\partial}{\partial t} a_{t}=\square a_{t}+\nabla_{Z} a_{t} \quad \text { on }[0, T] \times M
$$

Then

$$
N_{t}:=/ /_{t}^{-1} a_{T-t}\left(X_{t}\right), \quad 0 \leq t \leq T
$$

is a local martingale.
Proof. Indeed we have

$$
d N_{t} \stackrel{\mathrm{~m}}{=} / /_{t}^{-1} \underbrace{\left(\square a_{T-t}+\nabla_{Z} a_{t}+\frac{\partial}{\partial t} a_{T-t}\right)}_{=0}\left(X_{t}\right) d t=0
$$

where $\stackrel{m}{=}$ denotes equality modulo differentials of local martingales.
We are now going to look at operators $L^{\mathscr{R}}$ on $\Gamma(E)$ which differ from $\square$ by a zeroorder term, in other words,

$$
\begin{equation*}
\square-L^{\mathscr{R}}=\mathscr{R} \quad \text { where } \mathscr{R} \in \Gamma(\operatorname{End} E) \tag{2.3.6}
\end{equation*}
$$

Thus, by definition, the action $\mathscr{R}_{x}: E_{x} \rightarrow E_{x}$ is linear for each $x \in M$.
Example 2.3.4. A typical example is $E=\Lambda^{p} T^{*} M$ and $A^{p}(M)=\Gamma\left(\Lambda^{p} T^{*} M\right)$ with $p \geq 1$. The de Rham-Hodge Laplacian

$$
\Delta^{(p)}=-\left(d^{*} d+d d^{*}\right): A^{p}(M) \rightarrow A^{p}(M)
$$

then takes the form

$$
\Delta^{(p)} \alpha=\square \alpha-\mathscr{R} \alpha
$$

where $\mathscr{R}$ is given by the Weitzenböck decomposition. In the special case $p=1$, one obtains $\mathscr{R} \alpha=\operatorname{Ric}\left(\cdot, \alpha^{\sharp}\right)$ where Ric: $T M \oplus T M \rightarrow \mathbb{R}$ is the Ricci tensor.

Definition 2.3.5. Fix $x \in M$ and let $X_{t}$ be a diffusion to $L=\Delta+Z$, starting at $x$. Let $Q_{t}$ be the $\operatorname{Aut}\left(E_{x}\right)$-valued process defined by the following linear pathwise differential equation

$$
\frac{d}{d t} Q_{t}=-Q_{t} \mathscr{R}_{/ / t}, \quad Q_{0}=\operatorname{id}_{E_{x}}
$$

where

$$
\mathscr{R}_{/ / t}:=/ / t^{-1} \circ \mathscr{R}_{X_{t}} \circ / /_{t} \in \operatorname{End}\left(E_{x}\right)
$$

and $/ / t$ is parallel transport in $E$ along $X$.
Proposition 2.3.6. Let $L^{\mathscr{R}}=\square-\mathscr{R}$ be as in Eq. (2.3.6) and $X_{t}$ be a diffusion to $L=\Delta+Z$, starting at $x$. Then, for any $a \in \Gamma(E)$,

$$
d\left(Q_{t} / / /^{-1} a\left(X_{t}\right)\right)=\sum_{i=1}^{n} Q_{t} / /_{t}^{-1}\left(\nabla_{U_{t} e_{i}} a\right) d B_{t}^{i}+Q_{t} / /_{t}^{-1}\left(\square a+\nabla_{Z} a-\mathscr{R} a\right)\left(X_{t}\right) d t
$$

Proof. Let $n_{t}:=/{ }_{t}^{-1} a\left(X_{t}\right)$. Then

$$
\begin{aligned}
d\left(Q_{t} n_{t}\right) & =\left(d Q_{t}\right) n_{t}+Q_{t} d n_{t} \\
& =-Q_{t} / /_{t}^{-1} \mathscr{R}_{X_{t}} / /{ }_{t} n_{t} d t+Q_{t} d n_{t} \\
& =-Q_{t} / /_{t}^{-1}(\mathscr{R} a)\left(X_{t}\right) d t+Q_{t} d n_{t} .
\end{aligned}
$$

The claim thus follows from Proposition 2.3.2.
Corollary 2.3.7. Fix $T>0$ and let $X_{t}(x)$ be a diffusion to $L=\Delta+Z$, starting at $x$. Suppose that $a_{t}$ solves

$$
\left\{\begin{array}{l}
\frac{\partial}{\partial t} a_{t}=\left(\square-\mathscr{R}+\nabla_{Z}\right) a_{t} \quad \text { on }[0, T] \times M, \\
\left.a_{t}\right|_{t=0}=a \in \Gamma(E) .
\end{array}\right.
$$

Then

$$
\begin{equation*}
N_{t}:=Q_{t} / /{ }_{t}^{-1} a_{T-t}\left(X_{t}(x)\right), \quad 0 \leq t \leq T, \tag{2.3.7}
\end{equation*}
$$

is a local martingale, starting at $a_{T}(x)$.
In particular, if Eq. (2.3.7) is a true martingale, we arrive at the formula

$$
a_{T}(x)=\mathbb{E}\left[Q_{T} / /_{T}^{-1} a\left(X_{T}(x)\right)\right], \quad a \in \Gamma(E) .
$$

Proof. Indeed, we have

$$
d N_{t} \stackrel{\mathrm{~m}}{=} Q_{t} / / /_{t}^{-1} \underbrace{\left(\left(\square+\nabla_{Z}-\mathscr{R}\right) a_{T-t}+\frac{\partial}{\partial t} a_{T-t}\right)}_{=0}\left(X_{t}\right) d t=0
$$

Remark 2.3.8. Note that

$$
\frac{d}{d t} Q_{t}=-Q_{t} \mathscr{R}_{/ / t}, \quad \text { with } Q_{0}=\operatorname{id}_{E_{x}}
$$

implies the obvious estimate

$$
\left\|Q_{t}\right\|_{\mathrm{op}} \leq \exp \left(-\int_{0}^{t} \underline{\mathscr{R}}\left(X_{s}(x)\right) d s\right)
$$

where $\underline{\mathscr{R}}(x)=\inf \left\{\left\langle\mathscr{R}_{x} v, w\right\rangle: v, w \in E_{x},\|v\|=\|w\|=1\right\}$.
2.3.3. A formula for the differential. In the sequel, we consider the special case $E=T^{*} M$. Thus $\Gamma(E)$ is the space of differential 1-forms on $M$. We identify vector fields $V \in \Gamma(T M)$ and 1-forms $\alpha \in \Gamma\left(T^{*} M\right)$ via the metric:

$$
V \longleftrightarrow V^{b}, \quad \alpha \longleftrightarrow \alpha^{\#}
$$

hence our results immediately apply to vector fields as well.
DEFINITION 2.3.9. Let $Z \in \Gamma(T M)$ be a vector field on $M$.
(1) The divergence of $Z$ is denoted by

$$
\operatorname{div} Z \in \mathcal{S}^{\infty}(M), \quad \operatorname{div} Z=\operatorname{trace}\left(v \mapsto \nabla_{v} Z\right)
$$

Thus, $(\operatorname{div} Z)(x)=\sum_{i}\left\langle\nabla_{e_{i}} X, e_{i}\right\rangle$ where $\left(e_{i}\right)$ is an orthonormal base for $T_{x} M$.
(2) The adjoint $Z^{*}$ of $Z$ is given by the relation

$$
Z^{*} f=-Z f-(\operatorname{div} Z) f, \quad f \in C^{\infty}(M)
$$

We have the identities:

$$
\begin{aligned}
\langle Z, \nabla f\rangle_{L^{2}(T M)} & =-\langle\operatorname{div} Z, f\rangle_{L^{2}(M)}, \quad \text { if } f \text { is compactly supported, } \\
\langle Z f, h\rangle_{L^{2}(M)} & =\left\langle f, Z^{*} h\right\rangle_{L^{2}(M)}, \quad \text { if } f \text { or } h \text { is compactly supported. }
\end{aligned}
$$

(3) Finally,

$$
\operatorname{Ric}_{Z}(X, Y):=\operatorname{Ric}(X, Y)-\left\langle\nabla_{X} Z, Y\right\rangle, \quad X, Y \in \Gamma(T M)
$$

Similarly, for $\alpha \in \Gamma\left(T^{*} M\right)$, let

$$
(\operatorname{div} \alpha)(x)=\operatorname{trace}\left(T_{x} M \xrightarrow{\nabla \alpha} T_{x}^{*} M \xrightarrow{\#} T_{x} M\right) .
$$

Thus $\operatorname{div} Y=\operatorname{div} Y^{b}$ and $\operatorname{div} \alpha=\operatorname{div} \alpha^{\#}$. In the same way, let $\operatorname{Ric}_{Z}(\alpha):=\operatorname{Ric}_{Z}\left(\cdot, \alpha^{\sharp}\right)$ for any $\alpha \in \Gamma\left(T^{*} M\right)$.

Lemma 2.3.10 (Commutation rules). Let $Z \in \Gamma(T M)$.
(1) For the exterior differential $d$, we have

$$
d(\Delta+Z)=\left(\square-\operatorname{Ric}_{Z}+\nabla_{Z}\right) d
$$

(2) For the codifferential $d^{*}=-\operatorname{div}$, we have

$$
\left(\Delta+Z^{*}\right) \operatorname{div}=\operatorname{div}\left(\square-\operatorname{Ric}_{Z}^{*}+\nabla_{Z}^{*}\right),
$$

where the formal adjoint of $\nabla_{Z}$ (acting on 1-forms) is $\nabla_{Z}^{*} \alpha=-\nabla_{Z} \alpha-(\operatorname{div} Z) \alpha$.
Proof. Indeed, we have

$$
\begin{aligned}
d(\Delta+Z) f & =d\left(-d^{*} d f+(d f) Z\right) \\
& =\Delta^{(1)} d f+\nabla_{Z} d f+\langle\nabla \cdot Z, \nabla f\rangle \\
& =\left(\square+\nabla_{Z}\right)(d f)-\operatorname{Ric}_{Z}(\cdot, \nabla f) \\
& =\left(\square-\operatorname{Ric}_{Z}+\nabla_{Z}\right)(d f) .
\end{aligned}
$$

The formula in (2) is just dual to (1).
Now, let $X=X(x)$ be a diffusion to $\Delta+Z$ on $M$, starting at $X_{0}(x)=x, U$ be a horizontal lift of $X$ to $\mathrm{O}(T M)$, and $B=U_{0} \int_{U} \vartheta$ (taking values in $T_{x} M$ ) the martingale part of the anti-development of $X$.

By Itô's formula, we have

$$
\begin{aligned}
& d(f \circ X)=\sum_{i=1}^{n} d f\left(U e_{i}\right) U_{0}^{-1} d B^{i}+(\Delta f+Z f)(X) d t \\
& d(F \circ U)=\sum_{i=1}^{n}\left(H_{i} F\right)(U) U_{0}^{-1} d B^{i}+\left(\Delta^{\mathrm{hor}} F+\bar{Z} F\right)(U) d t
\end{aligned}
$$

for $f \in C^{2}(M)$, resp. $F \in C^{2}(\mathbf{O}(T M))$.
Proposition 2.3.11. Let $Q_{t}$ be the $\operatorname{Aut}\left(T_{x}^{*} M\right)$-valued process defined by

$$
\frac{d}{d t} Q_{t}=-Q_{t}\left(\operatorname{Ric}_{Z}\right)_{/ / t}, \quad Q_{0}=\operatorname{id}_{T_{x}^{*} M}
$$

Let

$$
P_{t} f(x)=\mathbb{E}\left[f\left(X_{t}(x)\right], \quad f \in C_{b}^{\infty}(M)\right.
$$

be the semigroup generated by $\Delta+Z$, acting on bounded functions on $M$.
Fix $T>0$. Then,

$$
\begin{equation*}
Q_{t} / /_{t}^{-1}\left(d P_{T-t} f\right)_{X_{t}(x)}, \quad 0 \leq t \leq T \tag{2.3.8}
\end{equation*}
$$

is a local martingale in $T_{x}^{*} M$, starting at $\left(d P_{T} f\right)_{x}$.
From the fact that $N_{t}:=Q_{t} / /_{t}^{-1}\left(d P_{T-t} f\right)_{X_{t}(x)}$ is a local martingale, we deduce that, for any adapted process $\ell$ with paths in the Cameron-Martin space $\mathbb{H}\left(\mathbb{R}_{+} ; T_{x} M\right)$, the process

$$
\begin{aligned}
n_{t} & =\left(N_{t}, \ell_{t}\right)-\int_{0}^{t}\left(N_{s}, d \ell_{s}\right) \\
& =\left(\left(d P_{T-t} f\right)_{X_{t}(x)}, / / t Q_{t}^{\mathrm{tr}} \ell_{t}\right)-\int_{0}^{t}\left(\left(d P_{T-s} f\right)_{X_{s}(x)}, / / s Q_{s}^{\mathrm{tr}} \dot{\ell}_{s}\right) d s
\end{aligned}
$$

is a local martingale as well; here $(\cdot, \cdot)$ denotes the pairing between $T_{x}^{*} M$ and $T_{x} M$. Interpreting the last term on the right as quadratic covariation of two (local) martingales, we see that also

$$
\tilde{n}_{t}=\left(d P_{T-t} f\right)_{X_{t}(x)} / /_{t} Q_{t}^{\mathrm{tr}} \ell_{t}-\int_{0}^{t}\left(\left(d P_{T-s} f\right)_{X_{s}(x)}, / / s d B_{s}\right) \int_{0}^{t}\left\langle Q_{t}^{\mathrm{tr}} \dot{\ell}_{s}, d B_{s}\right\rangle
$$

is a local martingale. However taking into account that

$$
\left(P_{T-t} f\right)\left(X_{t}(x)\right)=\int_{0}^{t}\left(d P_{T-s} f\right)_{X_{s}(x)} / /_{s} d B_{s}
$$

we finally see that

$$
\begin{equation*}
\tilde{n}_{t}=\left(d P_{T-t} f\right)_{X_{t}(x)} / / t Q_{t}^{\mathrm{tr}} \ell_{t}-\left(P_{T-t} f\right)\left(X_{t}(x)\right) \int_{0}^{t}\left\langle Q_{s}^{\mathrm{tr}} \dot{\ell}_{s}, d B_{s}\right\rangle, \quad 0 \leq t \leq T \tag{2.3.9}
\end{equation*}
$$

is a local martingale. The idea now is to choose $\ell_{t}$ such that first the local martingale $\tilde{n}_{t}$ is a true martingale, and secondly such that

$$
\ell_{0}=v \quad \text { and } \quad \ell_{T}=0
$$

Then, the equality $\mathbb{E}\left[\tilde{n}_{0}\right]=\mathbb{E}\left[\tilde{n}_{T}\right]$ gives the formula

$$
\left(d P_{T} f\right)_{x} v=\mathbb{E}\left[f\left(X_{T}(x)\right) \int_{0}^{T}\left\langle Q_{s}^{\mathrm{tr}} \dot{\ell}_{s}, d B_{s}\right\rangle\right]
$$

This gives the following result.
THEOREM 2.3.12 (Gradient formula). Let $u$ be the solution to the heat equation

$$
\frac{\partial}{\partial t} u=\frac{1}{2} \Delta u,\left.\quad u\right|_{t=0}=f
$$

Let $X=X(x)$ be the diffusion to $\Delta+Z$ starting from $x$, which is assumed to be nonexplosive. Then, for $v \in T_{x} M$,

$$
\left(d P_{T} f\right)_{x} v=-\mathbb{E}\left[f\left(X_{T}(x)\right) \int_{0}^{\tau \wedge T}\left\langle\Theta_{s} \dot{\ell}_{s}, d B_{s}\right\rangle\right]
$$

where

- $\Theta_{t}$ is the $\operatorname{Aut}\left(T_{x} M\right)$-valued process defined by

$$
\frac{d}{d t} \Theta_{t}=-\operatorname{Ric} / /{ }_{t}\left(\cdot, \Theta_{t}\right)^{\sharp}+(\nabla . Z)_{/ / t} \Theta_{t}, \quad \Theta_{0}=\operatorname{id}_{T_{x} M}
$$

- $\tau=\tau_{D}(x)$ is the first exit time of $X$ from some relatively compact neighbourhood $D$ of $x$;
- $B$ is a Brownian motion in $T_{x} M$;
- $\ell_{t}$ is any adapted process in $T_{x} M$ with absolutely continuous paths such that (for some $\varepsilon>0$ )

$$
\ell_{0}=v, \quad \ell_{\tau}=0 \quad \text { and } \quad\left(\int_{0}^{\tau \wedge T}\left|\dot{\ell}_{t}\right|^{2} d t\right)^{1 / 2} \in L^{1+\varepsilon}
$$

2.3.4. A formula for the codifferential. Let now $X=X(x)$ be a diffusion to $\Delta-Z$ on $M$, starting at $X_{0}(x)=x, U$ be a horizontal lift of $X$ to $\mathrm{O}(T M)$, and $B=U_{0} \int_{U} \vartheta$ (taking values in $T_{x} M$ ) the martingale part of the anti-development of $X$.

