## 1 Linear Ordinary Differential Equations

The fundamental theorem for ordinary differential equations. Consider a function

$$
\begin{equation*}
\boldsymbol{f}: I \times U \ni(t, \boldsymbol{x}) \longmapsto \boldsymbol{f}(t, \boldsymbol{x}) \in \mathbb{R}^{m} \tag{1.1}
\end{equation*}
$$

of class $C^{1}$, where $I \subset \mathbb{R}$ is an interval and $U \subset \mathbb{R}^{m}$ is a domain in the Euclidean space $\mathbb{R}^{m}$. For any fixed $t_{0} \in I$ and $\boldsymbol{x}_{0} \in U$, the condition

$$
\begin{equation*}
\frac{d}{d t} \boldsymbol{x}(t)=\boldsymbol{f}(t, \boldsymbol{x}(t)), \quad \boldsymbol{x}\left(t_{0}\right)=\boldsymbol{x}_{0} \tag{1.2}
\end{equation*}
$$

of an $\mathbb{R}^{m}$-valued function $t \mapsto \boldsymbol{x}(t)$ is called the initial value problem of ordinary differential equation for unknown function $\boldsymbol{x}(t)$. A function $\boldsymbol{x}: I \rightarrow U$ satisfying (1.2) is called a solution of the initial value problem.
Fact 1.1 (The existence theorem for ODE's). Let $\boldsymbol{f}: I \times U \rightarrow \mathbb{R}^{m}$ be a $C^{1}$-function as in (1.1). Then, for any $\boldsymbol{x}_{0} \in U$ and $t_{0} \in I$, there exists a positive number $\varepsilon$ and a $C^{1}$-function $\boldsymbol{x}: I \cap\left(t_{0}-\right.$ $\left.\varepsilon, t_{0}+\varepsilon\right) \rightarrow U$ satisfying (1.2).

Consider two solutions $\boldsymbol{x}_{j}: J_{j} \rightarrow U(j=1,2)$ of (1.2) defined on subintervals $J_{j} \subset I$ containing $t_{0}$. Then the function $\boldsymbol{x}_{2}$ is said to be an extension of $\boldsymbol{x}_{1}$ if $J_{1} \subset J_{2}$ and $\left.\boldsymbol{x}_{2}\right|_{J_{1}}=\boldsymbol{x}_{1}$. A solution $\boldsymbol{x}$ of (1.2) is said to be maximal if there are no non-trivial extension of it.
Fact 1.2 (The uniqueness for ODE's). The maximal solution of (1.2) is unique.
Fact 1.3 (Smoothness of the solutions). If $\boldsymbol{f}: I \times U \rightarrow \mathbb{R}^{m}$ is of class $C^{r}(r=1, \ldots, \infty)$, the solution of (1.2) is of class $C^{r+1}$. Here, $\infty+1=\infty$, as a convention.

Let $V \subset \mathbb{R}^{k}$ be another domain of $\mathbb{R}^{k}$ and consider a $C^{\infty}$-function

$$
\begin{equation*}
\boldsymbol{h}: I \times U \times V \ni(t, \boldsymbol{x} ; \boldsymbol{\alpha}) \mapsto \boldsymbol{h}(t, \boldsymbol{x} ; \boldsymbol{\alpha}) \in \mathbb{R}^{m} \tag{1.3}
\end{equation*}
$$

For fixed $t_{0} \in I$, we denote by $\boldsymbol{x}\left(t ; \boldsymbol{x}_{0}, \boldsymbol{\alpha}\right)$ the (unique, maximal) solution of (1.2) for $\boldsymbol{f}(t, \boldsymbol{x})=$ $\boldsymbol{h}(t, \boldsymbol{x} ; \boldsymbol{\alpha})$. Then
Fact 1.4. The map $\left(t, \boldsymbol{x}_{0} ; \boldsymbol{\alpha}\right) \mapsto \boldsymbol{x}\left(t ; \boldsymbol{x}_{0}, \boldsymbol{\alpha}\right)$ is of class $C^{\infty}$.
Example 1.5. (1) Let $m=1, I=\mathbb{R}, U=\mathbb{R}$ and $f(t, x)=\lambda x$, where $\lambda$ is a constant. Then $x(t)=x_{0} \exp (\lambda t)$ defined on $\mathbb{R}$ is the maximal solution to

$$
\frac{d}{d t} x(t)=f(t, x(t))=\lambda x(t), \quad x(0)=x_{0}
$$

(2) Let $m=2, I=\mathbb{R}, U=\mathbb{R}^{2}$ and $\boldsymbol{f}(t ;(x, y))=\left(y,-\omega^{2} x\right)$, where $\omega$ is a constant. Then

$$
\binom{x(t)}{y(t)}=\binom{x_{0} \cos \omega t+\frac{y_{0}}{\omega} \sin \omega t}{-x_{0} \omega \sin \omega t+y_{0} \cos \omega t}
$$

is the unique solution of

$$
\frac{d}{d t}\binom{x(t)}{y(t)}=\binom{y(t)}{-\omega^{2} x(t)}, \quad\binom{x(0)}{y(0)}=\binom{x_{0}}{y_{0}}
$$

defined on $\mathbb{R}$. This differential equation can be considered a single equation

$$
\frac{d^{2}}{d t^{2}} x(t)=-\omega^{2} x(t), \quad x(0)=x_{0}, \quad \frac{d x}{d t}(0)=y_{0}
$$

of order 2 .
(3) Let $m=1, I=\mathbb{R}, U=\mathbb{R}$ and $f(t, x)=1+x^{2}$. Then $x(t)=\tan t$ defined on $\left(-\frac{\pi}{2}, \frac{\pi}{2}\right)$ is the unique maximal solution of the initial value problem

$$
\frac{d x}{d t}=1+x^{2}, \quad x(0)=0
$$

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Linear Ordinary Differential Equations. The ordinary differential equation (1.2) is said to be linear if the function (1.1) is a linear function in $\boldsymbol{x}$, that is, a linear differential equation is in a form

$$
\frac{d}{d t} \boldsymbol{x}(t)=A(t) \boldsymbol{x}(t)+\boldsymbol{b}(t)
$$

where $A(t)$ and $\boldsymbol{b}(t)$ are $m \times m$-matrix-valued and $\mathbb{R}^{m}$-valued functions in $t$.
For the sake of later use, we consider, in this lecture, the special form of linear differential equation for matrix-valued unknown functions as follows: Let $\mathrm{M}_{n}(\mathbb{R})$ be the set of $n \times n$-matrices with real components, and take functions

$$
\Omega: I \longrightarrow \mathrm{M}_{n}(\mathbb{R}), \quad \text { and } B: I \longrightarrow \mathrm{M}_{n}(\mathbb{R})
$$

where $I \subset \mathbb{R}$ is an interval. Identifying $\mathrm{M}_{n}(\mathbb{R})$ with $\mathbb{R}^{n^{2}}$, we assume $\Omega$ and $B$ are continuous functions (with respect to the topology of $\mathbb{R}^{n^{2}}=\mathrm{M}_{n}(\mathbb{R})$ ). Then we can consider the linear ordinary differential equation for matrix-valued unknown $X(t)$ as

$$
\begin{equation*}
\frac{d X(t)}{d t}=X(t) \Omega(t)+B(t), \quad X\left(t_{0}\right)=X_{0} \tag{1.4}
\end{equation*}
$$

where $X_{0}$ is given constant matrix.
Then, the fundamental theorem of linear ordinary equation states that the maximal solution of (1.4) is defined on whole $I$. To prove this, we prepare some materials related to matrix-valued functions.

Preliminaries: Matrix Norms. Denote by $\mathrm{M}_{n}(\mathbb{R})$ the set of $n \times n$-matrices with real components, which can be identified the vector space $\mathbb{R}^{n^{2}}$. In particular, the Euclidean norm of $\mathbb{R}^{n^{2}}$ induces a norm

$$
\begin{equation*}
|X|_{\mathrm{E}}=\sqrt{\operatorname{tr}\left(X^{T} X\right)}=\sqrt{\sum_{i, j=1}^{n} x_{i j}^{2}} \tag{1.5}
\end{equation*}
$$

on $\mathrm{M}_{n}(\mathbb{R})$. On the other hand, we let

$$
\begin{equation*}
|X|_{\mathrm{M}}:=\sup \left\{\frac{|X \boldsymbol{v}|}{|\boldsymbol{v}|} ; \boldsymbol{v} \in \mathbb{R}^{n} \backslash\{\mathbf{0}\}\right\} \tag{1.6}
\end{equation*}
$$

where $|\cdot|$ denotes the Euclidean norm of $\mathbb{R}^{n}$.
Lemma 1.6. (1) The map $X \mapsto|X|_{\mathrm{M}}$ is a norm of $\mathrm{M}_{n}(\mathbb{R})$.
(2) For $X, Y \in \mathrm{M}_{n}(\mathbb{R})$, it holds that $|X Y|_{\mathrm{M}} \leqq|X|_{\mathrm{M}}|Y|_{\mathrm{M}}$.
(3) Let $\lambda=\lambda(X)$ be the maximum eigenvalue of semi-positive definite symmetric matrix $X^{T} X$. Then $|X|_{\mathrm{M}}=\sqrt{\lambda}$ holds.
(4) $(1 / \sqrt{n})|X|_{\mathrm{E}} \leqq|X|_{\mathrm{M}} \leqq|X|_{\mathrm{E}}$.
(5) The map $|\cdot|_{\mathrm{M}}: \mathrm{M}_{n}(\mathbb{R}) \rightarrow \mathbb{R}$ is continuous with respect to the Euclidean norm.

Proof. Since $|X \boldsymbol{v}| /|\boldsymbol{v}|$ is invariant under scalar multiplications to $\boldsymbol{v}$, we have $|X|_{\mathrm{M}}=\sup \{|X \boldsymbol{v}| ; \boldsymbol{v} \in$ $\left.S^{n-1}\right\}$, where $S^{n-1}$ is the unit sphere in $\mathbb{R}^{n}$. Since $S^{n-1} \ni \boldsymbol{x} \mapsto|A \boldsymbol{x}| \in \mathbb{R}$ is a continuous function defined on a compact space, it takes the maximum. Thus, the right-hand side of (1.6) is welldefined. It is easy to verify that $|\cdot|_{\mathrm{M}}$ satisfies the axiom of the norm ${ }^{1}$.

[^0]Since $A:=X^{T} X$ is positive semi-definite, the eigenvalues $\lambda_{j}(j=1, \ldots, n)$ are non-negative real numbers. In particular, there exists an orthonormal basis [ $\boldsymbol{a}_{j}$ ] of $\mathbb{R}^{n}$ satisfying $A \boldsymbol{a}_{j}=\lambda_{j} \boldsymbol{a}_{j}$ $(j=1, \ldots, n)$. Let $\lambda$ be the maximum eigenvalue of $A$, and write $\boldsymbol{v}=v_{1} \boldsymbol{a}_{1}+\cdots+v_{n} \boldsymbol{a}_{n}$. Then it holds that

$$
\langle X \boldsymbol{v}, X \boldsymbol{v}\rangle=\lambda_{1} v_{1}^{2}+\cdots+\lambda_{n} v_{n}^{2} \leqq \lambda\langle\boldsymbol{v}, \boldsymbol{v}\rangle,
$$

where $\langle$,$\rangle is the Euclidean inner product of \mathbb{R}^{n}$. The equality of this inequality holds if and only if $\boldsymbol{v}$ is the $\lambda$-eigenvector, proving (3). Noticing the norm (1.5) is invariant under conjugations $X \mapsto$ $P^{T} X P(P \in \mathrm{O}(n))$, we obtain $|X|_{\mathrm{E}}=\sqrt{\lambda_{1}^{2}+\cdots+\lambda_{n}^{2}}$ by diagonalizing $X^{T} X$ by an orthogonal matrix $P$. Then we obtain (4). Hence two norms $|\cdot|_{\mathrm{E}}$ and $|\cdot|_{\mathrm{M}}$ induce the same topology as $\mathrm{M}_{n}(\mathbb{R})$. In particular, we have (5).

## Preliminaries: Matrix-valued Functions.

Lemma 1.7. Let $X$ and $Y$ be $C^{\infty}$-maps defined on a domain $U \subset \mathbb{R}^{m}$ into $\mathrm{M}_{n}(\mathbb{R})$. Then
(1) $\frac{\partial}{\partial u_{j}}(X Y)=\frac{\partial X}{\partial u_{j}} Y+X \frac{\partial Y}{\partial u_{j}}$,
(2) $\frac{\partial}{\partial u_{j}} \operatorname{det} X=\operatorname{tr}\left(\widetilde{X} \frac{\partial X}{\partial u_{j}}\right)$, and
(3) $\frac{\partial}{\partial u_{j}} X^{-1}=-X^{-1} \frac{\partial X}{\partial u_{j}} X^{-1}$,
where $\widetilde{X}$ is the cofactor matrix of $X$, and we assume in (3) that $X$ is a regular matrix.
Proof. The formula (1) holds because the definition of matrix multiplication and the Leibnitz rule, Denoting ${ }^{\prime}=\partial / \partial u_{j}$,

$$
O=(\mathrm{id})^{\prime}=\left(X^{-1} X\right)^{\prime}=\left(X^{-1}\right) X^{\prime}+\left(X^{-1}\right)^{\prime} X
$$

implies (3), where id is the identity matrix.
Decompose the matrix $X$ into column vectors as $X=\left(\boldsymbol{x}_{1}, \ldots, \boldsymbol{x}_{n}\right)$. Since the determinant is multi-linear form for $n$-tuple of column vectors, it holds that

$$
(\operatorname{det} X)^{\prime}=\operatorname{det}\left(\boldsymbol{x}_{1}^{\prime}, \boldsymbol{x}_{2}, \ldots, \boldsymbol{x}_{n}\right)+\operatorname{det}\left(\boldsymbol{x}_{1}, \boldsymbol{x}_{2}^{\prime}, \ldots, \boldsymbol{x}_{n}\right)+\cdots+\operatorname{det}\left(\boldsymbol{x}_{1}, \boldsymbol{x}_{2}, \ldots, \boldsymbol{x}_{n}^{\prime}\right)
$$

Then by cofactor expansion of the right-hand side, we obtain (2).
Proposition 1.8. Assume two $C^{\infty}$ matrix-valued functions $X(t)$ and $\Omega(t)$ satisfy

$$
\begin{equation*}
\frac{d X(t)}{d t}=X(t) \Omega(t), \quad X\left(t_{0}\right)=X_{0} \tag{1.7}
\end{equation*}
$$

Then

$$
\begin{equation*}
\operatorname{det} X(t)=\left(\operatorname{det} X_{0}\right) \exp \int_{t_{0}}^{t} \operatorname{tr} \Omega(\tau) d \tau \tag{1.8}
\end{equation*}
$$

holds. In particular, if $X_{0} \in \mathrm{GL}(n, \mathbb{R}),{ }^{2}$ then $X(t) \in \mathrm{GL}(n, \mathbb{R})$ for all $t$.
Proof. By (2) of Lemma 1.7, we have

$$
\begin{aligned}
\frac{d}{d t} \operatorname{det} X(t) & =\operatorname{tr}\left(\widetilde{X}(t) \frac{d X(t)}{d t}\right)=\operatorname{tr}(\widetilde{X}(t) X(t) \Omega(t)) \\
& =\operatorname{tr}(\operatorname{det} X(t) \Omega(t))=\operatorname{det} X(t) \operatorname{tr} \Omega(t)
\end{aligned}
$$

Here, we used the relation $\widetilde{X} X=X \widetilde{X}=(\operatorname{det} X) \operatorname{id}^{3}$. Hence $\frac{d}{d t}\left(\rho(t)^{-1} \operatorname{det} X(t)\right)=0$, where $\rho(t)$ is the right-hand side of (1.8).