By Itô's formula, we have

$$
\begin{aligned}
& d(f \circ X)=\sum_{i=1}^{n} d f\left(U e_{i}\right) U_{0}^{-1} d B^{i}+(\Delta f-Z f)(X) d t \\
& d(F \circ U)=\sum_{i=1}^{n}\left(H_{i} F\right)(U) U_{0}^{-1} d B^{i}+\left(\Delta^{\mathrm{hor}} F-\bar{Z} F\right)(U) d t
\end{aligned}
$$

for $f \in C^{2}(M)$, resp. $F \in C^{2}(\mathrm{O}(T M))$.
Proposition 2.3.13. Let $X=X(x)$ be a diffusion to $\Delta-Z$ on $M$, starting at $X_{0}(x)=x$, which for simplicity is assumed to be non-explosive. In terms of

$$
\mathscr{R}_{Z}:=\operatorname{Ric}_{Z}^{*}+\operatorname{div} Z \in \operatorname{End}\left(T^{*} M\right)
$$

where $\operatorname{div} Z$ acts fiberwise as multiplication operator, let $Q_{t}$ be the $\operatorname{Aut}\left(T_{x}^{*} M\right)$-valued process defined by

$$
\frac{d}{d t} Q_{t}=-Q_{t}\left(\mathscr{R}_{Z}\right)_{/ / t}, \quad Q_{0}=\mathrm{id}_{T_{x}^{*} M}
$$

Fix $T>0$ and let $\alpha_{t}$ solve

$$
\left\{\begin{array}{l}
\frac{\partial}{\partial t} \alpha_{t}=\left(\square-\mathscr{R}_{Z}+\nabla_{Z}\right) \alpha_{t} \quad \text { on }[0, T] \times M \\
\left.\alpha_{t}\right|_{t=0}=\alpha \in \Gamma\left(T^{*} M\right) .
\end{array}\right.
$$

Then, for given $x \in M$,

$$
\begin{equation*}
n_{t}:=\left(\operatorname{div} \alpha_{T-t}\right)\left(X_{t}(x)\right) \exp \left(\int_{0}^{t}(\operatorname{div} Z)\left(X_{s}(x)\right) d s\right), \quad 0 \leq t \leq T \tag{2.3.10}
\end{equation*}
$$

is a real local martingale, starting at $\left(\operatorname{div} \alpha_{T}\right)(x)$.
Let $\ell$ be an adapted process with paths in $\mathbb{H}\left(\mathbb{R}_{+} ; \mathbb{R}\right)$ which takes the form

$$
\begin{equation*}
\ell_{t}=\exp \left(-\int_{0}^{t}(\operatorname{div} Z)\left(X_{s}(x)\right) d s\right) h_{t} \tag{2.3.11}
\end{equation*}
$$

where $h_{0}=0$ and $h_{T}=1$. Starting with the local martingale $n_{t}$ given by (2.3.10), we observe that also

$$
\begin{equation*}
\tilde{n}_{t}=n_{t} \ell_{t}-\int_{0}^{t} n_{s} d \ell_{s} \equiv \int_{0}^{t} \ell_{s} d n_{s}+n_{0} \ell_{0}, \quad 0 \leq t \leq T \tag{2.3.12}
\end{equation*}
$$

is a local martingale. In explicit terms it writes as

$$
\begin{equation*}
\tilde{n}_{t}=\left(\operatorname{div} \alpha_{T-t}\right)\left(X_{t}(x)\right) h_{t}+\int_{0}^{t}\left(\operatorname{div} \alpha_{T-s}\right)\left(X_{s}(x)\right)\left[(\operatorname{div} Z)\left(X_{s}(x)\right) h_{s}-\dot{h}_{s}\right] d s \tag{2.3.13}
\end{equation*}
$$

Assuming that (2.3.12) is actually a true martingale, by taking expectations, we get from $\mathbb{E}\left[\tilde{n}_{T}\right]=0$ the formula

$$
\begin{equation*}
\mathbb{E}\left[(\operatorname{div} \alpha)\left(X_{T}(x)\right)\right]=\mathbb{E}\left[\int_{0}^{T} n_{t} d \ell_{t}\right] \tag{2.3.14}
\end{equation*}
$$

It remains to evaluate the right-hand-side of Eq. (2.3.14). In this way we obtain

$$
\left.\begin{array}{rl}
n_{t} d \ell_{t} & =-\left(\operatorname{div} \alpha_{T-t}\right)\left(X_{t}(x)\right)\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right] d t \\
& =-\sum_{i=1}^{n}\left\langle\nabla_{U_{t} e_{i}} \alpha_{T-t}, U_{t} e_{i}\right\rangle\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right] d t \\
& =-\sum_{i=1}^{n}\left\langle/ / /_{t}^{-1} \nabla_{U_{t} e_{i}} \alpha_{T-t}, U_{0} e_{i}\right\rangle\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right] d t \\
& =-\sum_{i=1}^{n}\left\langle Q_{t} / / t\right.
\end{array}{ }^{-1} \nabla_{U_{t} e_{i}} \alpha_{T-t},\left(Q_{t}^{-1}\right)^{*} U_{0} e_{i}\right\rangle\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right] d t . ~ \$
$$

Using the fact that $N_{t}:=Q_{t} / /_{t}^{-1} a_{T-t}\left(X_{t}(x)\right)$ is a local martingale, indeed

$$
d\left(Q_{t} / /_{t}^{-1} a_{T-t}\left(X_{t}(x)\right)\right)=\sum_{i=1}^{n} Q_{t} / /_{t}^{-1}\left(\nabla_{U_{t} e_{i}} a_{T-t}\right) d B_{t}^{i}
$$

we get

$$
\begin{aligned}
& \int_{0}^{T} n_{t} d \ell_{t} \\
& =-\left\langle\sum_{i=1}^{n} \int_{0}^{T} Q_{t} / /_{t}^{-1}\left(\nabla_{U_{t} e_{i}} a_{T-t}\right) d B_{t}^{i}, \int_{0}^{T}\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right]\left(Q_{t}^{-1}\right)^{*} d B_{t}\right\rangle \\
& =-\left\langle N_{T}, \int_{0}^{T}\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right]\left(Q_{t}^{-1}\right)^{*} d B_{t}\right\rangle
\end{aligned}
$$

$$
=-\left\langle Q_{T} / /_{T}^{-1} \alpha\left(X_{T}\right), \int_{0}^{T}\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right]\left(Q_{t}^{-1}\right)^{*} d B_{t}\right\rangle
$$

Using the identification of differential forms and vector fields via the metric, we obtain the following result.

THEOREM 2.3.14. Let $M$ be a Riemannian manifold and $Z$ a smooth vector field on $M$. Let $X=X(x)$ be a diffusion to $\Delta-Z$ on $M$, starting at $X_{0}(x)=x$, which is assumed to be non-explosive. Let $T>0$ and $h$ be an adapted process with paths in $\mathbb{H}([0, T] ; \mathbb{R})$ such that $h_{0}=0$ and $h_{T}=1$, and such that (2.3.13) is a true martingale. Then for all smooth vector fields $V$ on $M$,

$$
\begin{aligned}
& \mathbb{E}\left[(\operatorname{div} V)\left(X_{T}(x)\right)\right] \\
& \quad=-\mathbb{E}\left[\left\langle Q_{T} / /_{T}^{-1} V\left(X_{T}(x)\right), \int_{0}^{T}\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right]\left(Q_{t}^{-1}\right)^{*} d B_{t}\right\rangle\right]
\end{aligned}
$$

where $Q$ is the $\operatorname{Aut}\left(T_{x} M\right)$-valued process defined by the following pathwise differential equation:

$$
\frac{d}{d t} Q_{t}=-Q_{t}\left(\mathscr{R}_{Z}\right)_{/ / t} \quad \text { with } Q_{0}=\operatorname{id}_{T_{x} M}
$$

COROLLARY 2.3.15. We keep notations and assumptions of Theorem 2.3.14. Using the relation

$$
\operatorname{div}(f V)=V f+f \operatorname{div} V, \quad f \in C^{\infty}(M), V \in \Gamma(T M)
$$

we get the formula

$$
\begin{align*}
& \mathbb{E}\left[(V f)\left(X_{T}(x)\right)\right]=-\mathbb{E}\left[f\left(X_{T}(x)\right)(\operatorname{div} V)\left(X_{T}(x)\right)\right]  \tag{2.3.15}\\
& -\mathbb{E}\left[f\left(X_{T}(x)\right)\left\langle Q_{T} / /_{T}^{-1} V\left(X_{T}(x)\right), \int_{0}^{T}\left[(\operatorname{div} Z)\left(X_{t}(x)\right) h_{t}-\dot{h}_{t}\right]\left(Q_{t}^{-1}\right)^{*} d B_{t}\right\rangle\right]
\end{align*}
$$

where the right-hand side of Eq. (2.3.15) does not contain any derivatives of $f$.
2.3.5. General remarks. Our approach is based on martingale arguments; integration by parts is done at the level of local martingales. We get the wanted formulae then by taking expectations, under conditions which assure the local martingales to be true martingales.

REMARK 2.3.16. The formulae allow the choice of a finite energy process $\left(\ell_{t}\right)_{t \in[0, T]}$. Depending on the type of the wanted formulae, conditions are imposed on the left endpoint $\ell_{0}$, or the right endpoint $\ell_{t}$.
(i) The argument leading to the gradient formula of Theorem 2.3.12 is based on the fact that the local martingale (2.3.9) is a true martingale. Since the condition on $\ell_{t}$ is imposed on the left endpoint, this can always be achieved, for instance, by taking $\ell_{s}=0$ for $s \geq \tau \wedge T$ where $\tau$ is the first exit time of some relatively compact neighbourhood of $x$. No bounds on the geometry are needed; even explosion in finite times of the underlying diffusion can be allowed.
(ii) Imposing however the conditions $\ell_{0}=0$ and $\ell_{T}=v$ in (2.3.9) would lead to a formula for

$$
\mathbb{E}\left[(d f)_{X_{T}(x)} / /_{T} Q_{T}^{\operatorname{tr}} v\right]
$$

which clearly requires strong assumptions. Note that if the local martingale (2.3.8) is a true martingale, we get the formula

$$
\left(d P_{T} f\right)_{x}=\mathbb{E}\left[Q_{T} / /_{T}^{-1}(d f)_{X_{T}(x)}\right]
$$

For such a formula to hold, clearly $X_{t}(x)$ needs to be non-explosive.
(iii) For the formula in Theorem 2.3.14 the conditions on $\ell_{t}$, resp. $h_{t}$, given by Eq. (2.3.11) are on the right-hand-side. Thus assumptions are needed to make (2.3.13) a true martingale.