[^1]Corollary 1.9. If $\Omega(t)$ in (1.7) satisfies $\operatorname{tr} \Omega(t)=0$, $\operatorname{det} X(t)$ is constant. In particular, if $X_{0} \in \operatorname{SL}(n, \mathbb{R}), X$ is a function valued in $\operatorname{SL}(n, \mathbb{R})^{4}$.
Proposition 1.10. Assume $\Omega(t)$ in (1.7) is skew-symmetric for all $t$, that is, $\Omega^{T}+\Omega$ is identically O. If $X_{0} \in \mathrm{O}(n)$ (resp. $\left.X_{0} \in \mathrm{SO}(n)\right)^{5}$, then $X(t) \in \mathrm{O}(n)$ (resp. $\left.X(t) \in \mathrm{SO}(n)\right)$ for all $t$.

Proof. By (1) in Lemma 1.7,

$$
\begin{aligned}
\frac{d}{d t}\left(X X^{T}\right) & =\frac{d X}{d t} X^{T}+X\left(\frac{d X}{d t}\right)^{T} \\
& =X \Omega X^{T}+X \Omega^{T} X^{T}=X\left(\Omega+\Omega^{T}\right) X^{T}=O
\end{aligned}
$$

Hence $X X^{T}$ is constant, that is, if $X_{0} \in \mathrm{O}(n)$,

$$
X(t) X(t)^{T}=X\left(t_{0}\right) X\left(t_{0}\right)^{T}=X_{0} X_{0}^{T}=\mathrm{id}
$$

If $X_{0} \in \mathrm{O}(n)$, this proves the first case of the proposition. Since $\operatorname{det} A= \pm 1$ when $A \in \mathrm{O}(n)$, the second case follows by continuity of $\operatorname{det} X(t)$.

Preliminaries: Norms of Matrix-Valued functions. Let $I=[a, b]$ be a closed interval, and denote by $C^{0}\left(I, \mathrm{M}_{n}(\mathbb{R})\right)$ the set of continuous functions $X: I \rightarrow \mathrm{M}_{n}(\mathbb{R})$. For any positive number $k$, we define

$$
\begin{equation*}
\|X\|_{I, k}:=\sup \left\{e^{-k t}|X(t)|_{\mathrm{M}} ; t \in I\right\} \tag{1.9}
\end{equation*}
$$

for $X \in C^{0}\left(I, \mathrm{M}_{n}(\mathbb{R})\right)$. When $k=0,\|\cdot\|_{I, 0}$ is the uniform norm for continuous functions, which is complete. Similarly, one can prove the following in the same way:

Lemma 1.11. The norm $\|\cdot\|_{I, k}$ on $C^{0}\left(I, \mathrm{M}_{n}(\mathbb{R})\right)$ is complete.
Linear Ordinary Differential Equations. We prove the fundamental theorem for linear ordinary differential equations.

Proposition 1.12. Let $\Omega(t)$ be a $C^{\infty}$-function valued in $\mathrm{M}_{n}(\mathbb{R})$ defined on an interval $I$. Then for each $t_{0} \in I$, there exists the unique matrix-valued $C^{\infty}$-function $X(t)=X_{t_{0}, \text { id }}(t)$ such that

$$
\begin{equation*}
\frac{d X(t)}{d t}=X(t) \Omega(t), \quad X\left(t_{0}\right)=\mathrm{id} \tag{1.10}
\end{equation*}
$$

Proof. Uniqueness: Assume $X(t)$ and $Y(t)$ satisfy (1.10). Then

$$
Y(t)-X(t)=\int_{t_{0}}^{t}\left(Y^{\prime}(\tau)-X^{\prime}(\tau)\right) d \tau=\int_{t_{0}}^{t}(Y(\tau)-X(\tau)) \Omega(\tau) d \tau \quad\left({ }^{\prime}=\frac{d}{d t}\right)
$$

holds. Hence for an arbitrary closed interval $J \subset I$,

[^2]holds for $t \in J$. Thus, for an appropriate choice of $k \in \mathbb{R}$, it holds that
$$
\|Y-X\|_{J, k} \leqq \frac{1}{2}\|Y-X\|_{J, k},
$$
that is, $\|Y-X\|_{J, k}=0$, proving $Y(t)=X(t)$ for $t \in J$. Since $J$ is arbitrary, $Y=X$ holds on $I$. Existence: Let $J:=\left[t_{0}, a\right] \subset I$ be a closed interval, and define a sequence $\left\{X_{j}\right\}$ of matrix-valued functions defined on $I$ satisfying $X_{0}(t)=\mathrm{id}$ and
\[

$$
\begin{equation*}
X_{j+1}(t)=\mathrm{id}+\int_{t_{0}}^{t} X_{j}(\tau) \Omega(\tau) d \tau \quad(j=0,1,2, \ldots) . \tag{1.11}
\end{equation*}
$$

\]

Let $k:=2 \sup _{J}|\Omega|_{\mathrm{M}}$. Then

$$
\begin{aligned}
& \left|X_{j+1}(t)-X_{j}(t)\right|_{\mathrm{M}} \leqq \int_{t_{0}}^{t}\left|X_{j}(\tau)-X_{j-1}(\tau)\right|_{\mathrm{M}}|\Omega(\tau)|_{\mathrm{M}} d \tau \\
& \quad \leqq \frac{e^{k\left(t-t_{0}\right)}}{|k|} \sup _{J}|\Omega|_{\mathrm{M}}| | X_{j}-\left.X_{j-1}\right|_{J, k}
\end{aligned}
$$

for an appropriate choice of $k \in \mathbb{R}$, and hence $\left\|X_{j+1}-X_{j}\right\|_{J, k} \leqq \frac{1}{2}\left\|X_{j}-X_{j-1}\right\|_{J, k}$, that is, $\left\{X_{j}\right\}$ is a Cauchy sequence with respect to $\|\cdot\|_{J, k}$. Thus, by completeness (Lemma 1.11), it converges to some $X \in C^{0}\left(J, \mathrm{M}_{n}(\mathbb{R})\right)$. By (1.11), the limit $X$ satisfies

$$
X\left(t_{0}\right)=\mathrm{id}, \quad X(t)=\mathrm{id}+\int_{t_{0}}^{t} X(\tau) \Omega(\tau) d \tau .
$$

Applying the fundamental theorem of calculus, we can see that $X$ satisfies $X^{\prime}(t)=X(t) \Omega(t)$ $\left(^{\prime}=d / d t\right)$. Since $J$ can be taken arbitrarily, existence of the solution on $I$ is proven.

Finally, we shall prove that $X$ is of class $C^{\infty}$. Since $X^{\prime}(t)=X(t) \Omega(t)$, the derivative $X^{\prime}$ of $X$ is continuous. Hence $X$ is of class $C^{1}$, and so is $X(t) \Omega(t)$. Thus we have that $X^{\prime}(t)$ is of class $C^{1}$, and then $X$ is of class $C^{2}$. Iterating this argument, we can prove that $X(t)$ is of class $C^{r}$ for arbitrary $r$.

Corollary 1.13. Let $\Omega(t)$ be a matrix-valued $C^{\infty}$-function defined on an interval $I$. Then for each $t_{0} \in I$ and $X_{0} \in \mathrm{M}_{n}(\mathbb{R})$, there exists the unique matrix-valued $C^{\infty}$-function $X_{t_{0}, X_{0}}(t)$ defined on I such that

$$
\begin{equation*}
\frac{d X(t)}{d t}=X(t) \Omega(t), \quad X\left(t_{0}\right)=X_{0} \quad\left(X(t):=X_{t_{0}, X_{0}}(t)\right) \tag{1.12}
\end{equation*}
$$

In particular, $X_{t_{0}, X_{0}}(t)$ is of class $C^{\infty}$ in $X_{0}$ and $t$.
Proof. We rewrite $X(t)$ in Proposition 1.12 as $Y(t)=X_{t_{0}, \mathrm{id}}(t)$. Then the function

$$
\begin{equation*}
X(t):=X_{0} Y(t)=X_{0} X_{t_{0}, \mathrm{id}}(t), \tag{1.13}
\end{equation*}
$$

is desired one. Conversely, assume $X(t)$ satisfies the conclusion. Noticing $Y(t)$ is a regular matrix for all $t$ because of Proposition 1.8,

$$
W(t):=X(t) Y(t)^{-1}
$$

satisfies

$$
\frac{d W}{d t}=\frac{d X}{d t} Y^{-1}-X Y^{-1} \frac{d Y}{d t} Y^{-1}=X \Omega Y^{-1}-X Y^{-1} Y \Omega Y^{-1}=O
$$

Hence

$$
W(t)=W\left(t_{0}\right)=X\left(t_{0}\right) Y\left(t_{0}\right)^{-1}=X_{0} .
$$

Hence the uniqueness is obtained. The final part is obvious by the expression (1.13).

Proposition 1.14. Let $\Omega(t)$ and $B(t)$ be matrix-valued $C^{\infty}$-functions defined on $I$. Then for each $t_{0} \in I$ and $X_{0} \in \mathrm{M}_{n}(\mathbb{R})$, there exists the unique matrix-valued $C^{\infty}$-function defined on $I$ satisfying

$$
\begin{equation*}
\frac{d X(t)}{d t}=X(t) \Omega(t)+B(t), \quad X\left(t_{0}\right)=X_{0} \tag{1.14}
\end{equation*}
$$

Proof. Rewrite $X$ in Proposition 1.12 as $Y:=X_{t_{0}, \text { id }}$. Then

$$
\begin{equation*}
X(t)=\left(X_{0}+\int_{t_{0}}^{t} B(\tau) Y^{-1}(\tau) d \tau\right) Y(t) \tag{1.15}
\end{equation*}
$$

satisfies (1.14). Conversely, if $X$ satisfies (1.14), $W:=X Y^{-1}$ satisfies

$$
X^{\prime}=W^{\prime} Y+W Y^{\prime}=W^{\prime} Y+W Y \Omega, \quad X \Omega+B=W Y \Omega+B
$$

and then we have $W^{\prime}=B Y^{-1}$. Since $W\left(t_{0}\right)=X_{0}$,

$$
W=X_{0}+\int_{t_{0}}^{t} B(\tau) Y^{-1}(\tau) d \tau
$$

Thus we obtain (1.15).
Theorem 1.15. Let $I$ and $U$ be an interval and a domain in $\mathbb{R}^{m}$, respectively, and let $\Omega(t, \boldsymbol{\alpha})$ and $B(t, \boldsymbol{\alpha})$ be matrix-valued $C^{\infty}$-functions defined on $I \times U\left(\boldsymbol{\alpha}=\left(\alpha_{1}, \ldots, \alpha_{m}\right)\right)$. Then for each $t_{0} \in I$, $\boldsymbol{\alpha} \in U$ and $X_{0} \in \mathrm{M}_{n}(\mathbb{R})$, there exists the unique matrix-valued $C^{\infty}$-function $X(t)=X_{t_{0}, X_{0}, \boldsymbol{\alpha}}(t)$ defined on $I$ such that

$$
\begin{equation*}
\frac{d X(t)}{d t}=X(t) \Omega(t, \boldsymbol{\alpha})+B(t, \boldsymbol{\alpha}), \quad X\left(t_{0}\right)=X_{0} \tag{1.16}
\end{equation*}
$$

Moreover,

$$
I \times I \times \mathrm{M}_{n}(\mathbb{R}) \times U \ni\left(t, t_{0}, X_{0}, \boldsymbol{\alpha}\right) \mapsto X_{t_{0}, X_{0}, \boldsymbol{\alpha}}(t) \in \mathrm{M}_{n}(\mathbb{R})
$$

is a $C^{\infty}$-map.
Proof. Let $\widetilde{\Omega}(t, \tilde{\boldsymbol{\alpha}}):=\Omega\left(t+t_{0}, \boldsymbol{\alpha}\right)$ and $\widetilde{B}(t, \tilde{\boldsymbol{\alpha}})=B\left(t+t_{0}, \boldsymbol{\alpha}\right)$, and let $\widetilde{X}(t):=X\left(t+t_{0}\right)$. Then (1.16) is equivalent to

$$
\begin{equation*}
\frac{d \widetilde{X}(t)}{d t}=\widetilde{X}(t) \widetilde{\Omega}(t, \tilde{\boldsymbol{\alpha}})+\widetilde{B}(t, \tilde{\boldsymbol{\alpha}}), \quad \widetilde{X}(0)=X_{0} \tag{1.17}
\end{equation*}
$$

where $\tilde{\boldsymbol{\alpha}}:=\left(t_{0}, \alpha_{1}, \ldots, \alpha_{m}\right)$. There exists the unique solution $\widetilde{X}(t)=\widetilde{X}_{0, X_{0}, \tilde{\boldsymbol{\alpha}}}(t)$ of (1.17) for each $\tilde{\boldsymbol{\alpha}}$ because of Proposition 1.14. So it is sufficient to show differentiability with respect to the parameter $\tilde{\boldsymbol{\alpha}}$. We set $Z=Z(t)$ the unique solution of

$$
\begin{equation*}
\frac{d Z}{d t}=Z \widetilde{\Omega}+\widetilde{X} \frac{\partial \widetilde{\Omega}}{\partial \alpha_{j}}+\frac{\partial \widetilde{B}}{\partial \alpha_{j}}, \quad Z(0)=O \tag{1.18}
\end{equation*}
$$

Then it holds that $Z=\partial \widetilde{X} / \partial \alpha_{j}$. In particular, by the proof of Proposition 1.14, it holds that

$$
Z=\frac{\partial \widetilde{X}}{\partial \alpha_{j}}=\left(\int_{0}^{t}\left(\widetilde{X}(\tau) \frac{\partial \widetilde{\Omega}(\tau, \tilde{\boldsymbol{\alpha}})}{\partial \alpha_{j}}+\frac{\partial \widetilde{B}(\tau, \tilde{\boldsymbol{\alpha}})}{\partial \alpha_{j}}\right) Y^{-1}(\tau) d \tau\right) Y(t)
$$

Here, $Y(t)$ is the unique matrix-valued $C^{\infty}$-function satisfying $Y^{\prime}(t)=Y(t) \widetilde{\Omega}(t, \widetilde{\boldsymbol{\alpha}})$, and $Y(0)=$ id. Hence $\tilde{X}$ is a $C^{\infty}$-function in $(t, \tilde{\boldsymbol{\alpha}})$.

Fundamental Theorem for Space Curves. As an application, we prove the fundamental theorem for space curves. A $C^{\infty}$-map $\gamma: I \rightarrow \mathbb{R}^{3}$ defined on an interval $I \subset \mathbb{R}$ into $\mathbb{R}^{3}$ is said to be a regular curve if $\dot{\gamma} \neq \mathbf{0}$ holds on $I$. For a regular curve $\gamma(t)$, there exists a parameter change $t=t(s)$ such that $\tilde{\gamma}(s):=\gamma(t(s))$ satisfies $\left|\tilde{\gamma}^{\prime}(s)\right|=1$. Such a parameter $s$ is called the arc-length parameter.

Let $\gamma(s)$ be a regular curve in $\mathbb{R}^{3}$ parametrized by the arc-length satisfying $\gamma^{\prime \prime}(s) \neq \mathbf{0}$ for all $s$. Then

$$
\boldsymbol{e}(s):=\gamma^{\prime}(s), \quad \boldsymbol{n}(s):=\frac{\gamma^{\prime \prime}(s)}{\left|\gamma^{\prime \prime}(s)\right|}, \quad \boldsymbol{b}(s):=\boldsymbol{e}(s) \times \boldsymbol{n}(s)
$$

forms a positively oriented orthonormal basis $\{\boldsymbol{e}, \boldsymbol{n}, \boldsymbol{b}\}$ of $\mathbb{R}^{3}$ for each $s$. Regarding each vector as column vector, we have the matrix-valued function

$$
\begin{equation*}
\mathcal{F}(s):=(\boldsymbol{e}(s), \boldsymbol{n}(s), \boldsymbol{b}(s)) \in \mathrm{SO}(3) \tag{1.19}
\end{equation*}
$$

in $s$, which is called the Frenet frame associated to the curve $\gamma$. Under the situation above, we set

$$
\kappa(s):=\left|\gamma^{\prime \prime}(s)\right|>0, \quad \tau(s):=-\left\langle\boldsymbol{b}^{\prime}(s), \boldsymbol{n}(s)\right\rangle
$$

which are called the curvature and torsion, respectively, of $\gamma$. Using these quantities, the Frenet frame satisfies

$$
\frac{d \mathcal{F}}{d s}=\mathcal{F} \Omega, \quad \Omega=\left(\begin{array}{ccc}
0 & -\kappa & 0  \tag{1.20}\\
\kappa & 0 & -\tau \\
0 & \tau & 0
\end{array}\right)
$$

Proposition 1.16. The curvature and the torsion are invariant under the transformation $\boldsymbol{x} \mapsto$ $A \boldsymbol{x}+\boldsymbol{b}$ of $\mathbb{R}^{3}\left(A \in \mathrm{SO}(3), \boldsymbol{b} \in \mathbb{R}^{3}\right)$. Conversely, two curves $\gamma_{1}(s)$, $\gamma_{2}(s)$ parametrized by arclength parameter have common curvature and torsion, there exist $A \in \mathrm{SO}(3)$ and $\boldsymbol{b} \in \mathbb{R}^{3}$ such that $\gamma_{2}=A \gamma_{1}+\boldsymbol{b}$.

Proof. Let $\kappa, \tau$ and $\mathcal{F}_{1}$ be the curvature, torsion and the Frenet frame of $\gamma_{1}$, respectively. Then the Frenet frame of $\gamma_{2}=A \gamma_{1}+\boldsymbol{b}\left(A \in \mathrm{SO}(3), \boldsymbol{b} \in \mathbb{R}^{3}\right)$ is $\mathcal{F}_{2}=A \mathcal{F}_{1}$. Hence both $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$ satisfy (1.20), and then $\gamma_{1}$ and $\gamma_{2}$ have common curvature and torsion.

Conversely, assume $\gamma_{1}$ and $\gamma_{2}$ have common curvature and torsion. Then the frenet frame $\mathcal{F}_{1}$, $\mathcal{F}_{2}$ both satisfy (1.20). Let $\mathcal{F}$ be the unique solution of (1.20) with $\mathcal{F}\left(t_{0}\right)=$ id. Then by the proof of Corollary 1.13, we have $\mathcal{F}_{j}(t)=\mathcal{F}_{j}\left(t_{0}\right) \mathcal{F}(t)(j=1,2)$. In particular, since $\mathcal{F}_{j} \in \operatorname{SO}(3)$, $\mathcal{F}_{2}(t)=A \mathcal{F}_{1}(t)\left(A:=\mathcal{F}_{2}\left(t_{0}\right) \mathcal{F}_{1}\left(t_{0}\right)^{-1} \in \mathrm{SO}(3)\right)$. Comparing the first column of these, $\gamma_{2}^{\prime}(s)=$ $A \gamma_{1}^{\prime}(t)$ holds. Integrating this, the conclusion follows.

Theorem 1.17 (The fundamental theorem for space curves).
Let $\kappa(s)$ and $\tau(s)$ be $C^{\infty}$-functions defined on an interval I satisfying $\kappa(s)>0$ on $I$. Then there exists a space curve $\gamma(s)$ parametrized by arc-length whose curvature and torsion are $\kappa$ and $\tau$, respectively. Moreover, such a curve is unique up to transformation $\boldsymbol{x} \mapsto A \boldsymbol{x}+\boldsymbol{b}(A \in \mathrm{SO}(3)$, $\boldsymbol{b} \in \mathbb{R}^{3}$ ) of $\mathbb{R}^{3}$.

Proof. We have already shown the uniqueness in Proposition 1.16. We shall prove the existence: Let $\Omega(s)$ be as in (1.20), and $\mathcal{F}(s)$ the solution of (1.20) with $\mathcal{F}\left(s_{0}\right)=$ id. Since $\Omega$ is skewsymmetric, $\mathcal{F}(s) \in \mathrm{SO}(3)$ by Proposition 1.10. Denoting the column vectors of $\mathcal{F}$ by $\boldsymbol{e}, \boldsymbol{n}, \boldsymbol{b}$, and let

$$
\gamma(s):=\int_{s_{0}}^{s} \boldsymbol{e}(\sigma) d \sigma
$$

Then $\mathcal{F}$ is the Frenet frame of $\gamma$, and $\kappa$, and $\tau$ are the curvature and torsion of $\gamma$, respectively.

## Exercises

1-1 Find the maximal solution of the initial value problem

$$
\frac{d x}{d t}=x(1-x), \quad x(0)=a
$$

where $b$ is a real number.
1-2 Find an explicit expression of a space curve $\gamma(s)$ parametrized by the arc-length $s$, whose curvature $\kappa$ and torsion $\tau$ satisfy

$$
\kappa=\tau=\frac{1}{\sqrt{2}\left(1+s^{2}\right)}
$$

## 2 Integrability Conditions

Let $U \subset \mathbb{R}^{m}$ be a domain of $\left(\mathbb{R}^{m} ; u^{1}, \ldots, u^{m}\right)$ and consider an $m$-tuple of $n \times n$-matrix valued $C^{\infty}$-maps

$$
\begin{equation*}
\Omega_{j}: \mathbb{R}^{m} \supset U \longrightarrow \mathrm{M}_{n}(\mathbb{R}) \quad(j=1, \ldots, m) . \tag{2.1}
\end{equation*}
$$

In this section, we consider an initial value problem of a system of linear partial differential equations

$$
\begin{equation*}
\frac{\partial X}{\partial u^{j}}=X \Omega_{j} \quad(j=1, \ldots, m), \quad X\left(\mathrm{P}_{0}\right)=X_{0} \tag{2.2}
\end{equation*}
$$

where $\mathrm{P}_{0}=\left(u_{0}^{1}, \ldots, u_{0}^{m}\right) \in U$ is a fixed point, $X$ is an $n \times n$-matrix valued unknown, and $X_{0} \in$ $\mathrm{M}_{n}(\mathbb{R})$.
Proposition 2.1. If a $C^{\infty}$ _map $X: U \rightarrow \mathrm{M}_{n}(\mathbb{R})$ defined on a domain $U \subset \mathbb{R}^{m}$ satisfies (4.1) with $X_{0} \in \operatorname{GL}(n, \mathbb{R})$, then $X(\mathrm{P}) \in \mathrm{GL}(n, \mathbb{R})$ for all $\mathrm{P} \in U$. In addition, if $\Omega_{j}(j=1, \ldots, m)$ are skew-symmetric and $X_{0} \in \mathrm{SO}(n)$, then $X(\mathrm{P}) \in \mathrm{SO}(n)$ holds for all $\mathrm{P} \in U$.
Proof. Since $U$ is connected, there exists a continuous path $\gamma_{0}:[0,1] \rightarrow U$ such that $\gamma_{0}(0)=\mathrm{P}_{0}$ and $\gamma_{0}(1)=$ P. By Whitney's approximation theorem (cf. Theorem 6.21 in [Lee13]), there exists a smooth path $\gamma:[0,1] \rightarrow U$ joining $\mathrm{P}_{0}$ and P approximating $\gamma_{0}$. Since $\hat{X}:=X \circ \gamma$ satisfies (2.4) with $\hat{X}(0)=X_{0}$, Proposition 1.8 yields that $\operatorname{det} \hat{X}(1) \neq 0$ whenever $\operatorname{det} X_{0} \neq 0$. Moreover, if $\Omega_{j}$ 's are skew-symmetric, so is $\Omega_{\gamma}(t)$ in (2.4). Thus, by Proposition 1.10, we obtain the latter half of the proposition.

Proposition 2.2. If a matrix-valued $C^{\infty}$ function $X: U \rightarrow \mathrm{GL}(n, \mathbb{R})$ satisfies (4.1), it holds that

$$
\begin{equation*}
\frac{\partial \Omega_{j}}{\partial u^{k}}-\frac{\partial \Omega_{k}}{\partial u^{j}}=\Omega_{j} \Omega_{k}-\Omega_{k} \Omega_{j} \tag{2.3}
\end{equation*}
$$

for each $(j, k)$ with $1 \leqq j<k \leqq m$.
Proof. Differentiating (4.1) by $u^{k}$, we have

$$
\frac{\partial^{2} X}{\partial u^{k} \partial u^{j}}=\frac{\partial X}{\partial u^{k}} \Omega_{j}+X \frac{\partial \Omega_{j}}{\partial u^{k}}=X\left(\frac{\partial \Omega_{j}}{\partial u^{k}}+\Omega_{k} \Omega_{j}\right) .
$$

On the other hand, switching the roles of $j$ and $k$, we get

$$
\frac{\partial^{2} X}{\partial u^{j} \partial u^{k}}=X\left(\frac{\partial \Omega_{k}}{\partial u^{j}}+\Omega_{j} \Omega_{k}\right) .
$$

Since $X$ is of class $C^{\infty}$, the left-hand sides of these equalities coincide, and so are the right-hand sides. Since $X \in \operatorname{GL}(n, \mathbb{R})$, the conclusion follows.

The equality (2.3) is called the integrability condition or compatibility condition of (4.1).
The chain rule yields the following:
Lemma 2.3. Let $X: U \rightarrow \mathrm{M}_{n}(\mathbb{R})$ be a $C^{\infty}$-map satisfying (4.1). Then for each smooth path $\gamma: I \rightarrow U$ defined on an interval $I \subset \mathbb{R}, \hat{X}:=X \circ \gamma: I \rightarrow \mathrm{M}_{n}(\mathbb{R})$ satisfies the ordinary differential equation

$$
\begin{equation*}
\frac{d \hat{X}}{d t}(t)=\hat{X}(t) \Omega_{\gamma}(t) \quad\left(\Omega_{\gamma}(t):=\sum_{j=1}^{m} \Omega_{j} \circ \gamma(t) \frac{d u^{j}}{d t}(t)\right) \tag{2.4}
\end{equation*}
$$

on I, where $\gamma(t)=\left(u^{1}(t), \ldots, u^{m}(t)\right)$.
20. June, 2023. Revised: 27. June, 2023)

Lemma 2.4. Let $\Omega_{j}: U \rightarrow \mathrm{M}_{n}(\mathbb{R})(j=1, \ldots, m)$ be $C^{\infty}$-maps defined on a domain $U \subset \mathbb{R}^{m}$ which satisfy (2.3). Then for each smooth map

$$
\sigma: D \ni(t, w) \longmapsto \sigma(t, w)=\left(u^{1}(t, w), \ldots, u^{m}(t, w)\right) \in U
$$

defined on a domain $D \subset \mathbb{R}^{2}$, it holds that

$$
\begin{equation*}
\frac{\partial T}{\partial w}-\frac{\partial W}{\partial t}-T W+W T=0 \tag{2.5}
\end{equation*}
$$

where

$$
\begin{equation*}
T:=\sum_{j=1}^{m} \widetilde{\Omega}_{j} \frac{\partial u^{j}}{\partial t}, \quad W:=\sum_{j=1}^{m} \widetilde{\Omega}_{j} \frac{\partial u^{j}}{\partial w} \quad\left(\widetilde{\Omega}_{j}:=\Omega_{j} \circ \sigma\right) \tag{2.6}
\end{equation*}
$$