### 2.4. Brownian Bridges and Brownian Loops

Let $(M, g)$ be a complete connected Riemannian manifold of dimension $n$. Given $x \in M$ we denote by $X$. $(x)$ Brownian motion on $M$ starting from $x$ at time zero. Now fix $t>0$. Given $x, y \in M$ we want to define a Brownian bridge $X^{t}(x, y)$ from $x$ to $y$ of lifetime $t$. Intuitively, a Brownian bridge $X^{t}(x, y)$ is a Brownian motion on $M$ starting from $x$ at time 0 but "conditioned to hit $y$ at time $t$ ". In particular, for $s<t$ we require the properties:

$$
\begin{aligned}
\mathbb{P}\left\{X_{s}(x) \in A\right\} & =\int 1_{A}(z) \mathbb{P}\left\{X_{s}(x) \in d z\right\}=\int 1_{A}(z) p(s, x, z) \operatorname{vol}(d z) \\
\mathbb{P}\left\{X_{s}^{t}(x, y) \in A\right\} & =\int 1_{A}(z) \frac{p(s, x, z) p(t-s, z, y)}{p(t, x, y)} \operatorname{vol}(d z)
\end{aligned}
$$

where $p(t, x, y)$ denotes the heat kernel on $M$.
Lemma 2.4.1 (Local heat kernel estimates). Let $M$ be complete.
(1) For each point $x_{0} \in M$ and each $R>0$, there exist constants $c_{i}, k_{i}>0($ for $i=1,2)$ such that

$$
k_{1} s^{-n / 2} \exp \left(-c_{1} \frac{d(x, y)^{2}}{s}\right) \leq p(s, x, y) \leq k_{2} s^{-n / 2} \exp \left(-c_{2} \frac{d(x, y)^{2}}{s}\right)
$$

for all $(s, x, y) \in] 0, R\left[\times B\left(x_{0}, R\right) \times B\left(x_{0}, R\right)\right.$.
(2) For each point $x_{0} \in M$ and each $R>0$, there exist constant $c>0$ such that

$$
|(\nabla \log p(s, \cdot, y))(x)| \leq c\left(\frac{d(x, y)}{s}+\frac{1}{\sqrt{s}}\right)
$$

for all $(s, x, y) \in] 0, R\left[\times B\left(x_{0}, R\right) \times B\left(x_{0}, R\right)\right.$.
Let $\left(C\left(\mathbb{R}_{+}, \mathbb{R}^{n}\right), \mathscr{F}, \mathbb{P}\right)$ be the standard Wiener space and $\left(\mathscr{F}_{s}\right)_{s \geq 0}$ be the standard filtration (with the usual conditions). Denote by

$$
B_{s}:=\operatorname{pr}_{s}: C\left(\mathbb{R}_{+}, \mathbb{R}^{n}\right) \rightarrow \mathbb{R}^{n}, \quad s \geq 0
$$

the standard Wiener process. Construct Brownian motion $X$ on $(M, g)$ starting from $x$ in the usual way as $X=\pi \circ U$, where

$$
\left\{\begin{array}{l}
d U=\sum_{i=1}^{n} L_{i}(U) \circ d B^{i} \\
U_{0}=u \in \pi^{-1} x \in \mathrm{O}(T M)
\end{array}\right.
$$

Let

$$
h:[0, t[\times M \rightarrow] 0, \infty[
$$

be a smooth space-time harmonic function, in the sense that

$$
\left(\frac{\partial}{\partial s}+\frac{1}{2} \Delta\right) h \equiv 0
$$

In the sequel we always take

$$
\begin{equation*}
h(s, \cdot):=\frac{p(t-s, \cdot, y)}{p(t, x, y)} . \tag{2.4.1}
\end{equation*}
$$

Furthermore let $\mathscr{F}_{s}^{0}:=\sigma\left\{\operatorname{pr}_{r}: r \leq s\right\} \subset \mathscr{F}_{s}$.
Lemma 2.4.2. Let $h(s, \cdot)$ be given by (2.4.1). Then

$$
Z_{s}:=h\left(s, X_{s}(x)\right), \quad 0 \leq s<t
$$

is a martingale and by Kolmogorov's theorem there exists a unique probability measure $\mathbb{P}^{h}$ on $C\left(\left[0, t\left[, \mathbb{R}^{n}\right), \sigma\left\{\operatorname{pr}_{s}: s<t\right\}\right.\right.$ such that

$$
\left.\frac{d \mathbb{P}^{h}}{d \mathbb{P}^{2}}\right|_{\mathscr{F}_{s}^{0}}=Z_{s} \quad \mathbb{P} \text {-a.s. } \quad(0 \leq s<t) .
$$

Corollary 2.4.3. Let

$$
h(s, \cdot):=\frac{p(t-s, \cdot, y)}{p(t, x, y)},
$$

and let $\mathbb{P}_{x, y}^{t}:=\mathbb{P}^{h}$ on $C\left(\left[0, t\left[, \mathbb{R}^{n}\right)\right.\right.$ as constructed above. The following properties hold:
(1) For each $s<t$,

$$
\mathbb{P}_{x, y}^{t}\left\{X_{s}(x) \in A\right\}=\mathbb{E}\left[\left(1_{A} \circ X_{s}(x)\right) Z_{s}\right]=\mathbb{E}\left[1_{A} \circ X_{s}(x) \frac{p\left(t-s, X_{s}(x), y\right)}{p(t, x, y)}\right]
$$

(2) The process $\left.X(x)\right|_{[0, t[ }$ is a $\mathbb{P}_{x, y}^{t}$-semimartingale (with respect to the local completion of $\left(\mathscr{F}_{s}^{0}\right)_{0 \leq s<t}$ with respect to $\left.\mathbb{P}_{x, y}^{t}\right)$.
(3) The process

$$
\left(\frac{1}{h\left(s, X_{s}(x)\right)}\right)_{0 \leq s<t}
$$

is a positive $\mathbb{P}_{x, y}^{t}$-martingale, and hence it has $a \mathbb{P}_{x, y}^{t}$-a.s. limit as $s \uparrow t$.
(4) Since for $x \neq y$,

$$
p(t-s, x, y) \rightarrow 0 \text { as } s \uparrow t
$$

we observe that $X_{s}(x) \rightarrow y \mathbb{P}_{x, y}^{t}$-a.s. as $s \uparrow t$.
Notation 2.4.4. The process $\left(X_{s}(x)\right)_{0 \leq s \leq t}$ with respect to $\mathbb{P}_{x, y}^{t}$ is called Brownian bridge (pinned Brownian motion) on $M$ from $x$ to $y$ with lifetime $t$. The distribution $\mu_{x, y}^{t}$ of $\left.X(x)\right|_{[0, t]}$ on

$$
\{\omega \in C([0, t], M): w(0)=x, w(t)=y\}
$$

is called Brownian bridge measure.
PROPOSITION 2.4.5.
(1) By construction each Brownian bridge from $x$ to $y$ with lifetime $t$ is a continuous Markov process $\left(Y_{s}\right)_{0 \leq s \leq t}$ such that

$$
\mathbb{P}\left\{Y_{s} \in d z \mid Y_{0}=x\right\}=\frac{p(s, x, z) p(t-s, z, y)}{p(t, x, y)} \operatorname{vol}(d z)
$$

Hence the time-reversed Markov process $\widehat{Y}_{s}:=Y_{t-s}$ is a bridge from from y to $x$ in the sense that its distribution is given by $\mu_{y, x}^{t}$.
(2) Since $\left(X_{s}(x)\right)_{s<t}$ is a $\mathbb{P}_{x, y}^{t}$-semimartingale on $M$, it is the stochastic development of an $\mathbb{R}^{n}$-valued semimartingale. Indeed, it is the stochastic development of the $\mathbb{P}$ Brownian motion $B$ on $\mathbb{R}^{n}$ considered as $\mathbb{P}_{x, y}^{t}$-semimartingale.

This leads to the problem of finding the Doob-Meyer decomposition of $\left(B_{s}\right)_{s<t}$ with respect to $\mathbb{P}_{x, y}^{t}$. First, recall that

$$
Z_{s}=h\left(s, X_{s}(x)\right), \quad 0 \leq s<t
$$

is a $\mathbb{P}$-martingale. Thus by Itô's formula,

$$
d Z=\sum_{i}\left\langle\nabla h\left(s, X_{s}(x)\right), U_{s} e_{i}\right\rangle_{T_{X_{s}} M} d B^{i}
$$

where

$$
\nabla h(s, z)=(\nabla h(s, \cdot))_{z}
$$

Hence, we have

$$
d Z=Z\langle c, d B\rangle_{\mathbb{R}^{n}}
$$

where

$$
c_{s}^{i}=\frac{\left\langle\nabla h\left(s, X_{s}(x)\right), U e_{i}\right\rangle}{h\left(s, X_{s}(x)\right)}
$$

or equivalently

$$
c_{s}=\frac{U_{s}^{-1} \nabla h\left(s, \pi\left(U_{s}\right)\right)}{h\left(s, \pi\left(U_{s}\right)\right)}
$$

Thus by Girsanov's theorem we find that

$$
\tilde{B}_{s}:=B_{s}-\int_{0}^{s} c_{r} d r
$$

is a $\mathbb{P}_{x, y}^{t}-\mathrm{BM}\left(\mathbb{R}^{n}\right)$ on the interval $[0, t[$. This means in particular that

$$
B_{s}=\tilde{B}_{s}+\int_{0}^{s} U_{r}^{-1} \nabla \log h\left(r, X_{r}(x)\right) d r
$$

is the $\mathbb{P}_{x, y}^{t}$-Doob-Meyer decomposition of $\left(B_{s}\right)_{s<t}$. Note that on the right-hand-side $\tilde{B}_{s}$ has a limit $\mathbb{P}_{x, y}^{t}$-a.s., as $s \uparrow t$, and the integral converges absolutely $\mathbb{P}_{x, y}^{t}$-a.s. as $s \uparrow t$. Therefore $B$ is a $\mathbb{P}_{x, y}^{t}$-semimartingale on $[0, t]$.

COROLLARY 2.4.6.
a) The process $X(x)$ is a $\mathbb{P}_{x, y}^{t}$-semimartingale on $[0, t]$ with anti-development given by

$$
\mathscr{A}(X(x))=\mathrm{BM}\left(T_{x} M\right)+\int_{0}^{\cdot} / /_{0, r}^{-1}(\nabla \log h)\left(r, X_{r}(x)\right) d r
$$

where

$$
(\nabla \log h)(r, \cdot)=\frac{\nabla p(t-r, \cdot, y)}{p(t-r, \cdot, y)}
$$

b) With respect to the Brownian bridge measure $\mathbb{P}_{x, y}^{t}$ we have

$$
\begin{aligned}
d U & =\sum_{i} L_{i}(U) \circ d B^{i} \\
& =\sum_{i} L_{i}(U) \circ\left(d \tilde{B}^{i}+c_{s}^{i} d s\right) \\
& =\sum_{i} L_{i}(U) \circ d \tilde{B}^{i}+L_{0}\left(s, U_{s}\right) d s, \quad \pi U_{0}=x
\end{aligned}
$$

where $L_{0}(s, \cdot) \in \Gamma(T \mathrm{O}(T M))$ denotes the horizontal lift of the (time-dependent) vector field
$\nabla \log p(t-s, \cdot, y) \in \Gamma(T M), \quad 0 \leq s<t$,
i.e., $L_{0}(s, u)=h_{u}(\nabla \log p(t-s, \pi(u), y)) \in T_{u} \mathrm{O}(T M)$.

REMARK 2.4.7. Let $\left(U_{s}\right)_{0 \leq s<t}$ be a solution of

$$
d U=\sum_{i} L_{i}(U) \circ d W^{i}+L_{0}\left(s, U_{s}\right) d s, \quad U_{0}=u \in \pi^{-1}\{x\}
$$

where $W$ is a Brownian motion on $\mathbb{R}^{n}$, the limit $U_{t}:=\lim _{s \uparrow t} U_{s}$ exists a.s. and

$$
X_{s}:=\pi \circ U_{s}, \quad 0 \leq s \leq t
$$

is Brownian bridge on $M$ from $x$ to $y$ of lifetime $t$ (in particular, $X_{s}$ is a semimartingale on the interval [ $0, t$ ] with the property $X_{s} \rightarrow y$ for $s \uparrow t$ almost surely). On $M$ we have

$$
d X_{s}=U_{s} \circ d W_{s}+\nabla \log p\left(t-s, X_{s}, y\right) d s, \quad X_{0}=x
$$

REMARK 2.4.8. Let $X$ be a Brownian motion on $(M, g)$ starting from $x$ (with respect to $\mathbb{P}$ ) and fix $t>0$.
(1) For all $A \in \sigma\left\{X_{r}: r \leq t\right\}$, we have

$$
\begin{aligned}
\mathbb{P}(A) & =\int \mathbb{P}_{x, y}^{t}(A)\left(P \circ X_{t}^{-1}\right)(d y) \\
& =\int \mathbb{P}_{x, y}^{t}(A) p(t, x, y) \operatorname{vol}(d y)
\end{aligned}
$$

Indeed, if $A \in \sigma\left\{X_{r}: r \leq s\right\}$ where $s<t$, then

$$
\mathbb{P}_{x, y}^{t}(A)=\mathbb{E}\left[1_{A} Z_{s}\right]=\mathbb{E}\left[1_{A} \frac{p\left(t-s, X_{s}, y\right)}{p(t, x, y)}\right]
$$

and
$\int P_{x, y}^{t}(A) p(t, x, y) \operatorname{vol}(d y)=\int \mathbb{E}\left[1_{A} p\left(t-s, X_{s}, y\right)\right] \operatorname{vol}(d y)=\mathbb{E}\left[1_{A}\right]=\mathbb{P}(A)$.
(2) The map

$$
y \mapsto P_{x, y}^{t}
$$

is a desintegration of $\left.\mathbb{P}\right|_{\left\{\sigma\left\{X_{r}: r \leq t\right\}\right.}$ with respect to Brownian motion $X$ starting at $x$ according to its values $y$ at time $t$. In an informal way, we write

$$
\mathbb{P}_{x, y}^{t}\{\ldots\}=\mathbb{P}\left\{\ldots \mid X_{t}(x)=y\right\}
$$

EXAMPLE 2.4.9. Let $M=\mathbb{R}^{n}$ and

$$
d X_{s}=d W_{s}+\frac{y-X_{s}}{t-s} d s, \quad X_{0}=x
$$

Then $X$ is a Brownian bridge from $x$ to $y$ of lifetime $t$.

## APPENDIX A

## Background on SDEs

## A.1. One-dimensional Stochastic Differential Equations

In this Section we collect some facts about stochastic differential equations in the one-dimensional case and develop a qualitative theory of one-dimensional SDEs which provides a useful tool for many geometric comparison theorems. Most of the results of this Section go back to the work of William Feller and have originally been formulated for one-dimensional diffusion processes.

We consider the situation of an Itô SDE of the form

$$
\begin{equation*}
d Y=\beta(t, Y) d t+\sigma(t, Y) d B \tag{A.1.1}
\end{equation*}
$$

with continuous coefficients $\beta, \sigma: \mathbb{R}_{+} \times \mathbb{R} \rightarrow \mathbb{R}$ and a one-dimensional Brownian motion $B$ as driving process. We consider first the case of global solutions of (A.1.1) of infinite lifetime.

THEOREM A.1.1. Let $Y^{1}$ and $Y^{2}$ be two solutions of (A.1.1) with $Y_{0}^{1}=Y_{0}^{2}$. Then also $Y^{1} \vee Y^{2}$ is a solution of (A.1.1) to the same initial condition if and only if the local time $L^{0}\left(Y^{2}-Y^{1}\right)$ of $Y^{2}-Y^{1}$ at 0 vanishes modulo indistinguishability.

Definition A.1.2. For any real $a$ the local time of a continuous real semimartingale $X \in \mathscr{S}$ at $a$ is given by

$$
\begin{equation*}
L_{t}^{a}(X):=\left|X_{t}-a\right|-\left|X_{0}-a\right|-\int_{0}^{t} \operatorname{sign}\left(X_{s}-a\right) d X_{s} \tag{A.1.2}
\end{equation*}
$$

where $\int_{0}^{\infty} \mathbb{1}_{\left\{\left|X_{s}\right| \neq a\right\}} d L_{s}^{a}(X)=0$ almost surely. Recall that sign $:=-\mathbb{1}_{]-\infty, 0]}+\mathbb{1}_{] 0, \infty}[$. It holds that

$$
\begin{align*}
& \left(X_{t}-a\right)_{+}=\left(X_{0}-a\right)_{+}+\int_{0}^{t} 1_{\left\{X_{s}>a\right\}} d X_{s}+\frac{1}{2} L_{t}^{a}(X)  \tag{A.1.3}\\
& \left(X_{t}-a\right)_{-}=\left(X_{0}-a\right)_{-}-\int_{0}^{t} 1_{\left\{X_{s} \leq a\right\}} d X_{s}+\frac{1}{2} L_{t}^{a}(X) \tag{A.1.4}
\end{align*}
$$

These formulae are well-known (e.g. [38, p. 222]) and usually refered to under the name "Tanaka formulae".