Proof. By the chain rule, we have

$$
\begin{aligned}
\frac{\partial T}{\partial w} & =\sum_{j, k=1}^{m} \frac{\partial \Omega_{j}}{\partial u^{k}} \frac{\partial u^{k}}{\partial w} \frac{\partial u^{j}}{\partial t}+\sum_{j=1}^{m} \widetilde{\Omega}_{j} \frac{\partial^{2} u^{j}}{\partial w \partial t} \\
\frac{\partial W}{\partial t} & =\sum_{j, k=1}^{m} \frac{\partial \Omega_{j}}{\partial u^{k}} \frac{\partial u^{k}}{\partial t} \frac{\partial u^{j}}{\partial w}+\sum_{j=1}^{m} \widetilde{\Omega}_{j} \frac{\partial^{2} u^{j}}{\partial t \partial w} \\
& =\sum_{j, k=1}^{m} \frac{\partial \Omega_{k}}{\partial u^{j}} \frac{\partial u^{j}}{\partial t} \frac{\partial u^{k}}{\partial w}+\sum_{j=1}^{m} \widetilde{\Omega}_{j} \frac{\partial^{2} u^{j}}{\partial t \partial w}
\end{aligned}
$$

Hence

$$
\begin{aligned}
\frac{\partial T}{\partial w} & -\frac{\partial W}{\partial t}=\sum_{j, k=1}^{m}\left(\frac{\partial \Omega_{j}}{\partial u^{k}}-\frac{\partial \Omega_{k}}{\partial u^{j}}\right) \frac{\partial u^{k}}{\partial w} \frac{\partial u^{j}}{\partial t} \\
& =\sum_{j, k=1}^{m}\left(\widetilde{\Omega}_{j} \widetilde{\Omega}_{k}-\widetilde{\Omega}_{k} \widetilde{\Omega}_{j}\right) \frac{\partial u^{k}}{\partial w} \frac{\partial u^{j}}{\partial t} \\
& =\left(\sum_{j=1}^{m} \widetilde{\Omega}_{j} \frac{\partial u^{j}}{\partial t}\right)\left(\sum_{k=1}^{m} \widetilde{\Omega}_{k} \frac{\partial u^{k}}{\partial w}\right)-\left(\sum_{k=1}^{m} \widetilde{\Omega}_{k} \frac{\partial u^{k}}{\partial w}\right)\left(\sum_{j=1}^{m} \widetilde{\Omega}_{j} \frac{\partial u^{j}}{\partial t}\right) \\
& =T W-W T .
\end{aligned}
$$

Thus (2.5) holds.

Integrability of linear systems. The main theorem in this section is the following theorem:
Theorem 2.5. Let $\Omega_{j}: U \rightarrow \mathrm{M}_{n}(\mathbb{R})(j=1, \ldots, m)$ be $C^{\infty}$-functions defined on a simply connected domain $U \subset \mathbb{R}^{m}$ satisfying (2.3). Then for each $\mathrm{P}_{0} \in U$ and $X_{0} \in \mathrm{M}_{n}(\mathbb{R})$, there exists the unique $n \times n$-matrix valued function $X: U \rightarrow \mathrm{M}_{n}(\mathbb{R})$ satisfying (4.1). Moreover,

- if $X_{0} \in \mathrm{GL}(n, \mathbb{R}), X(\mathrm{P}) \in \mathrm{GL}(n, \mathbb{R})$ holds on $U$,
- if $X_{0} \in \mathrm{SO}(n)$ and $\Omega_{j}(j=1, \ldots, m)$ are skew-symmetric matrices, $X \in \mathrm{SO}(n)$ holds on $U$.

Proof. The latter half is a direct conclusion of Proposition 2.1. We show the existence of $X$ : Take a smooth path $\gamma:[0,1] \rightarrow U$ joining $\mathrm{P}_{0}$ and P . Then by Theorem 1.15 , there exists a unique $C^{\infty}{ }_{-} \operatorname{map} \hat{X}:[0,1] \rightarrow \mathrm{M}_{n}(\mathbb{R})$ satisfying (2.4) with initial condition $\hat{X}(0)=X_{0}$.

We shall show that the value $\hat{X}(1)$ does not depend on choice of paths joining $\mathrm{P}_{0}$ and P . To show this, choose another smooth path $\tilde{\gamma}$ joining $\mathrm{P}_{0}$ and P . Since $U$ is simply connected, there
exists a homotopy between $\gamma$ and $\tilde{\gamma}$, that is, there exists a continuous map $\sigma_{0}:[0,1] \times[0,1] \ni$ $(t, w) \mapsto \sigma(t, w) \in U$ satisfying

$$
\begin{align*}
\sigma_{0}(t, 0) & =\gamma(t), & \sigma_{0}(t, 1) & =\tilde{\gamma}(t),  \tag{2.7}\\
\sigma_{0}(0, w) & =\mathrm{P}_{0}, & & \sigma_{0}(1, w)
\end{align*}
$$

Then, by Whitney's approximation theorem (Theorem 6.21 in [Lee13]) again, there exists a smooth map $\sigma:[0,1] \times[0,1] \rightarrow U$ satisfying the same boundary conditions as (2.7):

$$
\begin{align*}
\sigma(t, 0) & =\gamma(t), & \sigma(t, 1) & =\tilde{\gamma}(t),  \tag{2.8}\\
\sigma(0, w) & =\mathrm{P}_{0}, & & \sigma(1, w)
\end{align*}
$$

We set $T$ and $W$ as in (2.6). For each fixed $w \in[0,1]$, there exists $X_{w}:[0,1] \rightarrow \mathrm{M}_{n}(\mathbb{R})$ such that

$$
\frac{d X_{w}}{d t}(t)=X_{w}(t) T(t, w), \quad X_{w}(0)=X_{0} .
$$

Since $T(t, w)$ is smooth in $t$ and $w$, the map

$$
\check{X}:[0,1] \times[0,1] \ni(t, w) \mapsto X_{w}(t) \in \mathrm{M}_{n}(\mathbb{R})
$$

is a smooth map, because of smoothness in parameter $\alpha$ in Theorem 1.15. To show that $\hat{X}(1)=$ $\check{X}(1,0)$ does not depend on choice of paths, it is sufficient to show that

$$
\begin{equation*}
\frac{\partial \check{X}}{\partial w}=\check{X} W \tag{2.9}
\end{equation*}
$$

holds on $[0,1] \times[0,1]$. In fact, by (2.8), $W(1, w)=0$ for all $w \in[0,1]$, and then (2.9) implies that $\check{X}(1, w)$ is constant.

We prove (2.9): By definition, it holds that

$$
\begin{equation*}
\frac{\partial \check{X}}{\partial t}=\check{X} T, \quad \check{X}(0, w)=X_{0} \tag{2.10}
\end{equation*}
$$

for each $w \in[0,1]$. Hence by (2.5),

$$
\begin{aligned}
\frac{\partial}{\partial t} \frac{\partial \check{X}}{\partial w} & =\frac{\partial^{2} \check{X}}{\partial t \partial w}=\frac{\partial^{2} \check{X}}{\partial w \partial t}=\frac{\partial}{\partial w}(\check{X} T) \\
& =\frac{\partial \check{X}}{\partial w} T+\check{X} \frac{\partial T}{\partial w}=\frac{\partial \check{X}}{\partial w} T+\check{X}\left(\frac{\partial W}{\partial t}+T W-W T\right) \\
& =\frac{\partial \check{X}}{\partial w} T+\check{X} \frac{\partial W}{\partial t}+\frac{\partial \check{X}}{\partial t} W-\check{X} W T \\
& =\frac{\partial}{\partial t}(\check{X} W)+\left(\frac{\partial \check{X}}{\partial w}-\check{X} W\right) T .
\end{aligned}
$$

So, the function $Y_{w}(t):=\partial \check{X} / \partial w-\check{X} W$ satisfies the ordinary differential equation

$$
\frac{d Y_{w}}{d t}(t)=Y_{w}(t) T(t, w), \quad Y_{w}(0)=O
$$

for each $w \in[0,1]$. Thus, by the uniqueness of the solution, $Y_{w}(t)=O$ holds on $[0,1] \times[0,1]$. Hence we have (2.9).

Thus, $\hat{X}(1)$ depends only on the end point P of the path. Hence we can set $X(\mathrm{P}):=\hat{X}(1)$ for each $\mathrm{P} \in U$, and obtain a map $X: U \rightarrow \mathrm{M}_{n}(\mathbb{R})$. Finally we show that $X$ is the desired solution. The initial condition $X\left(\mathrm{P}_{0}\right)=X_{0}$ is obviously satisfied. On the other hand, if we set

$$
Z(\delta):=X\left(u^{1}, \ldots, u^{j}+\delta, \ldots, u^{m}\right),
$$

$Z(\delta)$ satisfies the equation (2.4) for the path $\gamma(\delta):=\left(u^{1}, \ldots, u^{j}+\delta, \ldots, u^{m}\right)$ with $Z(0)=X(\mathrm{P})$. Since $\Omega_{\gamma}=\Omega_{j}$,

$$
\frac{\partial X}{\partial u^{j}}(\mathrm{P})=\left.\frac{d Z}{d \delta}\right|_{\delta=0}=Z(0) \Omega_{j}(\mathrm{P})=X(\mathrm{P}) \Omega_{j}(\mathrm{P})
$$

which completes the proof.

## Application: Poincaré's lemma.

Theorem 2.6 (Poincaré's lemma). If a differential 1-form

$$
\omega=\sum_{j=1}^{m} \alpha_{j}\left(u^{1}, \ldots, u^{m}\right) d u^{j}
$$

defined on a simply connected domain $U \subset \mathbb{R}^{m}$ is closed, that is, $d \omega=0$ holds, then there exists a $C^{\infty}$-function $f$ on $U$ such that $d f=\omega$. Such a function $f$ is unique up to additive constants.

Proof. Since

$$
d \omega=\sum_{i<j}\left(\frac{\partial \alpha_{j}}{\partial u^{i}}-\frac{\partial \alpha_{i}}{\partial u^{j}}\right) d u^{i} \wedge d u^{j},
$$

the assumption is equivalent to

$$
\begin{equation*}
\frac{\partial \alpha_{j}}{\partial u^{i}}-\frac{\partial \alpha_{i}}{\partial u^{j}}=0 \quad(1 \leqq i<j \leqq m) . \tag{2.11}
\end{equation*}
$$

Consider a system of linear partial differential equations with unknown $\xi$, a $1 \times 1$-matrix valued function (i.e. a real-valued function), as

$$
\begin{equation*}
\frac{\partial \xi}{\partial u^{j}}=\xi \alpha_{j} \quad(j=1, \ldots, m), \quad \xi\left(u_{0}^{1}, \ldots, u_{0}^{m}\right)=1 . \tag{2.12}
\end{equation*}
$$

Then it satisfies (2.3) because of (2.11). Hence by Theorem 4.5, there exists a smooth function $\xi\left(u^{1}, \ldots, u^{m}\right)$ satisfying (2.12). In particular, Proposition 1.8 yields $\xi=\operatorname{det} \xi$ never vanishes. Hence $\xi\left(u_{0}^{1}, \ldots, u_{0}^{m}\right)=1>0$ means that $\xi>0$ holds on $U$. Letting $f:=\log \xi$, we have the function $f$ satisfying $d f=\omega$.

Next, we show the uniqueness: if two functions $f$ and $g$ satisfy $d f=d g=\omega$, it holds that $d(f-g)=0$. Hence by connectivity of $U, f-g$ must be constant.

Application: Conjugation of Harmonic functions. In this paragraph, we identify $\mathbb{R}^{2}$ with the complex plane $\mathbb{C}$. It is well-known that a smooth function

$$
\begin{equation*}
f: U \ni u+\mathrm{i} v \longmapsto \xi(u, v)+\mathrm{i} \eta(u, v) \in \mathbb{C} \quad(\mathrm{i}=\sqrt{-1}) \tag{2.13}
\end{equation*}
$$

defined on a domain $U \subset \mathbb{C}$ is holomorphic if and only if it satisfies the following relation, called the Cauchy-Riemann equations:

$$
\begin{equation*}
\frac{\partial \xi}{\partial u}=\frac{\partial \eta}{\partial v}, \quad \frac{\partial \xi}{\partial v}=-\frac{\partial \eta}{\partial u} . \tag{2.14}
\end{equation*}
$$

Definition 2.7. A function $f: U \rightarrow \mathbb{R}$ defined on a domain $U \subset \mathbb{R}^{2}$ is said to be harmonic if it satisfies

$$
\Delta f=f_{u u}+f_{v v}=0 .
$$

The operator $\Delta$ is called the Laplacian.

Proposition 2.8. If function $f$ in (2.13) is holomorphic, $\xi(u, v)$ and $\eta(u, v)$ are harmonic functions.

Proof. By (2.14), we have

$$
\xi_{u u}=\left(\xi_{u}\right)_{u}=\left(\eta_{v}\right)_{u}=\eta_{v u}=\eta_{u v}=\left(\eta_{u}\right)_{v}=\left(-\xi_{v}\right)_{v}=-\xi_{v v}
$$

Hence $\Delta \xi=0$. Similarly,

$$
\eta_{u u}=\left(-\xi_{v}\right)_{u}=-\xi_{v u}=-\xi_{u v}=-\left(\xi_{u}\right)_{v}=-\left(\eta_{v}\right)_{v}=-\eta_{v v}
$$

Thus $\Delta \eta=0$.
Theorem 2.9. Let $U \subset \mathbb{C}=\mathbb{R}^{2}$ be a simply connected domain and $\xi(u, v)$ a $C^{\infty}$-function harmonic on $U^{6}$. Then there exists a $C^{\infty}$ harmonic function $\eta$ on $U$ such that $\xi(u, v)+\mathrm{i} \eta(u, v)$ is holomorphic on $U$.
Proof. Let $\alpha:=-\xi_{v} d u+\xi_{u} d v$. Then by the assumption,

$$
d \alpha=\left(\xi_{v v}+\xi_{u u}\right) d u \wedge d v=0
$$

holds, that is, $\alpha$ is a closed 1-form. Hence by simple connectivity of $U$ and the Poincaré's lemma (Theorem 4.8), there exists a function $\eta$ such that $d \eta=\eta_{u} d u+\eta_{v} d v=\alpha$. Such a function $\eta$ satisfies (2.14) for given $\xi$. Hence $\xi+\mathrm{i} \eta$ is holomorphic in $u+\mathrm{i} v$.

Example 2.10. A function $\xi(u, v)=e^{u} \cos v$ is harmonic. Set

$$
\alpha:=-\xi_{v} d u+\xi_{u} d v=e^{u} \sin v d u+e^{u} \cos v d v
$$

Then $\eta(u, v)=e^{u} \sin v$ satisfies $d \eta=\alpha$. Hence

$$
\xi+\mathrm{i} \eta=e^{u}(\cos v+\mathrm{i} \sin v)=e^{u+\mathrm{i} v}
$$

is holomorphic in $u+\mathrm{i} v$.
Definition 2.11. The harmonic function $\eta$ in Theorem 2.9 is called the conjugate harmonic function of $\xi$.

## Exercises

2-1 Let $\xi(u, v):=\log \sqrt{u^{2}+v^{2}}$ be a function defined on $U:=\mathbb{R}^{2} \backslash\{(0,0)\}$.
(1) Show that $\xi$ is harmonic on $U$.
(2) Find the conjugate harmonic function $\eta$ of $\xi$ on

$$
V=\mathbb{R}^{2} \backslash\{(u, 0) \mid u \leqq 0\} \subset U
$$

(3) Show that there exists no conjugate harmonic function of $\xi$ defined on $U$.

2-2 Consider a linear system of partial differential equations for $3 \times 3$-matrix valued unknown $X$ on a domain $U \subset \mathbb{R}^{2}$ as

$$
\frac{\partial X}{\partial u}=X \Omega, \quad \frac{\partial X}{\partial v}=X \Lambda, \quad\left(\Omega:=\left(\begin{array}{ccc}
0 & -\alpha & -h_{1}^{1} \\
\alpha & 0 & -h_{1}^{2} \\
h_{1}^{1} & h_{1}^{2} & 0
\end{array}\right), \quad \Lambda:=\left(\begin{array}{ccc}
0 & -\beta & -h_{2}^{1} \\
\beta & 0 & -h_{2}^{2} \\
h_{2}^{1} & h_{2}^{2} & 0
\end{array}\right)\right)
$$

where $(u, v)$ are the canonical coordinate system of $\mathbb{R}^{2}$, and $\alpha, \beta$ and $h_{j}^{i}(i, j=1,2)$ are smooth functions defined on $U$. Write down the integrability conditions in terms of $\alpha, \beta$ and $h_{j}^{i}$.

[^3]
## 3 Differential Forms

Let $M$ be an $n$-dimensional manifold and denote by $\mathcal{F}(M)$ and $\mathfrak{X}(M)$ the set of smooth function and the set of smooth vector fields on $M$, respectively.

Lie brackets A vector field $X \in \mathfrak{X}(M)$ can be considered as a differential operator acting on $\mathcal{F}(M)$ as $(X f)(p)=X_{p} f$. By definition it satisfies the Leibniz rule

$$
\begin{equation*}
X(f g)=f(X g)+g(X f) \quad(X \in \mathfrak{X}(M), f, g \in \mathcal{F}(M)) \tag{3.1}
\end{equation*}
$$

For two vector fields $X, Y \in \mathfrak{X}(M)$, set

$$
\begin{equation*}
[X, Y]: \mathcal{F}(M) \ni f \longmapsto X(Y f)-Y(X f) \in \mathcal{F}(M) \tag{3.2}
\end{equation*}
$$

Then $[X, Y]$ also satisfies the Leibnitz rule (3.1), and gives a vector field on $M$. The map

$$
[,]: \mathfrak{X}(M) \times \mathfrak{X}(M) \ni(X, Y) \mapsto[X, Y] \in \mathfrak{X}(M)
$$

is called the Lie bracket on $\mathfrak{X}(M)$. One can easily show that the product [, ] is bilinear, skew symmetric and satisfies the Jacobi identity

$$
\begin{equation*}
[X,[Y, Z]]+[Y,[Z, X]]+[Z,[X, Y]]=\mathbf{0} \tag{3.3}
\end{equation*}
$$

that is, $(\mathfrak{X}(M),[]$,$) is a Lie algebra (of infinite dimension). By the Leibniz rule, it holds that$

$$
\begin{equation*}
[f X, Y]=f[X, Y]-(Y f) X, \quad[X, f Y]=f[X, Y]+(X f) Y \quad(X, Y \in \mathfrak{X}(M), f \in \mathcal{F}(M)) \tag{3.4}
\end{equation*}
$$

Tensors. For each $p \in M$, the dual space $T_{p}^{*} M$ of $T_{p} M$ is the liner space consisting of all linear maps from $T_{p} M$ to $\mathbb{R}$.

Lemma 3.1. Let $\left(x^{1}, \ldots, x^{n}\right)$ be a local coordinate system of $M$ around $p$, and set

$$
\left(\frac{\partial}{\partial x^{j}}\right)_{p}: \mathcal{F}(M) \ni f \mapsto \frac{\partial f}{\partial x^{j}}(p), \quad\left(d x^{j}\right)_{p}: T_{p} M \rightarrow \mathbb{R} \quad \text { with } \quad\left(d x^{j}\right)_{p}\left(\left(\frac{\partial}{\partial x^{k}}\right)_{p}\right)=\delta_{k}^{j}
$$

for $j, k=1, \ldots, n$. Then $\left\{\left(\partial / \partial x^{j}\right)_{p}\right\}_{j=1, \ldots, n}$ and $\left\{\left(d x^{j}\right)_{p}\right\}_{j=1, \ldots, n}$ are a basis of $T_{p} M$ and $T_{p}^{*} M$, respectively, where $\delta_{k}^{j}$ denotes Kronecker's delta symbol.

We let

$$
T_{p}^{*} M \otimes T_{p}^{*} M \quad\left(\text { resp. } \quad T_{p}^{*} M \otimes T_{p}^{*} M \otimes T_{p}^{*} M\right)
$$

the set of bilinear (resp. trilinear) maps of $T_{p} M \times T_{p} M$ (resp. $T_{p} M \times T_{p} M \times T_{p} M$ ) to $\mathbb{R}$. A section of the vector bundle
$T^{*} M \otimes T^{*} M:=\bigcup_{p \in M} T_{p}^{*} M \otimes T_{p}^{*} M \quad\left(\operatorname{resp} . T^{*} M \otimes T^{*} M \otimes T^{*} M:=\bigcup_{p \in M} T_{p}^{*} M \otimes T_{p}^{*} M \otimes T_{p}^{*} M\right)$
is called a covariant 2 (resp. 3)-tensor.
A section $\omega \in \Gamma\left(T^{*} M\right)$ of the cotangent bundle $T^{*} M$ is called a covariant 1-tensor or a 1-form. A one form $\omega$ induces a linear map

$$
\begin{equation*}
\omega: \mathfrak{X}(M) \ni X \longmapsto \omega(X) \in \mathcal{F}(M), \quad \text { where } \quad \omega(X)(p)=\omega_{p}\left(X_{p}\right) \tag{3.5}
\end{equation*}
$$

By definition, it holds that

$$
\frac{(3.6)}{\text { 27. June, 2023. Revised: 04. July, 2023) }} \quad \omega(f X)=f \omega(X) \quad(f \in \mathcal{F}(M), X \in \mathfrak{X}(M))
$$

Lemma 3.2. A linear map $\omega: \mathfrak{X}(M) \rightarrow \mathcal{F}(M)$ is a 1 -form if and only if (3.6) holds.
Proof. The "only if" part is trivial by definition. Assume a linear map $\omega: \mathfrak{X}(M) \rightarrow \mathcal{F}(M)$ satisfies (3.6). In fact, under a local coordinate system $\left(x^{1}, \ldots, x^{n}\right)$ around $p \in M$,

$$
\omega(X)(p)=\omega\left(\sum_{j=1}^{n} X^{j} \frac{\partial}{\partial x^{j}}\right)(p)=\sum_{j=1}^{n} X^{j}(p) \omega\left(\frac{\partial}{\partial x^{j}}\right)_{p}, \quad\left(X=\sum_{j=1}^{n} X^{j} \frac{\partial}{\partial x^{j}} .\right)
$$

holds. In other words, $\omega(X)(p)$ depend only on $X_{p}$. Hence $\omega$ induces a map $\omega_{p}: T_{p} M \rightarrow \mathbb{R}$.
Similarly, a covariant 2 (resp. 3) tensor $\alpha \in \Gamma\left(T^{*} M \otimes T^{*} M\right)$ (resp. $\beta \in \Gamma\left(T^{*} M \otimes T^{*} M \otimes\right.$ $\left.T^{*} M\right)$ )induces a bilinear (resp. trilinear) map $\alpha: \mathfrak{X}(M) \times \mathfrak{X}(M) \rightarrow \mathcal{F}(M)$. (resp. $\beta: \mathfrak{X}(M) \times$ $\mathfrak{X}(M) \times \mathfrak{X}(M) \rightarrow \mathcal{F}(M)$. By the same reason as Lemma 3.2, we have

Lemma 3.3. $A$ bilinear map $\alpha: \mathfrak{X}(M) \times \mathfrak{X}(M) \rightarrow \mathcal{F}(M)($ resp. $\beta: \mathfrak{X}(M) \times \mathfrak{X}(M) \times \mathfrak{X}(M) \rightarrow \mathcal{F}(M))$ is a a covariant 2 (resp. 3)-tensor if and only if

$$
\begin{aligned}
& \alpha(f X, Y)=\alpha(X, f Y)=f \alpha(X, Y) \\
& \quad(\text { resp. } \quad \beta(f X, Y, Z)=\beta(X, f Y, Z)=\beta(X, Y, f Z)=f \beta(X, Y, Z))
\end{aligned}
$$

holds for all $X, Y, Z \in \mathfrak{X}(M)$ and $f \in \mathcal{F}(M)$.
A covariant 2 (resp. 3 )-tensor $\alpha$ (resp. $\beta$ ) said to be skew-symmetric if

$$
\alpha(X, Y)=-\alpha(Y, X), \quad(\beta(X, Y, Z)=-\beta(Y, X, Z)=-\beta(X, Z, Y)=-\beta(Z, Y, X))
$$

holds for all $X, Y, Z \in \mathfrak{X}(M)$. We denote

$$
\wedge^{k}(M):= \begin{cases}\mathcal{F}(M) & (k=0)  \tag{3.7}\\ \Gamma\left(T^{*} M\right) & (k=1) \\ \left\{\omega \in \Gamma\left(T^{*} M \otimes T^{*} M\right) ; \omega \text { is skew-symmetric }\right\} & (k=2) \\ \left\{\omega \in \Gamma\left(T^{*} M \otimes T^{*} M \otimes T^{*} M\right) ; \omega \text { is skew-symmetric }\right\} & (k=3)\end{cases}
$$

An element of $\wedge^{k}(M)$ is called an $k$-form.

The Exterior products. The exterior product $\alpha \wedge \beta \in \wedge^{2}(M)$ of two 1-forms $\alpha, \beta \in \wedge^{1}(M)$ is defined as

$$
\begin{equation*}
(\alpha \wedge \beta)(X, Y):=\alpha(X) \beta(Y)-\alpha(Y) \beta(X) \tag{3.8}
\end{equation*}
$$

On the other hand, the exterior product of $\alpha$ and $\omega$ is defined as a 3 -form on $M$ by

$$
\begin{equation*}
(\alpha \wedge \omega)(X, Y, Z)=(\omega \wedge \alpha)(X, Y, Z):=\alpha(X, Y) \omega(Z)+\alpha(Y, Z) \omega(X)+\alpha(Z, X) \omega(Y) \tag{3.9}
\end{equation*}
$$

Then by a direct computation together with (3.8), it holds that

$$
\begin{equation*}
(\mu \wedge \omega) \wedge \lambda=\mu \wedge(\omega \wedge \lambda)(=: \mu \wedge \omega \wedge \lambda) \tag{3.10}
\end{equation*}
$$

for 1-forms $\mu, \omega$ and $\lambda$.

The Exterior derivative. Under a local coordinate system $\left(x^{1}, \ldots, x^{n}\right)$, a one form $\alpha$ and a two form $\omega$ are expressed as

$$
\alpha=\sum_{j=1}^{n} \alpha_{j} d x^{j}, \quad \omega=\sum_{1 \leqq i<j \leqq n} \omega_{i j} d x^{i} \wedge d x^{j}
$$

where $\alpha_{j}(j=1, \ldots, n)$ and $\omega_{i j}(1 \leqq i<j \leqq n)$ are smooth functions in $\left(x^{1}, \ldots, x^{n}\right)$. By Lemma 3.3 and the property (3.4) of the Lie brackets, we have
Lemma 3.4. For a function $f \in \mathcal{F}(M)=\wedge^{0}(M)$, a 1 -form $\alpha \in \wedge^{1}(M)$ and a 2 -form $\beta \in \wedge^{2}(M)$ )

$$
\begin{aligned}
d f & : \mathfrak{X}(M) \ni X \mapsto d f(X)=X f \in \mathcal{F}(M) \\
d \alpha: & \mathfrak{X}(M) \times \mathfrak{X}(M) \ni(X, Y) \mapsto X \alpha(Y)-Y \alpha(X)-\alpha([X, Y]) \in \mathcal{F}(M) \\
d \beta: & \mathfrak{X}(M) \times \mathfrak{X}(M) \times \mathfrak{X}(M) \ni(X, Y, Z) \mapsto \\
& X \beta(Y, Z)+Y \beta(Z, X)+Z \beta(X, Y)-\beta([X, Y], Z)-\beta([Y, Z], Z)-\beta([Z, X], Y)
\end{aligned}
$$

are a 1-form, a 2-form and a 3-form respectively.
Definition 3.5. For a function $f$, a 1 -form $\alpha$ and a 2 -form $\beta, d f, d \alpha$ and $d \beta$ are called the exterior derivatives of $f, \alpha$ and $\beta$, respectively.

Then, for one forms $\mu$ and $\omega$, we have

$$
\begin{equation*}
d d \omega=0, \quad d(\mu \wedge \omega)=d \mu \wedge \omega-\mu \wedge d \omega \tag{3.11}
\end{equation*}
$$

by the definition and the Jacobi identity (3.3).
The Riemannian connection. In the rest of this section, we let $(M, g)$ be an $n$-dimensional (pseudo) Riemannian manifold, and denote by $\langle$,$\rangle the inner product induced by g$.