Proof (of Theorem A.1.1). Letting $L^{0}\left(Y^{2}-Y^{1}\right)$ denote the local time of $Y^{2}-Y^{1}$ at 0 , we have

$$
\begin{aligned}
d\left(Y^{1} \vee Y^{2}\right)= & d Y^{1}+d\left(Y^{2}-Y^{1}\right)_{+} \\
= & \beta\left(t, Y^{1}\right) d t+\sigma\left(t, Y^{1}\right) d B+1_{\left\{Y^{2}>Y^{1}\right\}} d\left(Y^{2}-Y^{1}\right)+\frac{1}{2} d L^{0}\left(Y^{2}-Y^{1}\right) \\
= & \left(\beta\left(t, Y^{1}\right)+\left(\beta\left(t, Y^{2}\right)-\beta\left(t, Y^{1}\right)\right) 1_{\left\{Y^{2}>Y^{1}\right\}}\right) d t \\
& +\left(\sigma\left(t, Y^{1}\right)+\left(\sigma\left(t, Y^{2}\right)-\sigma\left(t, Y^{1}\right)\right) 1_{\left\{Y^{2}>Y^{1}\right\}}\right) d B+\frac{1}{2} d L^{0}\left(Y^{2}-Y^{1}\right)
\end{aligned}
$$

$$
=\beta\left(t, Y^{1} \vee Y^{2}\right) d t+\sigma\left(t, Y^{1} \vee Y^{2}\right) d B+\frac{1}{2} d L^{0}\left(Y^{2}-Y^{1}\right)
$$

from where the claim is seen.
THEOREM A.1.3. Suppose that for any two solutions $Y^{1}$ and $Y^{2}$ of (A.1.1) such that $Y_{0}^{1}=Y_{0}^{2}$, it holds that $L^{0}\left(Y^{1}-Y^{2}\right)=0$. Then if solutions to (A.1.1) are unique in distribution, they are even pathwise unique.

Proof. Indeed, letting $Y^{1}$ and $Y^{2}$ be two solutions satisfying $Y_{0}^{1}=Y_{0}^{2}$, by Theorem A.1.1 then also $Y^{1} \vee Y^{2}$ is a solution. For $t \geq 0$, by the uniqueness of solutions in law, both $Y_{t}^{1}$ and $Y_{t}^{1} \vee Y_{t}^{2}$ have the same law, with the consequence that $Y_{t}^{1} \geq Y_{t}^{2}$ almost surely; analogously one obtains $Y_{t}^{1} \leq Y_{t}^{2}$ almost surely. The claim thus follows by the continuity of the paths.

The following Lemma gives a criterion for the vanishing of the local time $L^{0}(X)$ of a semimartingales $X \in \mathscr{S}$ at 0 . To this end, we denote by $\rho$ a measurable function $\mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$such that $\int_{0}^{t} \rho(u)^{-1} d u=\infty$ for any $t>0$. We note this property shortly as $\int_{0+} \rho(u)^{-1} d u=\infty$.

Lemma A.1.4. Let $X \in \mathscr{S}$ and $\rho: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$be a measurable function such that $\int_{0+} \rho(u)^{-1} d u=\infty$. Suppose there exists $\varepsilon>0$ such that for any $t>0$

$$
\begin{equation*}
\int_{0}^{t} 1_{\left\{0<X_{s} \leq \varepsilon\right\}} \rho\left(X_{s}\right)^{-1} d[X]_{s}<\infty \quad \text { almost surely } . \tag{A.1.5}
\end{equation*}
$$

Then $L^{0}(X)=0$ modulo indistinguishability.
Proof. From the "occupation times formula" of the local time (e.g. [38, p. 224]) we have for fixed $t>0$ the relation

$$
\int_{0}^{t} 1_{\left\{0<X_{s} \leq \varepsilon\right\}} \rho\left(X_{s}\right)^{-1} d[X]_{s}=\int_{0}^{\varepsilon} \rho(a)^{-1} L_{t}^{a}(X) d a
$$

where $L_{t}^{a}(X)$ denotes again the local time of $X$ at $a$. Hence if $L_{t}^{0}(X)$ does not vanish almost surely, by means of the right continuity of $L_{t}^{a}(X)$ in $a$, the right-hand side of the last formula would be infinite with positive probability - in contradiction to assumption (A.1.5).

THEOREM A. 1.5 (Yamada-Watanabe). Let $\beta, \sigma: \mathbb{R}_{+} \times \mathbb{R} \rightarrow \mathbb{R}$ be continuous functions satisfying the following properties:
(i) $\sigma$ is bounded, and there exists a measurable function $\rho: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$satisfying $\int_{0+} \rho(u)^{-1} d u=\infty$ such that

$$
|\sigma(s, x)-\sigma(s, y)|^{2} \leq \rho(|x-y|)
$$

for all $s \in \mathbb{R}_{+}$and $x, y \in \mathbb{R}$.
(ii) $\beta$ is globally Lipschitz, i.e., for any $t \geq 0$ there is a constant $L_{t}$ such that

$$
|\beta(s, x)-\beta(s, y)| \leq L_{t}|x-y|
$$

for all $0 \leq s \leq t$ and $x, y \in \mathbb{R}$.
Then solutions of (A.1.1) are pathwise unique.

Proof. Let $Y^{1}$ and $Y^{2}$ be two solutions of (A.1.1) satisfying $Y_{0}^{1}=Y_{0}^{2}$. Then by $d\left(Y^{1}-Y^{2}\right)=\left(\beta\left(t, Y^{1}\right)-\beta\left(t, Y^{2}\right)\right) d t+\left(\sigma\left(t, Y^{1}\right)-\sigma\left(t, Y^{2}\right)\right) d B$, we obtain

$$
\begin{aligned}
\int_{0}^{t} \rho\left(Y_{s}^{1}-Y_{s}^{2}\right)^{-1} & 1_{\left\{Y_{s}^{1}>Y_{s}^{2}\right\}} d\left[Y^{1}-Y^{2}\right]_{s} \\
& =\int_{0}^{t} \rho\left(Y_{s}^{1}-Y_{s}^{2}\right)^{-1} 1_{\left\{Y_{s}^{1}>Y_{s}^{2}\right\}}\left(\sigma\left(s, Y_{s}^{1}\right)-\sigma\left(s, Y_{s}^{2}\right)\right)^{2} d s \leq t
\end{aligned}
$$

and by Lemma A.1.4 hence $L^{0}\left(Y^{1}-Y^{2}\right)=0$ modulo indistinguishability. Thus we have

$$
\begin{aligned}
\left|Y_{t}^{1}-Y_{t}^{2}\right|= & \int_{0}^{t} \operatorname{sign}\left(Y_{s}^{1}-Y_{s}^{2}\right) d\left(Y_{s}^{1}-Y_{s}^{2}\right) \\
= & \int_{0}^{t} \operatorname{sign}\left(Y_{s}^{1}-Y_{s}^{2}\right)\left(\beta\left(s, Y_{s}^{1}\right)-\beta\left(s, Y_{s}^{2}\right)\right) d s \\
& +\int_{0}^{t} \operatorname{sign}\left(Y_{s}^{1}-Y_{s}^{2}\right)\left(\sigma\left(s, Y_{s}^{1}\right)-\sigma\left(s, Y_{s}^{2}\right)\right) d B_{s}
\end{aligned}
$$

and consequently

$$
\mathbb{E}\left|Y_{t}^{1}-Y_{t}^{2}\right| \leq L_{t} \int_{0}^{t} \mathbb{E}\left|Y_{s}^{1}-Y_{s}^{2}\right| d s
$$

From this inequality we get $|E| Y_{t}^{1}-Y_{t}^{2} \mid=0$ by means of Gronwall's lemma along with the usual continuity argument.

Example A.1.6 (Girsanov). Solutions to the one-dimensional SDE

$$
\begin{equation*}
d Y=\sigma_{\alpha}(Y) d B, \quad Y_{0}=0 \tag{A.1.6}
\end{equation*}
$$

with $\sigma_{\alpha}(x)=|x|^{\alpha} \wedge 1$ are pathwise unique for $\alpha \geq 1 / 2$, and $Y \equiv 0$ is the only solution. This is an immediate consequence of Theorem A.1.5: for $1 / 2 \leq \alpha \leq 1$ it holds that $\left|\sigma_{\alpha}(x)-\sigma_{\alpha}(y)\right| \leq|x-y|^{\alpha}$, whereas $\sigma_{\alpha}$ is globally Lipschitz continuous for $\alpha \geq 1$. We are going to verify that pathwise uniqueness in Equation (A.1.6) is violated for $0<\alpha<1 / 2$. To this end, let

$$
T(t):=\int_{0}^{t}\left(\left|B_{s}\right|^{2 \alpha} \wedge 1\right)^{-1} d s, \quad t \in \mathbb{R}_{+}
$$

As $\mathbb{E}[T(t)]<\infty$, each $T(t)$ is finite almost surely, $t \mapsto T(t)$ is almost surely strictly monotone increasing, and because of $T(t) \geq t$ trivially $T(\infty)=\infty$ holds. Hence $\tau_{t}:=$ $T^{-1}(t) \equiv \inf \left\{s \in \mathbb{R}_{+}: T(s)>t\right\}$ defines a finite continuous time-change, and $Y_{t}:=B_{\tau_{t}}$ with respect to the time-changed filtration $\left(\mathscr{F}_{\tau_{t}}\right)_{t \in \mathbb{R}_{+}}$gives a (non-trivial) weak solution $Y$ of (A.1.6). Indeed, with the $\left(\mathscr{F}_{\tau_{t}}\right)$-Brownian motion

$$
\tilde{B}_{t}:=\int_{0}^{\tau_{t}}\left(\left|B_{s}\right|^{\alpha} \wedge 1\right)^{-1} d B_{s}, \quad t \in \mathbb{R}_{+}
$$

we have

$$
Y_{t}=B_{\tau_{t}}=\int_{0}^{\tau_{t}} d B_{s}=\int_{0}^{t}\left(\left|B_{\tau_{s}}\right|^{\alpha} \wedge 1\right) d \tilde{B}_{s}=\int_{0}^{t}\left(\left|Y_{s}\right|^{\alpha} \wedge 1\right) d \tilde{B}_{s}
$$

which shows the claim. Hence uniqueness in distribution, and in particular pathwise uniqueness of solutions to (A.1.6), does not hold for $0<\alpha<1 / 2$.

REMARK A.1.7. It may be surprising in Example A.1.6 that unique solvability of (A.1.6) is given in cases where uniqueness of solutions in the analogous ordinary differential equation is violated. For instance, the equation $y(t)=\int_{0}^{t}\left(|y(s)|^{\alpha} \wedge 1\right) d s$ has for
$\alpha \geq 1$ only the trivial solution $y \equiv 0$, as can be seen by the Gronwall lemma, whereas for $0<\alpha<1$ a further solution is given by

$$
y(t):= \begin{cases}(\beta t)^{1 / \beta} & \text { for } 0 \leq t \leq 1 / \beta \\ (1-1 / \beta)+t & \text { for } t \geq 1 / \beta\end{cases}
$$

where $\beta=1-\alpha$.
We want now to use the above techniques for the goal to derive comparison theorems of one-dimensional SDEs. To this end, we consider the situation of two SDEs

$$
\begin{align*}
d Y^{1} & =\beta^{1}\left(t, Y^{1}\right) d t+\sigma\left(t, Y^{1}\right) d B \\
d Y^{2} & =\beta^{2}\left(t, Y^{2}\right) d t+\sigma\left(t, Y^{2}\right) d B \tag{A.1.7}
\end{align*}
$$

with continuous functions $\beta^{1}, \beta^{2}, \sigma: \mathbb{R}_{+} \times \mathbb{R} \rightarrow \mathbb{R}$ and the same one-dimensional Brownian motion $B$ as driving process.

THEOREM A.1.8 (Comparison Theorem of Ikeda-Watanabe). Let $Y^{1}$ and $Y^{2}$ be solutions of (A.1.7) under the following conditions:
(i) Either $\beta^{1}$ or $\beta^{2}$ is globally Lipschitz, and it holds $\beta^{1} \geq \beta^{2}$,
(ii) $\sigma$ satisfies Condition (i) of Theorem A.1.5.

Then $Y_{0}^{1} \geq Y_{0}^{2}$ almost surely already implies $Y_{t}^{1} \geq Y_{t}^{2}$ almost surely for any $t>0$.
Proof. By the same argument as in the proof of Theorem A.1.5 we obtain again $L^{0}\left(Y^{1}-Y^{2}\right)=0$. Since $Y_{0}^{1} \geq Y_{0}^{2}$ almost surely, we have by the Tanaka formula (A.1.3)

$$
\begin{aligned}
\left(Y_{t}^{2}-Y_{t}^{1}\right)_{+}= & \int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left(\beta^{2}\left(s, Y_{s}^{2}\right)-\beta^{1}\left(s, Y_{s}^{1}\right)\right) d s \\
& +\int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left(\sigma\left(s, Y_{s}^{2}\right)-\sigma\left(s, Y_{s}^{1}\right)\right) d B_{s}
\end{aligned}
$$

and hence

$$
\phi(t):=\mathbb{E}\left[\left(Y_{t}^{2}-Y_{t}^{1}\right)_{+}\right]=\mathbb{E}\left[\int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left[\beta^{2}\left(s, Y_{s}^{2}\right)-\beta^{1}\left(s, Y_{s}^{1}\right)\right] d s\right]
$$

In case $\beta^{1}$ is globally Lipschitz, we obtain

$$
\begin{aligned}
\phi(t) & \leq \mathbb{E}\left[\int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left[\beta^{1}\left(s, Y_{s}^{2}\right)-\beta^{1}\left(s, Y_{s}^{1}\right)\right] d s\right] \\
& \leq L_{t} \mathbb{E}\left[\int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left|Y_{s}^{2}-Y_{s}^{1}\right| d s\right]=L_{t} \int_{0}^{t} \phi(s) d s
\end{aligned}
$$

and the claim follows in the usual way by Gronwall's Lemma. On the other hand, if $\beta^{2}$ is globally Lipschitz, then

$$
\begin{aligned}
\phi(t)=\mathbb{E} & {\left[\int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left(\beta^{2}\left(s, Y_{s}^{2}\right)-\beta^{2}\left(s, Y_{s}^{1}\right)\right) d s\right] } \\
& +\mathbb{E}\left[\int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left(\beta^{2}\left(s, Y_{s}^{1}\right)-\beta^{1}\left(s, Y_{s}^{1}\right)\right) d s\right] \\
\leq \mathbb{E} & {\left[\int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left(\beta^{2}\left(s, Y_{s}^{2}\right)-\beta^{2}\left(s, Y_{s}^{1}\right)\right) d s\right] } \\
\leq L_{t} \mathbb{E}[ & {\left[\int_{0}^{t} 1_{\left\{Y_{s}^{2}>Y_{s}^{1}\right\}}\left|Y_{s}^{2}-Y_{s}^{1}\right| d s\right] }
\end{aligned}
$$

and the claim is derived as in the first case.
We now want to discuss the asymptotic behaviour of solutions to one-dimensional SDEs. For such questions it is natural to consider SDEs on open real intervals and then solutions with lifetime.