Lemma 3.6. There exists the unique bilinear map $\nabla: \mathfrak{X}(M) \times \mathfrak{X}(M) \ni(X, Y) \mapsto \nabla_{X} Y \in \mathfrak{X}(M)$ satisfying

$$
\begin{equation*}
\nabla_{X} Y-\nabla_{Y} X=[X, Y], \quad X\langle Y, Z\rangle=\left\langle\nabla_{X} Y, Z\right\rangle+\left\langle X, \nabla_{X} Z\right\rangle \quad(X, Y, Z \in \mathfrak{X}(M)) \tag{3.12}
\end{equation*}
$$

Definition 3.7. The map $\nabla$ in Lemma 3.6 is called the Riemannian connection or the Levi-Civita connection of $(M, g)$.

Lemma 3.8. The Riemannian connection $\nabla$ satisfies

$$
\begin{equation*}
\nabla_{f X} Y=f \nabla_{X} Y, \quad \nabla_{X}(f Y)=(X f) Y+f \nabla_{X} Y \tag{3.13}
\end{equation*}
$$

Remark 3.9. A bilinear map $\nabla: \mathfrak{X}(M) \times \mathfrak{X}(M) \rightarrow \mathfrak{X}(M)$ satisfying (3.13) is called a linear connection or an affine connection.
Remark 3.10. By Lemmas 3.8 and 3.2, $X \mapsto \nabla_{X} Y$ determines a one form.

Orthonormal frames. For a sake of simplicity, we assume that $g$ is positive definite, in other words, $(M, g)$ is a Riemannian manifold.

Definition 3.11. Let $U \subset M$ be a domain of $M$. An $n$-tuple of vector fields $\left\{\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right\}$ on $U$ is called an orthonormal frame on $U$ if $\left\langle\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right\rangle=\delta_{i j}$. It is said to be positive if $M$ is oriented and $\left\{\boldsymbol{e}_{j}\right\}$ is compatible to the orientation on $M$.

Remark 3.12. For each $p \in M$, there exists a neighborhood $U$ of $p$ which admits an orthonormal frame on $U$.

Lemma 3.13. Let $\left\{\boldsymbol{e}_{j}\right\}$ and $\left\{\boldsymbol{v}_{j}\right\}$ be two orthonormal frames on $U \subset M$. Then there exists a smooth map

$$
\begin{equation*}
\Theta: U \longrightarrow \mathrm{O}(n) \quad \text { such that } \quad\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]=\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right] \Theta \tag{3.14}
\end{equation*}
$$

Moreover, if $\left\{\boldsymbol{e}_{j}\right\}$ and $\left\{\boldsymbol{v}_{j}\right\}$ determines the common orientation, $\Theta$ is valued on $\mathrm{SO}(n)$.
The map $\Theta$ in Lemma 3.13 is called a gauge transformation.
For an orthonormal frame $\left\{\boldsymbol{e}_{j}\right\}$ on $U$, we denote by $\left\{\omega^{j}\right\}_{j=1, \ldots, n}$ the dual frame of $\left\{\boldsymbol{e}_{j}\right\}$, that is, $\omega^{j} \in \wedge^{1}(U)$ such that

$$
\omega^{j}\left(\boldsymbol{e}_{k}\right)=\delta_{k}^{j}= \begin{cases}1 & (j=k) \\ 0 & \text { (otherwise) }\end{cases}
$$

In other words, $\omega^{j}(X)=\left\langle\boldsymbol{e}_{j}, X\right\rangle$.
Lemma 3.14. Two orthonormal frames $\left\{\boldsymbol{e}_{j}\right\}$ and $\left\{\boldsymbol{v}_{j}\right\}$ are related as (3.14). Then their duals $\left\{\omega^{j}\right\}$ and $\left\{\lambda^{j}\right\}$ satisfy

$$
\left(\begin{array}{c}
\lambda^{1} \\
\vdots \\
\lambda^{n}
\end{array}\right)=\Theta\left(\begin{array}{c}
\omega^{1} \\
\vdots \\
\omega^{n}
\end{array}\right)
$$

Proof.

$$
\left(\begin{array}{c}
\lambda^{1} \\
\vdots \\
\lambda^{n}
\end{array}\right)\left(\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right)=\left(\begin{array}{c}
\lambda^{1} \\
\vdots \\
\lambda^{n}
\end{array}\right)\left(\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right) \Theta=\Theta=\Theta\left(\begin{array}{c}
\omega^{1} \\
\vdots \\
\omega^{n}
\end{array}\right)\left(\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right)
$$

## Connection forms.

Definition 3.15. The connection form with respect to an orthonormal frame $\left\{\boldsymbol{e}_{j}\right\}$ is a $n \times n$-matrix valued one form $\Omega$ on $U$ defined by

$$
\Omega=\left(\begin{array}{cccc}
\omega_{1}^{1} & \omega_{2}^{1} & \ldots & \omega_{n}^{1} \\
\omega_{1}^{2} & \omega_{2}^{2} & \ldots & \omega_{n}^{2} \\
\vdots & \vdots & \ddots & \vdots \\
\omega_{1}^{n} & \omega_{2}^{n} & \ldots & \omega_{n}^{n}
\end{array}\right), \quad \omega_{j}^{k}:=\left\langle\nabla \boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right\rangle \in \wedge^{1}(U)
$$

By definition, we have $\nabla \boldsymbol{e}_{j}=\sum_{k=1}^{n} \omega_{j}^{k} \boldsymbol{e}_{k}$, that is, $\nabla\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]=\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right] \Omega$.
Lemma 3.16. $\omega_{j}^{k}=-\omega_{k}^{j}$.
Proof. $\omega_{j}^{k}=\left\langle\nabla \boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right\rangle=d\left\langle\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right\rangle-\left\langle\boldsymbol{e}_{j}, \nabla \boldsymbol{e}_{k}\right\rangle=-\omega_{k}^{j}$.

Lemma 3.17. $d \omega^{i}=\sum_{l=1}^{n} \omega^{l} \wedge \omega_{l}^{i}$.
Proof.

$$
\begin{aligned}
d \omega^{i}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right) & =\boldsymbol{e}_{j} \omega^{i}\left(\boldsymbol{e}_{k}\right)-\boldsymbol{e}_{k} \omega^{i}\left(\boldsymbol{e}_{j}\right)-\omega^{i}\left(\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right]\right)=-\omega^{i}\left(\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right]\right) \\
& =-\omega^{i}\left(\nabla \boldsymbol{e}_{j} \boldsymbol{e}_{k}-\nabla \boldsymbol{e}_{k} \boldsymbol{e}_{j}\right)=-\left\langle\nabla \boldsymbol{e}_{j} \boldsymbol{e}_{k}-\nabla \boldsymbol{e}_{k} \boldsymbol{e}_{j}, \boldsymbol{e}_{i}\right\rangle=-\omega_{k}^{i}\left(\boldsymbol{e}_{j}\right)+\omega_{j}^{i}\left(\boldsymbol{e}_{k}\right) \\
& =\sum_{l=1}^{n}\left(-\omega_{l}^{i}\left(\boldsymbol{e}_{j}\right) \omega^{l}\left(\boldsymbol{e}_{k}\right)+\omega_{l}^{i}\left(\boldsymbol{e}_{k}\right) \omega^{l}\left(\boldsymbol{e}_{j}\right)\right)=\sum_{l=1}^{n} \omega^{l} \wedge \omega_{l}^{i}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right) .
\end{aligned}
$$

## Exercises

3-1 Let $\left\{\boldsymbol{e}_{j}\right\}$ and $\left\{\boldsymbol{v}_{j}\right\}$ be two orthonormal frames on a domain $U$ of a Riemannian $n$-manifold $M$, which are related as (3.14). Show that the connection forms $\Omega$ of $\left\{\boldsymbol{e}_{j}\right\}$ and $\Lambda$ of $\left\{\boldsymbol{v}_{j}\right\}$ satisfy $\Omega=\Theta^{-1} \Lambda \Theta+\Theta^{-1} d \Theta$.

3-2 Let $\mathbb{R}_{1}^{3}$ be the 3-dimensional Lorentz-Minkowski space and let $H^{2}(-1)$ the hyperbolic 2-space (i.e. the hyperbolic plane) of constant curvature -1 .
(1) Verify that

$$
\boldsymbol{f}(u, v):=(\cosh u, \cos v \sinh u, \sin v \sinh u)
$$

gives a local coordinate system on $U:=H^{2}(-1) \backslash\{(1,0,0)\}$, and

$$
\boldsymbol{e}_{1}:=(\sinh u, \cos v \cosh u, \sin v \cosh u), \quad \boldsymbol{e}_{2}:=(0,-\sin v, \cos v)
$$

forms a orthonormal frame on $U$.
(2) Compute the connection form(s) with respect to the orthonormal frame $\left\{\boldsymbol{e}_{1}, \boldsymbol{e}_{2}\right\}$.

## 4 Curvatre forms

### 4.1 Addendum to the previous section

Proposition 4.1 (The local expression of the Lie bracket). Let $\left(U ; x^{1}, \ldots, x^{n}\right)$ be a coordinate neighborhood of an n-manifold $M$. Then the Lie bracket of two vector fields

$$
X=\sum_{j=1}^{n} \xi^{j} \frac{\partial}{\partial x^{j}}, \quad Y=\sum_{j=1}^{n} \eta^{j} \frac{\partial}{\partial x^{j}}
$$

is expressed as

$$
[X, Y]=\sum_{j=1}^{n}\left(\xi^{k} \frac{\partial \eta^{j}}{\partial x^{k}}-\eta^{k} \frac{\partial \xi^{j}}{\partial x^{k}}\right) \frac{\partial}{\partial x^{j}}
$$

Proof. For a smooth function $f$ on $U$, it holds that

$$
\frac{\partial}{\partial x^{i}} \frac{\partial}{\partial x^{j}} f=\frac{\partial^{2} f}{\partial x^{i} \partial x^{j}}=\frac{\partial^{2} f}{\partial x^{j} \partial x^{i}}=\frac{\partial}{\partial x^{j}} \frac{\partial}{\partial x^{i}} f
$$

Hence $\left[\partial / \partial x^{i}, \partial / \partial x^{j}\right]=0$. Then the conclusion follows from bilinearlity of $[X, Y]$ and the formula

$$
[f X, Y]=f[X, Y]-(Y f) X, \quad[X, f Y]=f[X, Y]+(X f) Y
$$

for a smooth function $f$ and vector fields $X$ and $Y$.
Proposition 4.2 (A local expression of the connection forms). Let $U$ be a domain of a Riemannian n-manifold $(M, g)$ and $\left[e_{1}, \ldots, e_{n}\right]$ an orthonormal frame on $U$. Then the connection form $\omega_{i}^{j}$ with respect to the frame $\left[\boldsymbol{e}_{j}\right]$ is obtained as

$$
\omega_{i}^{j}\left(\boldsymbol{e}_{k}\right)=\frac{1}{2}\left(-\left\langle\left[\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right], \boldsymbol{e}_{k}\right\rangle+\left\langle\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right], \boldsymbol{e}_{i}\right\rangle+\left\langle\left[\boldsymbol{e}_{k}, \boldsymbol{e}_{i}\right], \boldsymbol{e}_{j}\right\rangle\right)
$$

where $\langle$,$\rangle denotes the inner product induced from g$.
Proof. By the definition of the Levi-Civita connection $\nabla$,

$$
\begin{aligned}
\omega_{i}^{j}\left(\boldsymbol{e}_{k}\right) & =\left\langle\nabla \boldsymbol{e}_{k} \boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right\rangle=\boldsymbol{e}_{k}\left\langle\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right\rangle-\left\langle\boldsymbol{e}_{i}, \nabla \boldsymbol{e}_{k} \boldsymbol{e}_{j}\right\rangle=-\left\langle\boldsymbol{e}_{i}, \nabla \boldsymbol{e}_{j} \boldsymbol{e}_{k}+\left[\boldsymbol{e}_{k}, \boldsymbol{e}_{j}\right]\right\rangle \\
& =-\boldsymbol{e}_{j}\left\langle\boldsymbol{e}_{i}, \boldsymbol{e}_{k}\right\rangle+\left\langle\nabla \boldsymbol{e}_{j} \boldsymbol{e}_{i}, \boldsymbol{e}_{k}\right\rangle-\left\langle\boldsymbol{e}_{i},\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right]\right\rangle \\
& =\left\langle\nabla \boldsymbol{e}_{i} \boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right\rangle+\left\langle\left[\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right], \boldsymbol{e}_{k}\right\rangle-\left\langle\boldsymbol{e}_{i},\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right]\right\rangle \\
& =\boldsymbol{e}_{i}\left\langle\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right\rangle-\left\langle\boldsymbol{e}_{j}, \nabla \boldsymbol{e}_{i} \boldsymbol{e}_{k}\right\rangle+\left\langle\left[\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right], \boldsymbol{e}_{k}\right\rangle-\left\langle\boldsymbol{e}_{i},\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right]\right\rangle \\
& =-\left\langle\boldsymbol{e}_{j}, \nabla \boldsymbol{e}_{k} \boldsymbol{e}_{i}\right\rangle-\left\langle\boldsymbol{e}_{j},\left[\boldsymbol{e}_{i}, \boldsymbol{e}_{k}\right]\right\rangle+\left\langle\left[\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right], \boldsymbol{e}_{k}\right\rangle-\left\langle\boldsymbol{e}_{i},\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right]\right\rangle \\
& =-\omega_{i}^{j}\left(\boldsymbol{e}_{k}\right)+\left\langle\left[\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right], \boldsymbol{e}_{k}\right\rangle-\left\langle\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right], \boldsymbol{e}_{i}\right\rangle+\left\langle\left[\boldsymbol{e}_{k}, \boldsymbol{e}_{i}\right], \boldsymbol{e}_{j}\right\rangle
\end{aligned}
$$

### 4.2 Preliminaries

Integrability condition, a review. Let $U$ be a domain of $\mathbb{R}^{m}$ with coordinate system $\left(x^{1}, \ldots, x^{m}\right)$, and consider a system of differential equations

$$
\begin{equation*}
\frac{\partial F}{\partial x^{l}}=F \Omega_{l} \quad(l=1, \ldots, m) \tag{4.1}
\end{equation*}
$$

with initial condition

$$
\begin{equation*}
F\left(\mathrm{P}_{0}\right)=F_{0} \in \mathrm{M}_{n}(\mathbb{R}), \quad \mathrm{P}_{0}=\left(x_{0}^{1}, \ldots, x_{0}^{m}\right) \in U \tag{4.2}
\end{equation*}
$$

where $F$ is an unknown map into the space of $n \times n$-real matrices $\mathrm{M}_{n}(\mathbb{R})$, and the coefficient matrices $\Omega_{l}(l=1, \ldots, m)$ are $\mathrm{M}_{n}(\mathbb{R})$-valued $C^{\infty}$-functions.
Lemma 4.3. If the initial condition $F_{0}$ in (4.2) is non-singular, i.e., $F_{0} \in \mathrm{GL}(n, \mathbb{R})^{7}, F$ satisfying

[^4](4.1) is a $\mathrm{GL}(n, \mathbb{R})$-valued function, that is, $F$ is invertible for each point on $U$.