Let $I=] c_{1}, c_{2}\left[\right.$ with $-\infty \leq c_{1}<c_{2} \leq \infty$ be an interval in $\mathbb{R}$ and let $a, b: I \rightarrow \mathbb{R}$ be continuous functions where $a>0$. We suppose that on $I$ a diffusion process $Y$ with infinitesimal generator $L=a D^{2}+b D$ (where $D=d / d t$ ) is given, which we assume to be realized as maximal solution to an SDE of the form

$$
\begin{equation*}
d Y=b(Y) d t+\sigma(Y) d B \tag{A.1.8}
\end{equation*}
$$

with $\sigma^{2}=2 a, \sigma>0$ and $B$ a one-dimensional BM. In particular, $Y$ is then a continuous $I$-valued semimartingale of maximal, but not necessarily infinite lifetime $\zeta$, such that

$$
d(f(Y))-(L f)(Y) d t \in d \mathscr{M}
$$

for any $C^{2}$-function $f: I \rightarrow \mathbb{R}$ of compact support $\operatorname{supp}(f) \subset I$. Note that on the set $\{\zeta<\infty\}$ the limit $\lim _{t \uparrow \zeta} Y_{t}$ exists almost surely with values in $\left\{c_{1}, c_{2}\right\}$. Thus we may extend $Y$ via $Y_{t}:=\lim _{s \uparrow \zeta} Y_{s}$ on $\{\zeta<\infty\}$ for $t \geq \zeta$ to a continuous process defined globally on $\mathbb{R}_{+}$. For $x \in \bar{I}$ we denote by $\tau_{x}=\inf \left\{t \geq 0: Y_{t}=x\right\}$ the hitting time of $x$ and call the boundary point $c_{i}$ accessible for $Y$ if $\mathbb{P}\left\{\tau_{c_{i}}<\infty\right\}>0$.

The problem is now to find conditions on the coefficients of (A.1.8) which characterize properties such as transience, recurrence, or infinite lifetime of $Y$. The process $Y$ is called transient if $Y$ eventually exits every compact subset in $I$ almost surely, and recurrent if $\mathbb{P}\left\{\tau_{x}<\infty\right\}=1$ for each $x \in I$. The lifetime of $Y$ is obviously given by $\zeta=\tau_{c_{1}} \wedge \tau_{c_{2}} ;$ transience of $Y$ means $Y_{t} \rightarrow\left\{c_{1}, c_{2}\right\}$ as $t \uparrow \zeta$, and is in particular satisfied if $\zeta<\infty$ almost surely.

Fixing $c$ such that $c_{1}<c<c_{2}$, we consider

$$
\begin{equation*}
H: I \rightarrow \mathbb{R}, \quad H(r)=\exp \left\{-\int_{c}^{r} \frac{b(\rho)}{a(\rho)} d \rho\right\} \tag{A.1.9}
\end{equation*}
$$

as well as the $\mathbb{R}$-valued functions $s, m, k$ on $I$ defined by

$$
s(r)=\int_{c}^{r} H(\rho) d \rho, \quad m(r)=\int_{c}^{r} \frac{1}{a(\rho) H(\rho)} d \rho, \quad k(r)=\int_{c}^{r} m(\rho) s(d \rho),
$$

which extend to $\mathbb{R} \cup\{ \pm \infty\}$-valued functions on $\bar{I}=I \cup\left\{c_{1}, c_{2}\right\}$; here $s(d \rho)$ denotes the Borel measure on $I$ with distribution function $s$.

THEOREM A.1.9. Let $a, b: I \rightarrow \mathbb{R}$ be continuous functions on an interval $I=] c_{1}, c_{2}[$ where $-\infty \leq c_{1}<c_{2} \leq \infty$ and set $\sigma^{2}=2 a$ with $\sigma>0$. For a one-dimensional Brownian motion $B$ and $y \in I$, let $Y$ be the maximal solution to the SDE

$$
d Y=b(Y) d t+\sigma(Y) d B, \quad Y_{0}=y
$$

With respect to a fixed $c \in I$ let $H$ be defined by (A.1.9). The following items hold true:
(i) The process $Y$ is either recurrent or transient, and in fact transient if and only if $s\left(c_{i}\right)=\int_{c}^{c_{i}} H(r) d r$ is finite for $i=1$ or 2 .
More precisely, one can distinguish the following four cases:
(1) If $s\left(c_{1}\right)=-\infty$ and $s\left(c_{2}\right)=\infty$ then

$$
\mathbb{P}\{\zeta=\infty\}=\mathbb{P}\left\{\inf _{0 \leq t<\infty} Y_{t}=c_{1}\right\}=\mathbb{P}\left\{\sup _{0 \leq t<\infty} Y_{t}=c_{2}\right\}=1
$$

(2) If $s\left(c_{1}\right)>-\infty$ and $s\left(c_{2}\right)=\infty$ then

$$
\mathbb{P}\left\{\lim _{t \uparrow \zeta} Y_{t}=c_{1}\right\}=\mathbb{P}\left\{\sup _{0 \leq t<\zeta} Y_{t}<c_{2}\right\}=1
$$

(3) If $s\left(c_{1}\right)=-\infty$ and $s\left(c_{2}\right)<\infty$ then

$$
\mathbb{P}\left\{\inf _{0 \leq t<\zeta} Y_{t}>c_{1}\right\}=\mathbb{P}\left\{\lim _{t \uparrow \zeta} Y_{t}=c_{2}\right\}=1
$$

(4) If $s\left(c_{1}\right)>-\infty$ and $s\left(c_{2}\right)<\infty$ then

$$
\mathbb{P}\left\{\lim _{t \uparrow \zeta} Y_{t}=c_{1}\right\}=1-\mathbb{P}\left\{\lim _{t \uparrow \zeta} Y_{t}=c_{2}\right\}=\frac{s\left(c_{2}\right)-s(y)}{s\left(c_{2}\right)-s\left(c_{1}\right)}
$$

(ii) (Feller's test for explosion) $Y$ has almost surely infinite lifetime if and only if

$$
k\left(c_{i}\right)=\int_{c}^{c_{i}} H(r)\left(\int_{c}^{r} \frac{1}{a(\rho) H(\rho)} d \rho\right) d r=\infty, \quad i=1 \text { and } 2 .
$$

Moreover $c_{i}$ is accessible for $Y$ if $k\left(c_{i}\right)<\infty$.
(iii) The lifetime of $Y$ is almost surely finite (i.e., $\mathbb{P}\{\zeta<\infty\}=1$ ) if and only if one of the following three cases is at hand:
(1) $k\left(c_{1}\right)<\infty$ and $k\left(c_{2}\right)<\infty$, or
(2) $k\left(c_{1}\right)<\infty$ and $s\left(c_{2}\right)=\infty$, or
(3) $k\left(c_{2}\right)<\infty$ and $s\left(c_{1}\right)=-\infty$.

In case (1) even $\mathbb{E}[\zeta]<\infty$ holds.
The following implications hold trivially:

$$
\begin{equation*}
s\left(c_{1}\right)=-\infty \Rightarrow k\left(c_{1}\right)=\infty, \quad s\left(c_{2}\right)=\infty \Rightarrow k\left(c_{2}\right)=\infty \tag{A.1.10}
\end{equation*}
$$

The crucial method to prove Theorem A. 1.9 will be to "rescale" the process $Y$ by composition with an isotone transformation $\varphi$ in such a way that $\varphi(Y)$ becomes a local martingale. One speaks then of a "natural scale" for the diffusion $Y$ and calls $\varphi$ a scale function for $Y$. We start by verifying that actually $s$ defines a scale function for $Y$.

Proof of Theorem A.1.9. (a) The function $s$ defines a $C^{2}$-diffeomorphism of $I=$ $] c_{1}, c_{2}$ [onto $] s\left(c_{1}\right), s\left(c_{2}\right)$ [such that $\tilde{Y}=s(Y)$ is a local martingale with lifetime $\zeta$. Indeed, by $L s=a s^{\prime \prime}+b s^{\prime}=a H^{\prime}+b H=0$, we have

$$
d \tilde{Y}=d(s(Y))=s^{\prime}(Y) d Y+\frac{1}{2} s^{\prime \prime}(Y) d Y d Y=H(Y) \sigma(Y) d B
$$

and thus $d \tilde{Y}=\phi(\tilde{Y}) d B$ where $\phi:=(H \sigma) \circ s^{-1}$. In particular, modulo a time change, $\tilde{Y}$ is a Brownian motion, i.e., there exists a stopped one-dimensional Brownian motion $W$ such that $\tilde{Y}_{t}=W_{T_{t}}$ almost surely where $T_{t}:=\int_{0}^{t} \phi\left(\tilde{Y}_{s}\right)^{2} d s$.

We want to note first that $Y$ leaves each compact subinterval of $I$ in finite time: it holds $\mathbb{P}\left\{\tau_{a, b}<\zeta\right\}=1$ where

$$
\tau_{a, b}:=\inf \left\{t \geq 0: Y_{t} \notin[a, b]\right\}, \quad c_{1}<a<b<c_{2}
$$

Indeed, on $N:=\left\{\tau_{a, b}=\zeta\right\}$ it holds hat $\zeta=\infty$ by the maximality of the solution; the paths of $Y$ hence stay in the interval $[a, b]$, and the ones of $\tilde{Y}$ in $[s(a), s(b)]$. As consequence of $\tilde{Y}_{t}=W_{T_{t}}$, we get $T_{\zeta}<\infty$ on $N$ almost surely. On the other hand, we have $\phi \geq \varepsilon>0$ on the compact interval $[s(a), s(b)]$, and hence $T_{\zeta}=\infty$ on $N$. Both facts together imply $\mathbb{P}(N)=0$.

From the discussion above the claims of part (i) of the Theorem follow immediately. Indeed, the maximal lifetime $\zeta$ of $Y$ on $] c_{1}, c_{2}[$ coincides with the maximal lifetime of $\tilde{Y}$
on ] $s\left(c_{1}\right), s\left(c_{2}\right)\left[\right.$, and $Y$ is transient if and only if $\tilde{Y}$ is transient, for which $T_{\zeta}<\infty$ must hold almost surely. Conversely, convergence of $\tilde{Y}_{t}$ as $t \uparrow \zeta$ holds on $\left\{T_{\zeta}<\infty\right\}$, that is convergence almost surely to $s\left(c_{1}\right)$ or $s\left(c_{2}\right)$, since otherwise the path of $Y$ would stay in a compact subinterval of $I$ with the consequence that $T_{\zeta}=\infty$ as seen above. Hence transience of $Y$ is given exactly if $T_{\zeta}<\infty$ almost surely, and this is the case if and only if $\tilde{Y}_{t}$ converges almost surely to $s\left(c_{1}\right)$ or $s\left(c_{2}\right)$ as $t \uparrow \zeta$, hence if and only if $s\left(c_{1}\right)$ or $s\left(c_{2}\right)$ is finite.

Let $y \in I$ be the starting point of $Y$ und let $x, z \in I$ such that $x<y<z$. Then

$$
\begin{equation*}
\mathbb{P}\left\{\tau_{x} \leq \tau_{z}\right\}=\mathbb{P}\left\{\tau_{x}<\tau_{z}\right\}=\frac{s(z)-s(y)}{s(z)-s(x)} \tag{A.1.11}
\end{equation*}
$$

as can be seen from the equality

$$
s(y)=\mathbb{E}\left[s\left(Y_{0}\right)\right]=\mathbb{E}\left[s\left(Y_{\tau_{x} \wedge \tau_{z}}\right)\right]=s(x) \mathbb{P}\left\{\tau_{x}<\tau_{z}\right\}+s(z) \mathbb{P}\left\{\tau_{x} \geq \tau_{z}\right\}
$$

Recall that $Y$ is recurrent if and only if for any $x \in I$,

$$
P\left\{\tau_{x}<\tau_{c_{i}}\right\}=1, \quad i=1,2
$$

From (A.1.11) we conclude that this is equivalent to $s\left(c_{1}\right)=-\infty, s\left(c_{2}\right)=\infty$; more precisely, we have

$$
\begin{aligned}
& \mathbb{P}\left\{\tau_{x}<\tau_{c_{2}}\right\}=1 \Longleftrightarrow \lim _{z \uparrow c_{2}} \mathbb{P}\left\{\tau_{x}<\tau_{z}\right\}=1 \Longleftrightarrow s\left(c_{2}\right)=\infty, \quad c_{1}<x<y \\
& \mathbb{P}\left\{\tau_{x}<\tau_{c_{1}}\right\}=1 \Longleftrightarrow \lim _{z \searrow c_{1}} \mathbb{P}\left\{\tau_{z} \leq \tau_{x}\right\}=0 \Longleftrightarrow s\left(c_{1}\right)=-\infty, \quad y<x<c_{2}
\end{aligned}
$$

This shows in particular the claimed criterion for recurrence. The items (1) - (4) of part (i) are immediate combinations of the above.
(b) For the analysis of explosions of $Y$ we construct a twice continuously differentiable function $\psi: I \rightarrow \mathbb{R}_{+}$such that $1+k \leq \psi \leq \exp (k)$ and such that $Z_{t}:=e^{-t}\left(\psi \circ Y_{t}\right)$ defines a local martingale on $[0, \zeta[$.