Proof. For each $\mathrm{P} \in U$, take a smooth path $\gamma(t):=\left(x^{1}(t), \ldots, x^{m}(t)\right)(0 \leqq t \leqq 1)$ with $\gamma(0)=\mathrm{P}_{0}$ and $\gamma(1)=\mathrm{P}$. Then the matrix-valued function $\hat{F}:=F \circ \gamma$ of one variable satisfies the ordinary differential equation

$$
\frac{d \hat{F}}{d t}=\hat{F} \hat{\Omega}, \quad \hat{\Omega}:=\sum_{l=1}^{m} \Omega_{l} \circ \gamma \frac{d x^{l}}{d t} .
$$

Hence $\varphi:=\operatorname{det} \hat{F}$ satisfies

$$
\frac{d \varphi}{d t}=\frac{d}{d t} \operatorname{det} \hat{F}=\operatorname{tr}\left(\tilde{\hat{F}} \frac{d \hat{F}}{d t}\right)=\operatorname{tr}(\tilde{\hat{F}} \hat{F} \hat{\Omega})=\operatorname{det} \hat{F} \operatorname{tr} \hat{\Omega}=\varphi \omega
$$

where $\widetilde{\hat{F}}$ denotes the cofactor matrix of $\hat{F}$ and $\omega:=\operatorname{tr} \hat{\Omega}$. So

$$
\operatorname{det} \hat{F}(t)=\varphi(t)=\varphi_{0} \exp \int_{0}^{t} \omega(\tau) d \tau \quad\left(\varphi_{0}:=\operatorname{det} F_{0}\right)
$$

proving the lemma.
As seen in the previous lectures the following integrability condition holds:
Lemma 4.4. If a $C^{\infty}$-map $F: U \rightarrow \mathrm{GL}(n, \mathbb{R})$ satisfies (4.1), then it hold on $U$ that

$$
\begin{equation*}
\frac{\partial \Omega_{l}}{\partial x^{k}}-\frac{\partial \Omega_{k}}{\partial x^{l}}+\Omega_{k} \Omega_{l}-\Omega_{l} \Omega_{k}=O \quad(1 \leqq k<l \leqq m) . \tag{4.3}
\end{equation*}
$$

The integrability condition (4.3) guarantees existence of the solution of (4.1) as follows
Theorem 4.5. Let $\Omega_{l}: U \rightarrow \mathrm{M}_{m}(\mathbb{R})(l=1, \ldots, n)$ be $C^{\infty}$-functions defined on a simply connected domain $U \subset \mathbb{R}^{n}$ satisfying (4.3) Then for each $\mathrm{P}_{0} \in U$ and $F_{0} \in \mathrm{M}_{m}(\mathbb{R})$, there exists the unique $m \times m$-matrix valued function $F: U \rightarrow \mathrm{M}_{m}(\mathbb{R})$ satisfying (4.1) and (4.2). Moreover,

- if $F_{0} \in \mathrm{GL}(m, \mathbb{R}), F(\mathrm{P}) \in \mathrm{GL}(m, \mathbb{R})$ holds on $U$,
- if $F_{0} \in \mathrm{SO}(m)$ and $\Omega_{l}$ 's are skew-symmetric matrices, $F(\mathrm{P}) \in \mathrm{SO}(m)$ holds on $U$.

Coordinate-free expressions Let $\Omega_{l}: U \rightarrow \mathrm{M}_{n}(\mathbb{R})(l=1, \ldots, m)$ be $C^{\infty}$-functions defined on a domain $U \subset \mathbb{R}^{m}$, and define $n \times n$-matrix $\Omega$ of 1 -forms as

$$
\Omega=\left(\begin{array}{cccc}
\omega_{1}^{1} & \omega_{2}^{1} & \ldots & \omega_{n}^{1}  \tag{4.4}\\
\omega_{1}^{2} & \omega_{2}^{2} & \ldots & \omega_{n}^{2} \\
\vdots & \vdots & \ddots & \vdots \\
\omega_{1}^{n} & \omega_{2}^{n} & \ldots & \omega_{n}^{n}
\end{array}\right):=\sum_{l=1}^{m} \Omega_{l} d x^{l}=\left(\begin{array}{cccc}
\sum \omega_{l, 1}^{1} d x^{l} & \sum \omega_{l, 2}^{1} d x^{l} & \ldots & \sum \omega_{l, n}^{1} d x^{l} \\
\sum \omega_{l, 1}^{l} d x^{l} & \sum \omega_{l, 2}^{2} d x^{l} & \ldots & \sum \omega_{l, n}^{2} d x^{l} \\
\vdots & \vdots & \ddots & \vdots \\
\sum \omega_{l, 1}^{n} d x^{l} & \sum \omega_{l, 2}^{n} d x^{l} & \ldots & \sum \omega_{l, n}^{n} d x^{l}
\end{array}\right),
$$

where $\Omega_{l}=\left(\omega_{l, j}^{i}\right)$. Then $\Omega$ is considered as a $\mathrm{M}_{n}(\mathbb{R})$-valued 1-form, and (4.1) is restated as

$$
\begin{equation*}
d F=F \Omega . \tag{4.5}
\end{equation*}
$$

Lemma 4.6. Under the situation above, the integrability condition (4.3) is equivalent to

$$
\begin{equation*}
d \Omega+\Omega \wedge \Omega=O, \quad \text { where } \quad \Omega \wedge \Omega=\left(\sum_{k=1}^{n} \omega_{k}^{i} \wedge \omega_{j}^{k}\right)_{i, j=1, \ldots, n} \tag{4.6}
\end{equation*}
$$

Proof. Assume $F$ be a solution of (4.5) with $F \in \mathrm{GL}(n, \mathbb{R})$. Then

$$
O=d d F=d(F \Omega)=d F \wedge \Omega+F d \Omega=F(\Omega \wedge \Omega+d \Omega)
$$

Thus, by using differential forms, we can state the system of partial differential equations (4.1) and its integrability condition (4.3) in coordinate-free form. The proof of Theorem 4.5 works not only simply connected domain $U \subset \mathbb{R}^{m}$ but also simply connected $m$-manifold, and thus, we have

Theorem 4.7. Let $\Omega$ be an $\mathrm{M}_{n}(\mathbb{R})$-valued 1 -form on a simply connected m-manifold $M$ satisfying (4.6). Then for each $\mathrm{P}_{0} \in M$ and $F_{0} \in \mathrm{M}_{n}(\mathbb{R})$, there exists the unique $n \times n$-matrix valued function $F: M \rightarrow \mathrm{M}_{n}(\mathbb{R})$ satisfying (4.5) with $F(\mathrm{P})=F_{0}$. Moreover,

- if $F_{0} \in \mathrm{GL}(n, \mathbb{R}), F(\mathrm{P}) \in \mathrm{GL}(n, \mathbb{R})$ holds on $M$,
- if $F_{0} \in \mathrm{SO}(n)$ and $\Omega$ is skew-symmetric, $F(\mathrm{P}) \in \mathrm{SO}(n)$ holds on $M$.

When $n=1$, that is, $\Omega$ is a usual 1-form, $\Omega \wedge \Omega$ always vanishes, and the integrability condition (4.6) is simply $d \Omega=0$. Then we have the following Poncaré's lemma ${ }^{8}$.

Theorem 4.8 (Poincaré's lemma). If a differential 1-form $\omega$ defined on a simply connected and connected m-manifold $M$ is closed, that is, $d \omega=0$ holds, then there exists a $C^{\infty}$-function $f$ on $M$ such that $d f=\omega$. Such a function $f$ is unique up to additive constants.
Proof. Since $\omega$ is closed, there exists a function $F$ on $M$ satisfying $d F=F \omega$ with initial condition $F\left(\mathrm{P}_{0}\right)=1$. By Lemma 4.3, $F$ does not vanish on $M$, that is, $F>0$. Hence $f:=\log F$ is a smooth function on $M$ satisfying $d f=d F / F=F \omega / F=\omega$. Take another function $g$ on $M$ satisfying $d g=\omega, d(f-g)=0$ holds. Then connectedness of $M$ infers that $f-g$ is constant.

### 4.3 Curvature form

Let $U$ be a domain of $n$-dimensional Riemannian manifold $(M, g)$. We let $\Omega$ be the connection form with respect to an orthonormal frame $\left[e_{1}, \ldots, e_{n}\right]$ on $U$, as defined in Definition 3.15.
Definition 4.9. We define a skew-symmetric matrix-valued 2-form by $K:=d \Omega+\Omega \wedge \Omega$ and call the curvature form with respect to the frame $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$.

Take an orthonormal frame $\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right]$ on $U$ and take a gauge transformation $\Theta: U \rightarrow \mathrm{O}(n)$ :

$$
\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]=\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right] \Theta
$$

Denoting the connection form and the curvature form with respect to $\left[\boldsymbol{v}_{j}\right]$ by $\widetilde{\Omega}$ and $\widetilde{K}$. Then
Proposition 4.10. (1) $\Omega=\Theta^{-1} \widetilde{\Omega} \Theta+\Theta^{-1} d \Theta$, (2) $K=\Theta^{-1} \widetilde{K} \Theta$.
Proof. Since

$$
\begin{aligned}
{\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right] \Omega } & =\nabla\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]=\nabla\left(\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right] \Theta\right)=\nabla\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right] \Theta+\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right] d \Theta \\
& =\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right] \widetilde{\Omega} \Theta+\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right] d \Theta=\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right] \Theta^{-1}(\widetilde{\Omega} \Theta+d \Theta)
\end{aligned}
$$

the first assertion is obtained. Next, noticing $d(\widetilde{\Omega} \Theta)=(d \widetilde{\Omega}) \Theta-\widetilde{\Omega} \wedge d \Theta, \widetilde{\Omega} \Theta^{-1} \wedge \Theta \widetilde{\Omega}=\widetilde{\Omega} \wedge \widetilde{\Omega}$, and so on, we have

$$
\begin{aligned}
d \Omega & +\Omega \wedge \Omega=d\left(\Theta^{-1} \widetilde{\Omega} \Theta+\Theta^{-1} d \Theta\right)+\left(\Theta^{-1} \widetilde{\Omega} \Theta+\Theta^{-1} d \Theta\right) \wedge\left(\Theta^{-1} \widetilde{\Omega} \Theta+\Theta^{-1} d \Theta\right) \\
= & -\Theta^{-1} d \Theta \Theta^{-1} \widetilde{\Omega} \Theta+\Theta^{-1} d \widetilde{\Omega} \Theta-\Theta^{-1} \widetilde{\Omega} \wedge d \Theta-\Theta^{-1} d \Theta \Theta^{-1} \wedge d \Theta \\
& +\Theta^{-1} \widetilde{\Omega} \Theta \wedge \Theta^{-1} \widetilde{\Omega} \Theta+\Theta^{-1} d \Theta \wedge \Theta^{-1} \widetilde{\Omega} \Theta+\Theta^{-1} \widetilde{\Omega} \Theta \wedge \Theta^{-1} d \Theta+\Theta^{-1} d \Theta \wedge \Theta^{-1} d \Theta \\
= & \Theta^{-1}(d \widetilde{\Omega}+\widetilde{\Omega} \wedge \widetilde{\Omega}) \Theta
\end{aligned}
$$

proving (2).

[^5]The goal of this section is to prove the following
Theorem 4.11. Let $U$ be a domain of a Riemannian n-manifold $(M, g)$ and $K$ the curvature form with respect to an orthonormal frame $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ on $U$. For a point $\mathrm{P} \in U$, there exists a local coordinate system $\left(x^{1}, \ldots, x^{n}\right)$ around P such that $\left[\partial / \partial x^{1}, \ldots, \partial / \partial x^{n}\right]$ is an orthonormal frame if and only if $K$ vanishes on a neighborhood of P .

Remark 4.12. By (2) of Proposition 4.10, the condition $K=0$ does not depend on choice of orthonormal frames. A Riemannian manifold $(M, g)$ said to be flat if $K=0$ holds on $M$.

Proof of Theorem 4.11. First, we shall show the "only if" part: Let $\left(x^{1}, \ldots, x^{n}\right)$ be a coordinate system such that $\left[\boldsymbol{e}_{j}:=\partial / \partial x^{j}\right]$ is an orthonormal frame. Since

$$
\left[\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right]=\left[\frac{\partial}{\partial x^{j}}, \frac{\partial}{\partial x^{k}}\right]=\mathbf{0}
$$

Proposition 4.2 yields that all components of the connection forms $\omega_{i}^{j}$ vanish. Hene we have $K=0$.
Conversely, assume $K=0$ for an orthonormal frame $\left[\boldsymbol{e}_{j}\right]$. Since the connection form $\Omega$ satisfies $d \Omega+\Omega \wedge \Omega=O$, there exists a matrix-valued function $\Theta: V \rightarrow \mathrm{SO}(n)$ satisfying $d \Theta=\Theta \Omega$, $\Theta(\mathrm{P})=$ id on a sufficiently small neighborhood $V$ of P , because of Theorem 4.5. Take a new orthonormal frame $\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right]:=\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right] \Theta^{-1}$. Then by (1) of Proposition 4.10, the connection form $\widetilde{\Omega}=\left(\tilde{\omega}_{i}^{j}\right)$ with respect to $\left[\boldsymbol{v}_{j}\right]$ vanishes identically. So by Lemma $3.17, d \omega^{i}=0$ holds for $i=1, \ldots, n$. Hence by the Poincaré Lemma (Theorem 4.8), there exists a smooth functions on a neighborhood $V$ of P . Such $\left(x^{1}, \ldots, x^{n}\right)$ is a desired coordinate system if $V$ is sufficiently small.