More specifically, let $\psi: I \rightarrow \mathbb{R}$ be the unique solution to the linear SDE

$$
\begin{equation*}
L \psi=\psi \text { on } I \text { with } \psi^{\prime}(c)=0 \text { and } \psi(c)=1 \tag{A.1.12}
\end{equation*}
$$

We want to show that $\psi$ has the intended properties. For a continuous function $u$ on $I$, let $M(u)$ be the function on $I$ defined by

$$
M(u)(r):=\int_{c}^{r} s(d \rho)\left(\int_{c}^{\rho} u(t) m(d t)\right)=\int_{c}^{r} H(\rho)\left(\int_{c}^{\rho} \frac{u(t)}{H(t) a(t)} d t\right) d \rho
$$

From the equation $a H^{\prime}+b H \equiv 0$ it follows immediately that $L M(u)=u$ on $I$; in particular, condition (A.1.12) is equivalent to the validity of the equation $\psi=1+M(\psi)$ on $I$. This leads to the presentation

$$
\begin{equation*}
\psi=\sum_{n=0}^{\infty} M^{n}(1) \tag{A.1.13}
\end{equation*}
$$

where $M^{0}(1):=1$ and $M^{n+1}(1)=M\left(M^{n}(1)\right)$. We have $k=M(1)$ by definition, trivially $M^{n}(1) \geq 0$, and one verifies inductively

$$
\begin{equation*}
M^{n}(1) \leq k^{n} / n!, \quad n=1,2, \ldots \tag{A.1.14}
\end{equation*}
$$

Indeed if $M^{n}(1) \leq \frac{k^{n}}{n!}$ then also

$$
\begin{aligned}
M^{n+1}(1)(r) & =\int_{c}^{r} s(d \rho)\left(\int_{c}^{\rho} M^{n}(1)(t) m(d t)\right) \leq \frac{1}{n!} \int_{c}^{r} s(d \rho)\left(\int_{c}^{\rho} k^{n}(t) m(d t)\right) \\
& \leq \frac{1}{n!} \int_{c}^{r} s(d \rho) k^{n}(\rho)\left(\int_{c}^{\rho} m(d t)\right)=\frac{1}{n!} \int_{c}^{r} k^{n}(\rho) k(d \rho)
\end{aligned}
$$

$$
=\frac{1}{n!} \int_{c}^{r} k^{n}(\rho) k^{\prime}(\rho) d \rho=\frac{k^{n+1}(r)}{(n+1)!} .
$$

By (A.1.14) the right-hand side $\psi^{*}:=\sum_{n=0}^{\infty} M^{n}(1)$ of (A.1.13) is well-defined and satisfies the estimate $1+k \leq \psi^{*} \leq \exp (k)$; on the other hand, as $1+M\left(\psi^{*}\right)=\psi^{*}$ on $I$, we get $\psi=\psi^{*}$. In particular, the representation (A.1.13) shows that $\psi$ is decreasing on $] c_{1}, c$ ] and increasing on $\left[c, c_{2}[\right.$.

The claim that $\left(Z_{t}\right)_{t<\zeta}$ with $Z_{t}=e^{-t} \psi\left(Y_{t}\right)$ is a local martingale is finally seen from Itô's formula using the fact that $a \psi^{\prime \prime}+b \psi^{\prime}=\psi$.
(c) If $k\left(c_{1}\right)=k\left(c_{2}\right)=\infty$, then $\mathbb{P}\{\zeta=\infty\}=1$ holds for the lifetime $\zeta \equiv \tau_{c_{1}} \wedge \tau_{c_{2}}$.

To see this, we chose a compact exhaustion $\left.\left[a_{n}, b_{n}\right] \uparrow\right] c_{1}, c_{2}$ [ where $a_{n}<y<b_{n}$ and consider $\sigma_{n}=\inf \left\{t \geq 0: Y_{t} \notin\left[a_{n}, b_{n}\right]\right\} \equiv \tau_{a_{n}} \wedge \tau_{b_{n}}$. Since $Z^{\sigma_{n}}$ is a martingale, we have for each $t \in \mathbb{R}_{+}$the equality $\psi(y)=E\left[e^{-\left(\sigma_{n} \wedge t\right)} \psi\left(Y_{\sigma_{n} \wedge t}\right)\right]$ and thus

$$
\mathbb{P}\left\{\sigma_{n}<t\right\} \leq \frac{e^{t} \psi(y)}{\psi\left(a_{n}\right) \wedge \psi\left(b_{n}\right)} \rightarrow 0 \quad \text { for } n \rightarrow \infty
$$

here we use the assumption $k\left(c_{1}\right)=k\left(c_{2}\right)=\infty$ along with the estimate $1+k \leq \psi$.
(d) If $k\left(c_{i}\right)<\infty$, then $c_{i}$ is accessible for $Y$, i.e., it holds that $\mathbb{P}\left\{\tau_{c_{i}}<\infty\right\}>0$. In particular, $\mathbb{P}\{\zeta<\infty\}>0$ holds if $k\left(c_{i}\right)$ is finite for $i=1$ or 2 .

Suppose for instance that $k\left(c_{1}\right)$ is finite; because of $\psi \leq \exp (k)$ then $\psi(x)$ stays bounded as $x \rightarrow c_{1}+$. Without restrictions we may assume that $c_{1}<y<c$. In addition to the sequence $\left(\sigma_{n}\right)_{n \in \mathbb{N}}$ of exit times described in (c), we consider the hitting time $\sigma_{0}:=\tau_{c}$ of $c$. Taking into account that $\sigma_{n}<\zeta$, we have

$$
\begin{aligned}
1 & <\psi(y)=\mathbb{E}\left[e^{-\left(\sigma_{n} \wedge \sigma_{0}\right)}\left(\psi \circ Y_{\sigma_{n} \wedge \sigma_{0}}\right)\right] \\
& =\mathbb{E}\left[1_{\left\{\sigma_{n} \geq \sigma_{0}\right\}} e^{-\sigma_{0}} \psi(c)+1_{\left\{\sigma_{n}<\sigma_{0}\right\}} e^{-\sigma_{n}}\left(\psi \circ Y_{\sigma_{n}}\right)\right] \\
& \leq 1+\psi\left(c_{1}+\right) \mathbb{E}\left[1_{\left\{\sigma_{n}<\sigma_{0}\right\}} e^{-\sigma_{n}}\right] \\
& \downarrow 1+\psi\left(c_{1}+\right) E\left[1_{\left\{\tau_{c_{1}}<\sigma_{0}\right\}} e^{-\tau_{c_{1}}}\right] \quad \text { as } n \rightarrow \infty,
\end{aligned}
$$

where we used that the function $\psi$ is decreasing on the subinterval $\left.] c_{1}, c\right]$; in particular, then $\mathbb{P}\left\{\tau_{c_{1}}<\sigma_{0}\right\}>0$ holds true. The case $k\left(c_{2}\right)<\infty$ is treated analogously. This completes the proof of part (ii) of the Theorem.
(e) In order to verify (iii) we show first: If $\mathbb{P}\{\zeta<\infty\}=1$ holds, then one of the cases (1), (2) or (3) is in force.

If $\mathbb{P}\{\zeta<\infty\}=1$, then $k\left(c_{1}\right)<\infty$ or $k\left(c_{2}\right)<\infty$ by (ii). Assume for instance $k\left(c_{1}\right)<\infty$, and suppose that none of the cases (1), (2), (3) is given. Taking (A.1.10) into account, we see then that

$$
s\left(c_{1}\right)>-\infty \quad \text { and } \quad s\left(c_{2}\right)<\infty, \quad k\left(c_{2}\right)=\infty
$$

Hence we are in situation (4) of part (i), and $\mathbb{P}\left\{\lim _{t \uparrow \zeta} Y_{t}=c_{2}\right\}>0$ holds true. We consider the time-changed process $\left(Z_{\tau_{t}}:=e^{-\tau_{t}} \psi\left(Y_{\tau_{t}}\right)\right)_{t \in \mathbb{R}_{+}}$where the time change $\left(\tau_{t}\right)_{t \in \mathbb{R}_{+}}$ stretches the stochastic interval $\left[0, \zeta\left[\right.\right.$ to $\mathbb{R}_{+} \times \Omega$. As $Z_{\tau_{t}}$ is a non-negative continuous local martingale, the limit

$$
Z_{\tau_{\infty}}=\lim _{t \uparrow \zeta} e^{-t} \psi\left(Y_{t}\right)
$$

almost surely exists in $\mathbb{R}$ and takes the value $e^{-\zeta} \psi\left(c_{2}+\right)$ on $\left\{\lim _{t \uparrow \zeta} Y_{t}=c_{2}\right\}$. Since $1+k \leq \psi \leq \exp (k)$ and $k\left(c_{2}\right)=\infty$, it holds that $\psi\left(c_{2}-\right)=\infty$, with the consequence that $\zeta=\infty$ on $\left\{\lim _{t \uparrow \zeta} Y_{t}=c_{2}\right\}$. This is however in contradiction to $\mathbb{P}\{\zeta<\infty\}=1$; hence either (1), (2) or (3) must hold true. One argues analogously in the case $k\left(c_{2}\right)<\infty$.
(f) It remains to show that $\mathbb{P}\{\zeta<\infty\}=1$ if in (iii) one of the cases (1), (2), (3) is given. We verify first that $\mathbb{E}[\zeta]<\infty$ under the assumption $k\left(c_{1}\right)<\infty, k\left(c_{2}\right)<\infty$.

To this end, we construct under the condition $k\left(c_{1}\right)<\infty, k\left(c_{2}\right)<\infty$ a twice continuously differentiable non-negative function $u: I \rightarrow \mathbb{R}$ such that

$$
L u \equiv-1 \quad \text { and } \quad u\left(c_{1}+\right)=u\left(c_{2}-\right)=0
$$

We set $u(r)=\int_{I} G(r, t) m(d t)$ where

$$
G(r, t):=\frac{\left(s(r \wedge t)-s\left(c_{1}\right)\right)\left(s\left(c_{2}\right)-s(r \vee t)\right)}{s\left(c_{2}\right)-s\left(c_{1}\right)}, \quad(r, t) \in I \times I
$$

Under the assumption that $k\left(c_{1}\right)<\infty, k\left(c_{2}\right)<\infty$ we have $s\left(c_{i}\right) \in \mathbb{R}$ for $i=1,2$ and hence $G$ is bounded. Furthermore, we have

$$
\begin{aligned}
u(r) & =\frac{s\left(c_{2}\right)-s(r)}{s\left(c_{2}\right)-s\left(c_{1}\right)} \int_{c_{1}}^{r}\left(s(t)-s\left(c_{1}\right)\right) m(d t)+\frac{s(r)-s\left(c_{1}\right)}{s\left(c_{2}\right)-s\left(c_{1}\right)} \int_{r}^{c_{2}}\left(s\left(c_{2}\right)-s(t)\right) m(d t) \\
& =\frac{s\left(c_{2}\right)-s(r)}{s\left(c_{2}\right)-s\left(c_{1}\right)} \int_{c_{1}}^{r}(m(r)-m(\rho)) s(d \rho)+\frac{s(r)-s\left(c_{1}\right)}{s\left(c_{2}\right)-s\left(c_{1}\right)} \int_{r}^{c_{2}}(m(\rho)-m(r)) s(d \rho) \\
& =-\frac{s\left(c_{2}\right)-s(r)}{s\left(c_{2}\right)-s\left(c_{1}\right)} \int_{c_{1}}^{r} m(\rho) s(d \rho)+\frac{s(r)-s\left(c_{1}\right)}{s\left(c_{2}\right)-s\left(c_{1}\right)} \int_{r}^{c_{2}} m(\rho) s(d \rho)
\end{aligned}
$$

where we used in the second line the conversion

$$
\begin{aligned}
\int_{r}^{c_{2}}\left(s\left(c_{2}\right)-s(t)\right) m(d t) & =\int_{r}^{c_{2}}\left(\int_{t}^{c_{2}} H(\rho) d \rho\right) m(d t) \\
& =\int_{r}^{c_{2}} \int_{r}^{c_{2}} H(\rho) 1_{\{t \leq \rho\}}(t, \rho) d \rho m(d t) \\
& =\int_{r}^{c_{2}}\left(\int_{r}^{\rho} m(d t)\right) H(\rho) d \rho=\int_{r}^{c_{2}}(m(\rho)-m(r)) s(d \rho)
\end{aligned}
$$

The computation above shows that under the assumption that $k\left(c_{1}\right)<\infty, k\left(c_{2}\right)<\infty$ the function $u$ is finite on $I$; evidently even bounded and twice continuously differentiable. In addition, one verifies $L u \equiv-1$; trivially $u\left(c_{1}+\right)=u\left(c_{2}-\right)=0$ holds. Itô's formula then gives

$$
d(u(Y))=u^{\prime}(Y) \sigma(Y) d B-d t
$$

Choosing now as in part (c) a compact exhaustion $\left.\left[a_{n}, b_{n}\right] \uparrow\right] c_{1}, c_{2}$ [ such that $a_{n}<y<b_{n}$ and considering the stopping times $\sigma_{n}=\inf \left\{t \geq 0: Y_{t} \notin\left[a_{n}, b_{n}\right]\right\}$, we obtain

$$
\mathbb{E}\left[u\left(Y_{t \wedge \sigma_{n}}\right)\right]=u(y)-\mathbb{E}\left[t \wedge \sigma_{n}\right]
$$

Hence we get $\mathbb{E}\left[t \wedge \sigma_{n}\right] \leq u(y)$, and as $t \rightarrow \infty, n \rightarrow \infty$, we conclude $\mathbb{E} \zeta \leq u(y)<\infty$.
(g) We show: the conditions $k\left(c_{1}\right)<\infty$ and $s\left(c_{2}\right)=\infty$ imply $\mathbb{P}\{\zeta<\infty\}=1$.

If $\left(a_{n}\right)_{n \in \mathbb{N}}$ is a sequence in $I$ such that $a_{n} \uparrow c_{2}$, then it holds $\mathbb{P}\left\{\tau_{c_{1}} \wedge \tau_{a_{n}}<\infty\right\}=1$ according to (f); on the other hand, $k\left(c_{1}\right)<\infty$ implies $s\left(c_{1}\right)>-\infty$, so that situation (2) of part (i) is given, which implies $\lim _{n \rightarrow \infty} \mathbb{P}\left\{\tau_{a_{n}}>\tau_{c_{1}}\right\}=1$. By the obvious identity $\left\{\tau_{c_{1}}<\infty\right\}=\bigcup_{n}\left\{\tau_{c_{1}}<\tau_{a_{n}}\right\}$ one obtains then $\mathbb{P}\left\{\tau_{c_{1}}<\infty\right\}=1$ and $\mathbb{P}\{\zeta<\infty\}=1$ as wanted.

Analogously, one verifies $\mathbb{P}\{\zeta<\infty\}=1$ in the remaining case $k\left(c_{2}\right)<\infty$ and $s\left(c_{1}\right)=-\infty$.

## A.2. Derivative Flows

Let $M$ be an $n$-dimensional smooth manifold and, for some $m \in \mathbb{N}$, let

$$
A: M \times \mathbb{R}^{m} \rightarrow T M, \quad(x, e) \mapsto A(x) e
$$

be a homomorphism of vector bundles over $M$. Thus, $A \in \Gamma\left(\mathbb{R}^{m} \otimes T M\right)$, i.e., the map $A(x): \mathbb{R}^{m} \rightarrow T_{x} M$ is linear for $x \in M$, and $A(\cdot) e \in \Gamma(T M)$ is a smooth vector field on $M$ for $e \in \mathbb{R}^{m}$. Consider the Stratonovich stochastic differential equation

$$
\begin{equation*}
d X=A(X) \circ d B+A_{0}(X) d t \tag{A.2.1}
\end{equation*}
$$

where $A_{0} \in \Gamma(T M)$ is an additional vector field, and $B$ an $\mathbb{R}^{m}$-valued Brownian motion on a filtered probability space $\left(\Omega, \mathscr{F}, \mathbb{P} ;\left(\mathscr{F}_{t}\right)_{t \in \mathbb{R}_{+}}\right)$satisfying the usual completeness conditions. There is a partial flow $X_{t}(\cdot), \zeta(\cdot)$ associated to (A.2.1) (see [29] for details) such that for each $x \in M$ the process $X_{t}(x), 0 \leq t<\zeta(x)$, is the maximal strong solution to (A.2.1) with starting point $X_{0}(x)=x$, defined up to the explosion time $\zeta(x)$; moreover, using the notation $X_{t}(x, \omega)=X_{t}(x)(\omega)$ and $\zeta(x, \omega)=\zeta(x)(\omega)$, if

$$
M_{t}(\omega)=\{x \in M: t<\zeta(x, \omega)\}
$$

then there exists a set $\Omega_{0} \subset \Omega$ of full measure such that for all $\omega \in \Omega_{0}$ :
(i) $M_{t}(\omega)$ is open in $M$ for each $t \geq 0$, i.e., $\zeta(\cdot, \omega)$ is lower semicontinuous on $M$.
(ii) $X_{t}(\cdot, \omega): M_{t}(\omega) \rightarrow M$ is a diffeomorphism onto an open subset of $M$.
(iii) The map $s \mapsto X_{s}(\cdot, \omega)$ is continuous from $[0, t]$ into $C^{\infty}\left(M_{t}(\omega), M\right)$ with its $C^{\infty}$-topology, for each $t>0$.
The solution processes $X=X(x)$ to A.2.1 are diffusions on $M$ with generator

$$
L=A_{0}+\frac{1}{2} \sum_{i=1}^{m} A_{i}^{2}
$$

where $A_{i}=A(\cdot) e_{i} \in \Gamma(T M), i=1, \ldots, m$.
Consider the special case that the system A.2.1 is non-degenerate (elliptic), in te sense that $A(x): \mathbb{R}^{m} \rightarrow T_{x} M$ is surjective for each $x$, or equivalently that $L$ is an elliptic operator. This non-degeneracy provides a Riemannian metric on $M$ such that $A(x) A(x)^{*}: T_{x} M \rightarrow T_{x} M$ is the identity on $T_{x} M$ for $x \in M$. Then $A(x)^{*}: T_{x} M \rightarrow$ $\mathbb{R}^{m}$ defines an isometric inclusion for each $x \in M$, i.e.,

$$
\langle u, v\rangle_{T_{x} M}=\left\langle A(x)^{*} u, A(x)^{*} v\right\rangle_{\mathbb{R}^{m}} \quad \text { for all } u, v \in T_{x} M
$$

With respect to this Riemannian metric, $L=\frac{1}{2} \Delta_{M}+Z$ where $Z$ is of first order, i.e., a vector field on $M$. Standard examples are the gradient Brownian systems when $M$ is immersed into some Euclidean space $\mathbb{R}^{m}$, and $A(x): \mathbb{R}^{m} \rightarrow T_{x} M$ is the orthogonal projection; for $A_{0}=0$ this construction gives Brownian motion on $M$ with respect to the induced metric, see [11].

For $x \in M$, let $T_{x} X_{t}: T_{x} M \rightarrow T_{X_{t}(x)} M$ be the differential of $X_{t}(\cdot)$ at $x$ (welldefined for all $\omega \in \Omega$ such that $\left.x \in M_{t}(\omega)\right)$ and $V_{t}=V_{t}(v)=\left(T_{x} X_{t}\right) v$ the derivative process to $X_{t}(\cdot)$ at $x$ in direction $v \in T_{x} M$. It is well-known that $V$ on $T M$ solves the formally differentiated $\operatorname{SDE}(\mathrm{A} .2 .1)$, i.e.,

$$
\begin{equation*}
d V=\left(T_{X} A\right) V \circ d B+\left(T_{X} A_{0}\right) V d t, \quad V_{0}=v \tag{A.2.2}
\end{equation*}
$$

with the same lifetime as $X(x)$, if $v \neq 0$. Using the metric and the corresponding LeviCivita connection on $M$, Eq. (A.2.2) is most concisely written as a covariant equation along $X$

$$
\begin{equation*}
D V=(\nabla A) V \circ d B+\left(\nabla A_{0}\right) V d t \tag{A.2.3}
\end{equation*}
$$

(see [11]); by definition, (A.2.3) means

$$
d \tilde{V}=/ /{ }_{0}^{-1} t(\nabla A) / /{ }_{0} t \tilde{V} \circ d B+/ / 0_{0}^{-1} t\left(\nabla A_{0}\right) / /{ }_{0} t \tilde{V} d t
$$

for $\tilde{V}_{t}=/ /_{0}^{-1} t V_{t}$ where $/ /{ }_{0} t: T_{X_{0}} M \rightarrow T_{X_{t}} M$ is parallel transport along the paths of $X$.

## Bibliography

1. Michael T. Anderson and Richard Schoen, Positive harmonic functions on complete manifolds of negative curvature, Ann. of Math. (2) 121 (1985), no. 3, 429-461. MR 794369
2. Gérard Ben Arous, Shigeo Kusuoka, and Daniel W. Stroock, The Poisson kernel for certain degenerate elliptic operators, J. Funct. Anal. 56 (1984), no. 2, 171-209. MR 738578 (85k:35093)
3. Jean-Michel Bismut, Large deviations and the Malliavin calculus, Progress in Mathematics, vol. 45, Birkhäuser Boston, Inc., Boston, MA, 1984. MR 755001
4. Isaac Chavel, Eigenvalues in Riemannian geometry, Pure and Applied Mathematics, vol. 115, Academic Press, Inc., Orlando, FL, 1984, Including a chapter by Burton Randol, With an appendix by Jozef Dodziuk. MR 768584
5. Michael Cranston, Wilfrid S. Kendall, and Peter March, The radial part of Brownian motion. II. Its life and times on the cut locus, Probab. Theory Related Fields 96 (1993), no. 3, 353-368. MR 1231929
6. Richard W. R. Darling, Martingales in manifolds-definition, examples, and behaviour under maps, Seminar on Probability, XVI, Supplement, Lecture Notes in Math., vol. 921, Springer, Berlin-New York, 1982, pp. 217-236. MR 658727
7. James Eells and K. David Elworthy, Wiener integration on certain manifolds, Problems in non-linear analysis (C.I.M.E., IV Ciclo, Varenna, 1970), 1971, pp. 67-94. MR 0346835
8. K. David Elworthy, Stochastic dynamical systems and their flows, Stochastic analysis (Proc. Internat. Conf., Northwestern Univ., Evanston, Ill., 1978), Academic Press, New York-London, 1978, pp. 79-95. MR 517235
9. $\qquad$ , Stochastic methods and differential geometry, Bourbaki Seminar, Vol. 1980/81, Lecture Notes in Math., vol. 901, Springer, Berlin-New York, 1981, pp. 95-110. MR 647491
10. _ Stochastic differential equations on manifolds, London Mathematical Society Lecture Note Series, vol. 70, Cambridge University Press, Cambridge-New York, 1982. MR 675100
11. , Geometric aspects of diffusions on manifolds, École d'Été de Probabilités de Saint-Flour XV-XVII, 1985-87, Lecture Notes in Math., vol. 1362, Springer, Berlin, 1988, pp. 277-425. MR 983375
12. Michel Émery, Stochastic calculus in manifolds, Universitext, Springer-Verlag, Berlin, 1989, With an appendix by P.-A. Meyer. MR 1030543
13. Sylvestre Gallot, Dominique Hulin, and Jacques Lafontaine, Riemannian geometry, third ed., Universitext, Springer-Verlag, Berlin, 2004. MR 2088027
14. R. E. Greene and H. Wu, Function theory on manifolds which possess a pole, Lecture Notes in Mathematics, vol. 699, Springer, Berlin, 1979. MR 521983
15. Victor Guillemin and Alan Pollack, Differential topology, Prentice-Hall, Inc., Englewood Cliffs, N.J., 1974. MR 0348781
16. Wolfgang Hackenbroch and Anton Thalmaier, Stochastische Analysis, Mathematische Leitfäden. [Mathematical Textbooks], B. G. Teubner, Stuttgart, 1994, Eine Einführung in die Theorie der stetigen Semimartingale. [An introduction to the theory of continuous semimartingales]. MR 1312827
17. Elton P. Hsu, Stochastic analysis on manifolds, Graduate Studies in Mathematics, vol. 38, American Mathematical Society, Providence, RI, 2002. MR 1882015 (2003c:58026)
18. Pei Hsu and Peter March, The limiting angle of certain Riemannian Brownian motions, Comm. Pure Appl. Math. 38 (1985), no. 6, 755-768. MR 812346
19. Kanji Ichihara, Curvature, geodesics and the Brownian motion on a Riemannian manifold. I. Recurrence properties, Nagoya Math. J. 87 (1982), 101-114. MR 676589
20. , Curvature, geodesics and the Brownian motion on a Riemannian manifold. II. Explosion properties, Nagoya Math. J. 87 (1982), 115-125. MR 676590
21. Nobuyuki Ikeda and Shojiro Manabe, Integral of differential forms along the path of diffusion processes, Publ. Res. Inst. Math. Sci. 15 (1979), no. 3, 827-852. MR 566084
22. Nobuyuki Ikeda and Shinzo Watanabe, Stochastic differential equations and diffusion processes, second ed., North-Holland Mathematical Library, vol. 24, North-Holland Publishing Co., Amsterdam; Kodansha, Ltd., Tokyo, 1989. MR 1011252
23. Wilfrid S. Kendall, The radial part of Brownian motion on a manifold: a semimartingale property, Ann. Probab. 15 (1987), no. 4, 1491-1500. MR 905343
24. , Martingales on manifolds and harmonic maps, Geometry of random motion (Ithaca, N.Y., 1987), Contemp. Math., vol. 73, Amer. Math. Soc., Providence, RI, 1988, pp. 121-157. MR 954635
25. Wilhelm P. A. Klingenberg, Riemannian geometry, second ed., De Gruyter Studies in Mathematics, vol. 1, Walter de Gruyter \& Co., Berlin, 1995. MR 1330918
26. Shoshichi Kobayashi and Katsumi Nomizu, Foundations of differential geometry. Vol I, Interscience Publishers, a division of John Wiley \& Sons, New York-London, 1963. MR 0152974
27. ___ Foundations of differential geometry. Vol. II, Interscience Tracts in Pure and Applied Mathematics, No. 15 Vol. II, Interscience Publishers John Wiley \& Sons, Inc., New York-London-Sydney, 1969. MR 0238225
28. P. E. Kopp, Martingales and stochastic integrals, Cambridge University Press, Cambridge, 1984. MR 774050
29. Hiroshi Kunita, Stochastic flows and stochastic differential equations, Cambridge Studies in Advanced Mathematics, vol. 24, Cambridge University Press, Cambridge, 1990. MR 1070361
30. Terry Lyons and Dennis Sullivan, Function theory, random paths and covering spaces, J. Differential Geom. 19 (1984), no. 2, 299-323. MR 755228
31. Paul Malliavin, Formules de la moyenne, calcul de perturbations et théorèmes d'annulation pour les formes harmoniques, J. Functional Analysis 17 (1974), 274-291. MR 0385932
32. $\qquad$ , Géométrie différentielle stochastique, Séminaire de Mathématiques Supérieures [Seminar on Higher Mathematics], vol. 64, Presses de l’Université de Montréal, Montreal, Que., 1978, Notes prepared by Danièle Dehen and Dominique Michel. MR 540035
33. _ Stochastic analysis, Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences], vol. 313, Springer-Verlag, Berlin, 1997. MR 1450093
34. Peter March, Brownian motion and harmonic functions on rotationally symmetric manifolds, Ann. Probab. 14 (1986), no. 3, 793-801. MR 841584
35. P.-A. Meyer, A differential geometric formalism for the Itô calculus, Stochastic integrals (Proc. Sympos., Univ. Durham, Durham, 1980), Lecture Notes in Math., vol. 851, Springer, Berlin, 1981, pp. 256-270. MR 620993
36._, Géométrie stochastique sans larmes, Seminar on Probability, XV (Univ. Strasbourg, Strasbourg, 1979/1980) (French), Lecture Notes in Math., vol. 850, Springer, Berlin, 1981, pp. 44-102. MR 622555
36. Philip E. Protter, Stochastic integration and differential equations, Stochastic Modelling and Applied Probability, vol. 21, Springer-Verlag, Berlin, 2005, Second edition. Version 2.1, Corrected third printing. MR 2273672 (2008e:60001)
37. Daniel Revuz and Marc Yor, Continuous martingales and Brownian motion, third ed., Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences], vol. 293, Springer-Verlag, Berlin, 1999. MR 1725357
38. Laurent Schwartz, Semi-martingales sur des variétés, et martingales conformes sur des variétés analytiques complexes, Lecture Notes in Mathematics, vol. 780, Springer, Berlin, 1980. MR 575167
39. $\qquad$ , Géométrie différentielle du 2 ème ordre, semi-martingales et équations différentielles stochastiques sur une variété différentielle, Seminar on Probability, XVI, Supplement, Lecture Notes in Math., vol. 921, Springer, Berlin-New York, 1982, pp. 1-148. MR 658722
40. Ichiro Shigekawa, On stochastic horizontal lifts, Z. Wahrsch. Verw. Gebiete 59 (1982), no. 2, 211-221. MR 650613
41. Daniel W. Stroock, An introduction to the analysis of paths on a Riemannian manifold, Mathematical Surveys and Monographs, vol. 74, American Mathematical Society, Providence, RI, 2000. MR 1715265
42. Daniel W. Stroock and S. R. S. Varadhan, On degenerate elliptic-parabolic operators of second order and their associated diffusions, Comm. Pure Appl. Math. 25 (1972), 651-713. MR 0387812 (52 \#8651)
43. Dennis Sullivan, The Dirichlet problem at infinity for a negatively curved manifold, J. Differential Geom. 18 (1983), no. 4, 723-732 (1984). MR 730924
44. Frank W. Warner, Foundations of differentiable manifolds and Lie groups, Graduate Texts in Mathematics, vol. 94, Springer-Verlag, New York-Berlin, 1983, Corrected reprint of the 1971 edition. MR 722297
45. David Williams, Probability with martingales, Cambridge Mathematical Textbooks, Cambridge University Press, Cambridge, 1991. MR 1155402
46. Shing Tung Yau, Some function-theoretic properties of complete Riemannian manifold and their applications to geometry, Indiana Univ. Math. J. 25 (1976), no. 7, 659-670. MR 417452
47. Marc Yor, Some aspects of Brownian motion. Part I, Lectures in Mathematics ETH Zürich, Birkhäuser Verlag, Basel, 1992, Some special functionals. MR 1193919
49._, Some aspects of Brownian motion. Part II, Lectures in Mathematics ETH Zürich, Birkhäuser Verlag, Basel, 1997, Some recent martingale problems. MR 1442263

## Index

$(r, s)$-tensor, 50
$G$-connection, 66
$G$-invariant splitting, 66
$H^{p}$-martingale, 105
$L$-diffusion, 12
$\Gamma$-Operator, 18
grad gradient, 57
$\nabla$-martingale, 52
as solution of an SDE, 53
Darling's characterization, 98
$b$-quadratic variation, 36
PDO
purely second order, 120
accessible boundary point, 195
action
effective, 63
free, 63
linear, 63
affine map, 89
pullback properties, 97
stochastic characterization, 93
anti-development, 74
associated fiber bundle, 65
atlas, 1
Bessel process, 123,124
BM-complete manifold, 103
BM-completeness
criterion, 176
boundary point
accessible, 195
Brownian motion
manifold-valued, 59
one-dimensional, 86
radial process, 156
Brownian motion
horizontal, 91
Brownian motion, 59
as solution of an SDE, 60
comparison theorem, 157,166
on a Riemannian manifold, 59
on models, 170
radial part, 156
a
$G$-connection, 66
litting, 66
$H^{p}$-martingale, 105
$L$-diffusion, 12
Г-Operator, 18
grad gradient, 57
$\nabla$-martingale, 52
as solution of an SDE, 53
Darling's characterization, 98
$b$-quadratic variation, 36
PDO
accessible boundary point, 195
action
ive, 63
free, 63
linear, 63
pullback properties, 97
stochastic characterization, 93
associated fiber bundle, 65
atlas, 1
Bessel process, 123, 124
BM-complete manifold, 103
-completeness
boundary
accessible, 195
Brownian motion manifold-valued, 59 radial process, 156

Brownian motion
horizontal, 91
Brownian motion, 59 comparison theorem, 157, 166 on models, 170 radial part, 156
recurrence, 168
strong Markov property, 159
transience, 168
bundle
normal, 32
principal, 64
bundle atlas, 6
bundle chart, 6
bundle homomorphism, 7
canonical one-form, 71
of a Lie group, 77
Cartan development, 81
Cartan's structural equations, 51
chain rule, 5
chart, 1
Christoffel symbol, 45
comparison theorem
Hessian, 146
comparison criterion
for BM-completeness, 175
for transience, 175
Comparison principle, 142
comparison theorem
for Brownian motions, 157
for one-dimensional SDE, 194
for the Laplacian; basic version, 147
for the radial process, 157
of Ikeda-Watanabe, 194
Rauch, 141
complex differential form, 41
composition formula, 90
confluence of martingales, 106
conjugate locus, 130
conjugate vector, 130
connection
in principal bundle, 66
Levi-Civita, 55
metrically complete, 50
Riemannian, 54
symmetric, 51
torsion-free, 51, 52
connection form, 67
convergence theorem
of Darling-Zheng, 102
convex geometry, 107
convex map, 89
coordinates
normal, 96
cotangent bundle, 10
cotangent space
of order $k, 112$
covariant derivative
along a curve, 44
covariant derivative, 44
in coordinates, 46
in direction $v, 45$
induced, 45
of a differential form, 52
of a section along a curve, 46
covering, 143
criterion for BM-completeness, 176
curvature, 50
constant, 134
negative, 134
positive, 134
Ricci, 133
Riemann, 132
scalar, 133
sectional, 133
curvature function
radial, 150
curvature tensor, 132
curvature identities, 132
curvature tensor, 51
cut locus, 130
cut point, 130

Darling-Zheng
convergence theorem, 102
derivation, 4
derivative
covariant, 44
development
Cartan, 81
stochastic, 80, 81
diffeomorphism, 3
diffeomorphism group, 6
differentiable atlas, 1
differentiable manifold, 2
differentiable map, 2
differentiable structure, 1
differential, 4
of order $k, 112$
differential equation
second order, 50
differential form, 10
differential form of order $k, 112$
diffusion
one-dimensional, 195
dilatation, 94
direct sum of connections, 88

Dirichlet problem, 166
drift of a PDO, 120
dual connection, 88
effective action, 63
Einstein manifold, 134
embedding, 5
energy density, 89
exhaustion function, 103
exit set
of a Brownian motion, 167
exponential function, 96
exponential map, 125
Feller's test, 196
fiber, 6
at a point, 6
fiber bundle, 6
associated, 65
associated with a principal bundle, 65
basis, 6
projection, 6
total space, 6
trivial, 6
with structure group, 63
fibration, 6
local trivial, 6
first fundamental form, 89
first variation
of arc length, 124
flow
line, 10
process, 12
formula
pullback, 116
Synge, 134
Tanaka, 191
frame
induced, 38
frame bundle, 66
orthonormal, 66
free action, 63
function
exponential, 96
fundamental form
second, 89
fundamental form
first, 89
second, 89
Gauss Lemma, 126
geodesic, 46
minimal, 128
geodesic ball, 127
geodesic polar coordinates, 126
geodesic sphere, 127
geodesic spray, 50
geodesic Variation, 135
geometry
convex, 107
germ, 4
gradient, 57
group of diffeomorphisms, 6
Hörmander form, 24
Hadamard-Cartan
Theorem, 143
harmonic map
pullback properties, 97
harmonic function
mean-value property, 166
harmonic map, 89
stochastic characterization, 93
harmonic morphism, 94
analytic characterization, 95
Hessian, 52, 89
Hessian comparison theorem, 146
homomorphism
of a vector bundle, 7
Hopf-Rinow, 128
horizontal semimartingale, 74
horizontal Laplacian, 71
horizontal Brownian motion, 91
horizontal curve, 49
horizontal lift, 47, 67, 69
in a principal bundle, 69
of a semimartingale, 74
horizontal space, 67
horizontal splitting, 47
horizontally conformal, 94
hyperbolic space
hyperboloid model, 155
hyperbolic space, 155
Poincaré-model, 155
hyperboloid model of hyperbolic space, 155
immersion, 5
index form, 135
induced frame, 38
induced basis system, 38
induced covariant derivative, 45
induced fibration, 8
induced form, 37
injectivity radius, 125
integral curve, 10
inverse function theorem, 3
isometry, 54
local, 54
Itô integral
along a semimartingale, 121
of a one-form, 84
Itô process, 22
Itô's formula, 83
geometric, 83
Jacobi field
proper, 136

Jacobi field, 136
Jacobi variation, 135
Lévy's characterization, 59
Laplace operator
Euclidean, 24
Laplace operator, 24
Laplace-Beltrami operator, 58
Laplacian, 24
horizontal, 71
Laplacian comparison theorem
basic version, 147
left-invariant SDE, 76
left-invariant vector field, 76

## Lemma

Gauss, 126
length of a curve, 54
Levi-Civita connection, 55 on $\mathbb{R}^{n}$, 56
Levi-Civita parallelism, 55
Lie product, 51
linear action, 63
linear connection
on a manifold, 69
linear connection, 49
Liouville manifold, 167
local isometry, 54
local diffeomorphism, 3
local flow, 10
local frame, 7
local trivial fibration, 6
local trivialization, 6
manifold
BM-complete, 103
differentiable, 2
Einstein, 134
metrically complete, 103
Riemannian, 54
rotationally symmetric, 147
map
affine, 89
convex, 89
horizontally conformal, 94
strictly convex, 89
totally geodesic, 89
Markov property
strong, 159
martingale
$H^{p}, 105$
manifold-valued, 52
on a manifold, 52
on a submanifold, 57
one-dimensional, 86
martingale convergence, 102
maximal integral curve, 10
maximal lifetime, 20
maximal solution, 28
maximum principle, 167
mean-value property, 166
measure class
harmonic, 166
metric
Riemannian, 54
metrically complete connection, 50 metrically complete manifold, 103
minimal geodesic, 128
minimum principle
for martingales, 106
model, 147
elementary properties, 152
morphism
harmonic, 94
non-confluence of martingales, 106
normal coordinates, 126
normal bundle, 32
normal coordinates, 96
normal geodesic, 125
nullspace of the index form, 136
one-dimensional martingale, 86
one-dimensional Brownian motion, 86 one-dimensional diffusion, 195
one-dimensional SDE, 191
one-dimensional semimartingale, 86
one-form
canonical, 71
orthonormal frame bundle, 66
parallel, 44
parallel section
along a curve, 46
parallel transport, 43, 46
along a semimartingale, 83
induced by a connection, 46
parallelizable, 71
partial differential operator, 24
pathwise uniqueness
criterion of Yamada-Watanabe, 192
PDO, 24
in Hörmander form, 24
Poincaré model of hyperbolic space, 155
pole of a Riemannian manifold, 129
pre-bundle
atlas, 7
chart, 7
principal bundle
associated with a fiber bundle, 65
principal bundle, 64
principle
comparison, 142
Laurent Schwartz, 115
principle of Laurent Schwartz, 115
process
Bessel, 124
product connection, 88
pullback
of a form, 37
of a fibration, 8
pullback connection, 88
pullback formula, 40, 116
for the $b$-quadratic variation, 38
for the Stratonovich integral of a form, 39
quadratic variation, 36
Riemannian, 58
radial process
comparison theorem, 166
radial curvature, 150
radial curvature function, 150
radial geodesic, 127
radial process, 123
comparison theorem, 157
radial vector field, 149
Rauch
comparison theorem, 141
recurrence, 195
recurrence of Brownian motion, 168
reduction of the structure group, 64
representation
of a group, 63
Ricci curvature, 133
Ricci identity, 54
Riemann
curvature, 132
Riemannian manifold stochastically complete, 103
Riemannian quadratic variation, 58
Riemannian manifold
parabolic, 169
Riemannian connection, 54, 69
characterization, 54
Riemannian manifold, 54
flat, 134
hyperbolic, 169
Riemannian metric, 54
Riemannian sectional curvature, 133
Riemannian volume measure, 159
rolling without slipping, 81
rotationally symmetric manifold, 147
scalar curvature, 133
scale function, 196
SDE, 27
elliptic, 62
left-invariant, 76
maximal solution, 28
on a manifold, 27
one-dimensional, 191
solution, 27
SDE; Example of Girsanov, 193
second fundamental form, 52, 89
second variation
of arc length, 134
section, 7

```
    along a map, }
    parallel, 44
sectional curvature, }13
semimartingale
    M-valued, up to m, 101
    as solution of an SDE, 34
    on a manifold, 19
    one-dimensional, }8
    up to }\infty,10
    winding in the plane, 41
    with lifetime, 19
solution of SDE
    existence for }M=\mp@subsup{\mathbb{R}}{}{n},2
    existence for general M,30
    uniqueness for }M=\mp@subsup{\mathbb{R}}{}{n},2
    uniqueness for general M,30
space
    hyperbolic, }15
sphere, 151
spray, 50
standard-horizontal vector field, 71
stochastic development, 80,81
stochastic differential equation
    on }\mp@subsup{\mathbb{R}}{}{n},2
    on a manifold, 27
stochastic parallel transport, 83
stochastically complete manifold, 103
Stratonovich integral, 24
Stratonovich differential, 24
Stratonovich integral
    of a one-form, 38
strictly convex map, }8
strong Markov property,159
structure equations
    Cartan, 51
structure group, }6
subbundle, 6
subharmonic, }8
submanifold, 2
    totally geodesic, 109
submersive map,47
Synge formula,134
Tanaka formula, 191
tangent space
    algebraic, 4
tangent bundle, }
tangent space, 3, 4
    geometric, 3
    of order }k,11
tangent vector, 4
tangential vector field along a curve, 8
tangentially equivalent, 3
tension field, }8
tensor, 50
    curvature, 132
    of type (r,s),50
tensor field, 50
```

test function, 3
test for explosion, 196
Theorem
of Hadamard-Cartan, 143
of Hopf-Rinow, 128
of Yamada-Watanabe, 192

## theorem

Levi-Civita, 55
time-change, 19
topological manifold, 1
torsion, 50
torsion tensor, 51
torsion-free, 51, 52
totally geodesic, 109
totally geodesic map, 89
transience, 195
transience of Brownian motion, 168
transition function, 6
transition map, 1
trivialization, 6
local, 6
typical fiber, 6
variation
$b$-quadratic, 36
free, 124
of a curve, 124
variation of a curve, 124
vector bundle, 6
induced, 8
pullback, 8
vector field, 9
along a map, 8
in coordinates, 10
left-invariant, 76
radial, 149
standard-vertical, 67
vector field of order $k, 112$
velocity vector, 6
vertical space, 67
volume measure
Riemannian, 159
Whitney embedding, 30
Whitney's embedding theorem, 30

## Notations

```
( \(M, g\) ) Riemannian manifold, 54
\(B_{r}(x)\) (geodesic ball about \(x\) of radius \(r\) ), 127
\(C^{\infty}(M ; N)\) (space of differentiable maps), 3
\(L(\alpha)\) (length of a curve), 54
\(R\) (Riemann curvature), 132
\(S_{r}(x)\) (geodesic sphere about \(x\) of radius \(r\) ), 127
\(T_{x} M\) (tangent space), 4
[ \(X, X]\) (Riemannian quadratic variation), 58
\(\Delta\) (Laplace-Beltrami operator), 58
\(\Gamma(E)\) (sections of a vector bundle \(E\) ), 7
\(\Gamma\left(f^{*} T M\right)\) (vector fields along a map), 8
\(\Gamma_{i j}^{k}\) (Christoffel symbols), 45
\(\mathrm{L}(T M)\) (frame bundle), 66
\(\mathrm{O}(T M)\) (orthonormal frame bundle), 66
\(\operatorname{Ric}^{M}\) (Ricci curvature), 133
Riem \({ }^{M}\) (sectional curvature), 133
\(\mathscr{A}\) (processes locally of bounded variation), 35
\(\mathscr{A}(X)\) (anti-development of \(X\) ), 74
\(\mathscr{M}\) (space of real local martingales), 35
\(\mathscr{S}\) (space of real semimartingales), 35
\(\operatorname{cut}(x)\) (cut locus), 130
\(\dot{\alpha}\) (tangential vector field), 8
\(\dot{\alpha}\) (velocity field along a curve), 6
\(\int b(d X, d X)\) ( \(b\)-quadratic variation), 36
\(\int_{X} \alpha\) (Stratonovich integral of a one-form, 38
\(\nabla_{D} \sigma\) (covariant derivative along a curve, 44
\(\nabla d f\) (Hesse form), 52
\(\nabla\) (covariant derivative), 44
\(\frac{\partial}{\partial h^{i}}\) (coordinate basis field), 5
\(\frac{\partial}{\partial r}\) (radial vector field), 149
\(\tau(f)\) (tension), 89
\(\mathrm{BM}(M, g)\) (Brownian motions on \((M, g)\) ), 59
Conj \((x)\) (conjugate locus), 130
vol (Riemannian volume measure), 159
\(d(x, y)\) (distance), 123
\(d f\) (differential of \(f\) ), 8
\(d f_{x}\) (differential of \(f\) at \(x\) ), 4
\(k^{M}\) (scalar curvature), 133
\(k_{\mathbb{M}}\) (radial curvature function), 150
```


[^0]:    This is a very preliminary draft of an introduction to Stochastic Differential Geometry. The first chapter develops notions of Stochastic Analysis on Manifolds and is based on a lecture course given at the University of Luxembourg. The text will be continued in parallel to the teaching of courses on Stochastic Riemannian Geometry at the University of Luxembourg.