## Exercises

4-1 Consider a Riemannian metric

$$
g=d r^{2}+\{\varphi(r)\}^{2} d \theta^{2} \quad \text { on } \quad U:=\left\{(r, \theta) ; 0<r<r_{0},-\pi<\theta<\pi\right\}
$$

where $r_{0} \in(0,+\infty]$ and $\varphi$ is a positive smooth function defined on $\left(0, r_{0}\right)$ with

$$
\lim _{r \rightarrow+0} \varphi(r)=0, \quad \lim _{r \rightarrow+0} \varphi^{\prime}(r)=1
$$

Find a function $\varphi$ such that $(U, g)$ is flat. (Hint: $[\partial / \partial r,(1 / \varphi) \partial / \partial \theta)]$ is an orthonormal frame.)
4-2 Compute the curvature form of $H^{2}(-1)$ with respect to an orthonormal frame $\left[\boldsymbol{e}_{1}, \boldsymbol{e}_{2}\right]$ as in Exercise 3-2.

## 5 The Sectional Curvature

### 5.1 Preliminaries

Exterior products of tangent vectors. Let $V$ be an $n$-dimensional vector space $(1 \leqq n<\infty)$ and denote by $V^{*}$ its dual. Then $\left(V^{*}\right)^{*}$ can be naturally identified with $V$ itself. In fact,

$$
I: V \ni \boldsymbol{v} \longmapsto I \boldsymbol{v} \in\left(V^{*}\right)^{*}:=\left\{A: V^{*} \rightarrow \mathbb{R} ; \text { linear }\right\}, \quad I \boldsymbol{v}(\alpha):=\alpha(\boldsymbol{v})
$$

is a linear map with trivial kernel. Then $I$ is an isomorphism because $\operatorname{dim}\left(V^{*}\right)^{*}=\operatorname{dim} V$.
We denote by $\wedge^{2} V:=\wedge^{2}\left(V^{*}\right)^{*}$ the set of skew-symmetric bilinear forms on $V^{*}$. For vectors $\boldsymbol{v}$, $\boldsymbol{w} \in V$, the exterior product of them is an element of $\wedge^{2} V$ defined as

$$
(\boldsymbol{v} \wedge \boldsymbol{w})(\alpha, \beta):=\alpha(\boldsymbol{v}) \beta(\boldsymbol{w})-\alpha(\boldsymbol{w}) \beta(\boldsymbol{v}) \quad\left(\alpha, \beta \in V^{*}\right)
$$

For a basis $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ on $V$,

$$
\begin{equation*}
\left\{\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{j} ; 1 \leqq i<j \leqq n\right\} \tag{5.1}
\end{equation*}
$$

is a basis of $\wedge^{2} V$. In particular $\operatorname{dim} \wedge^{2} V=\frac{1}{2} n(n-1)$. When $V$ is a vector space endowed with an inner product $\langle$,$\rangle and \left[e_{1}, \ldots, e_{n}\right]$ is an orthonormal basis, there exists the unique inner product, which is also denoted by $\langle$,$\rangle , of \wedge^{2} V$ such that (5.1) is an orthonormal basis. This definition of the inner product does not depend on choice of orthonormal bases of $V$. In fact, take another orthonormal basis $\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right]$ related with $\left[\boldsymbol{e}_{j}\right]$ by

$$
\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]=\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right] \Theta \quad \Theta=\left(\theta_{i}^{j}\right) \in \mathrm{O}(n)
$$

Since $\Theta^{T}=\Theta^{-1},\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right]=\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right] \Theta^{T}$ holds. Hence

$$
\boldsymbol{v}_{s} \wedge \boldsymbol{v}_{t}=\left(\sum_{i} \theta_{s}^{i} \boldsymbol{e}_{i}\right) \wedge\left(\sum_{j} \theta_{t}^{j} \boldsymbol{e}_{j}\right)=\sum_{i, j} \theta_{i}^{s} \theta_{j}^{t}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{j}\right)=\sum_{i<j}\left(\theta_{i}^{s} \theta_{j}^{t}-\theta_{j}^{s} \theta_{i}^{t}\right)\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{j}\right)
$$

and so

$$
\begin{aligned}
\left\langle\boldsymbol{v}_{s} \wedge \boldsymbol{v}_{t}, \boldsymbol{v}_{u} \wedge \boldsymbol{v}_{v}\right\rangle & =\sum_{i<j, k<l}\left(\theta_{i}^{s} \theta_{j}^{t}-\theta_{j}^{s} \theta_{i}^{t}\right)\left(\theta_{k}^{u} \theta_{l}^{v}-\theta_{l}^{u} \theta_{k}^{v}\right)\left\langle\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{j}, \boldsymbol{e}_{k} \wedge \boldsymbol{e}_{l}\right\rangle \\
& =\sum_{i<j, k<l}\left(\theta_{i}^{s} \theta_{j}^{t}-\theta_{j}^{s} \theta_{i}^{t}\right)\left(\theta_{k}^{u} \theta_{l}^{v}-\theta_{l}^{u} \theta_{k}^{v}\right) \delta_{i k} \delta_{j l}=\sum_{i<j}\left(\theta_{i}^{s} \theta_{j}^{t}-\theta_{j}^{s} \theta_{i}^{t}\right)\left(\theta_{i}^{u} \theta_{j}^{v}-\theta_{j}^{u} \theta_{i}^{v}\right) \\
& =\sum_{i<j}\left(\theta_{i}^{s} \theta_{j}^{t} \theta_{i}^{u} \theta_{j}^{v}-\theta_{j}^{s} \theta_{i}^{t} \theta_{i}^{u} \theta_{j}^{v}-\theta_{i}^{s} \theta_{j}^{t} \theta_{j}^{u} \theta_{i}^{v}+\theta_{j}^{s} \theta_{i}^{t} \theta_{j}^{u} \theta_{i}^{v}\right) \\
& =\sum_{i<j} \theta_{i}^{s} \theta_{j}^{t} \theta_{i}^{u} \theta_{j}^{v}+\sum_{i<j} \theta_{j}^{s} \theta_{i}^{t} \theta_{i}^{u} \theta_{j}^{v}-\sum_{i>j} \theta_{j}^{s} \theta_{i}^{t} \theta_{i}^{u} \theta_{j}^{v}+\sum_{i>j} \theta_{i}^{s} \theta_{j}^{t} \theta_{i}^{u} \theta_{j}^{v} \\
& =\sum_{i \neq j} \theta_{i}^{s} \theta_{j}^{t} \theta_{i}^{u} \theta_{j}^{v}-\sum_{i \neq j} \theta_{j}^{s} \theta_{i}^{t} \theta_{i}^{u} \theta_{j}^{v} \\
& =\sum_{i, j}\left(\theta_{i}^{s} \theta_{j}^{t} \theta_{i}^{u} \theta_{j}^{v}-\theta_{j}^{s} \theta_{i}^{t} \theta_{i}^{u} \theta_{j}^{v}\right)-\sum_{i}\left(\theta_{i}^{s} \theta_{i}^{t} \theta_{i}^{u} \theta_{i}^{v}-\theta_{i}^{s} \theta_{i}^{t} \theta_{i}^{u} \theta_{i}^{v}\right) \\
& =\delta^{s u} \delta^{t v}-\delta^{t u} \delta^{s v}
\end{aligned}
$$

because $\sum_{i} \theta_{i}^{s} \theta_{i}^{t}=\delta^{s t}$. So, if $s<t$ and $u<v$, the second term of the right-hand side vanishes. That is, $\left\{\boldsymbol{v}_{s} \wedge \boldsymbol{v}_{t} ; s<t\right\}$ is an orthonormal basis as well as $\left\{\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{j} ; i<j\right\}$ is.
12. July, 2023. xRevised: 18. July, 2023

Symmetric bilinear forms. Let $V$ be a real vector space. A bilinear map $q: V \times V \rightarrow \mathbb{R}$ is said to be symmetric if $q(\boldsymbol{v}, \boldsymbol{w})=q(\boldsymbol{w}, \boldsymbol{v})$ for all $\boldsymbol{v}, \boldsymbol{w} \in V$.
Lemma 5.1. Two symmetric bilinear forms $q$ and $q^{\prime}$ coincide with each other if and only if $q(\boldsymbol{v}, \boldsymbol{v})=q^{\prime}(\boldsymbol{v}, \boldsymbol{v})$ hold for all $\boldsymbol{v} \in V$.

Proof. By symmetricity, $q(\boldsymbol{v}, \boldsymbol{w})=\frac{1}{2}(q(\boldsymbol{v}+\boldsymbol{w}, \boldsymbol{v}+\boldsymbol{w})-q(\boldsymbol{v}, \boldsymbol{v})-q(\boldsymbol{w}, \boldsymbol{w}))$ holds.

### 5.2 Sectional Curvature

Let $U$ be a domain on a Riemannian $n$-manifold $(M, g)$, and $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ an orthonormal frame on $U$. Denote by $\left(\omega^{j}\right)_{j=1, \ldots, n}, \Omega=\left(\omega_{i}^{j}\right)_{i, j=1, \ldots, n}$ and $K=\left(\kappa_{i}^{j}\right)_{i=1, \ldots, n}:=d \Omega+\Omega \wedge \Omega$ the dual frame, the connection form and the curvature form with respect to the frame $\left[\boldsymbol{e}_{j}\right]$. Then Lemma 3.17 and Definition 4.9, we have

$$
\begin{equation*}
d \omega^{j}=\sum_{l} \omega^{l} \wedge \omega_{l}^{j}, \quad \kappa_{i}^{j}=d \omega_{i}^{j}+\sum_{l} \omega_{l}^{j} \wedge \omega_{i}^{l} \tag{5.2}
\end{equation*}
$$

Since $\Omega$ is a one form valued in the skew-symmetric matrices, so is $K$ :

$$
\begin{equation*}
\omega_{i}^{j}=-\omega_{j}^{i}, \quad \kappa_{i}^{j}=-\kappa_{j}^{i} . \tag{5.3}
\end{equation*}
$$

Proposition 5.2 (The first Bianchi identity). $\kappa_{j}^{i}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right)+\kappa_{k}^{i}\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{j}\right)+\kappa_{l}^{i}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right)=0$.
Proof. By (5.2) and (3.11),

$$
\begin{aligned}
0 & =d d \omega^{i}=d\left(\sum_{s} \omega^{s} \wedge \omega_{s}^{i}\right)=\sum_{s}\left(d \omega^{s} \wedge \omega_{s}^{i}-\omega^{s} \wedge \omega_{s}^{i}\right) \\
& =\sum_{s}\left(\sum_{m}\left(\omega^{m} \wedge \omega_{m}^{s}\right) \wedge \omega_{s}^{i}-\omega^{s} \wedge\left(\kappa_{s}^{i}-\sum_{m} \omega_{m}^{i} \wedge d \omega_{s}^{m}\right)\right) \\
& =\sum_{s, m} \omega^{m} \wedge \omega_{m}^{s} \wedge \omega_{s}^{i}+\sum_{s, m} \omega^{s} \wedge \omega_{m}^{i} \wedge \omega_{s}^{m}-\sum_{s} \omega^{s} \wedge \kappa_{s}^{i} \\
& =\sum_{s, m} \omega^{m} \wedge\left(\omega_{m}^{s} \wedge \omega_{s}^{i}+\omega_{s}^{i} \wedge \omega_{m}^{s}\right)-\sum_{s} \omega^{s} \wedge \kappa_{s}^{i}=-\sum_{s} \omega^{s} \wedge \kappa_{s}^{i}
\end{aligned}
$$

Hence

$$
\begin{aligned}
0 & =\sum_{s}\left(\omega^{s} \wedge \kappa_{s}^{i}\right)\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right)=\sum_{s}\left(\omega^{s}\left(\boldsymbol{e}_{j}\right) \kappa_{s}^{i}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right)+\omega^{s}\left(\boldsymbol{e}_{k}\right) \kappa_{s}^{i}\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{j}\right)+\omega^{s}\left(\boldsymbol{e}_{l}\right) \kappa_{s}^{i}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right)\right) \\
& =\sum_{s}\left(\delta_{j}^{s} \kappa_{s}^{i}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right)+\delta_{k}^{s} \kappa_{s}^{i}\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{j}\right)+\delta_{l}^{s} \kappa_{s}^{i}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right)\right) \\
& =\kappa_{j}^{i}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right)+\kappa_{k}^{i}\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{j}\right)+\kappa_{l}^{i}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right)
\end{aligned}
$$

proving the assertion.
Corollary 5.3. $\kappa_{j}^{i}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right)=\kappa_{l}^{k}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right)$.
Proof. By Proposition 5.2,

$$
\begin{aligned}
\kappa_{j}^{i}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right)+\kappa_{k}^{i}\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{j}\right)+\kappa_{l}^{i}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{k}\right) & =0 \\
\kappa_{k}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{l}\right)+\kappa_{i}^{j}\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{k}\right)+\kappa_{l}^{j}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{i}\right) & =0 \\
\kappa_{i}^{k}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{l}\right)+\kappa_{j}^{k}\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{i}\right)+\kappa_{l}^{k}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right) & =0
\end{aligned}
$$

Summing up these and noticing $\kappa_{i}^{j}=-\kappa_{j}^{i}$, we have the conclusion.

A quadratic form induced from the curvature form. We fix a point $p \in U$. Under the notation above, we can define a bilinear map

$$
\begin{equation*}
\boldsymbol{K}(\boldsymbol{\xi}, \boldsymbol{\eta}):=\sum_{i<j, k<l} \kappa_{i}^{j}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right) \xi^{k l} \eta^{i j}, \quad \boldsymbol{\xi}=\sum_{k<l} \xi^{k l} \boldsymbol{e}_{k} \wedge \boldsymbol{e}_{l}, \quad \boldsymbol{\eta}=\sum_{i<j} \eta^{i j} \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{j} \tag{5.4}
\end{equation*}
$$

on $\wedge^{2} T_{p} M$, where $\boldsymbol{e}_{j}, \kappa_{i}^{j} \ldots$ are considered tangent vectors, 2 -forms at the fixed point $p$. In fact, one can show that the definition (5.4) is independent of choice of orthonormal frames. As a immediate conclusion of Corollary 5.3, we have

Lemma 5.4. $K$ is symmetric.
Hence, taking Lemma 5.1 into an account, we define the sectional curvature as follows:
Definition 5.5. Let $\Pi_{p} \subset T_{p} M$ be a 2-dimensional linear subspace in $T_{p} M$. The sectional curvature of $(M, g)$ with respect to the plane $\Pi_{p}$ is a number

$$
K\left(\Pi_{p}\right):=\boldsymbol{K}(\boldsymbol{v} \wedge \boldsymbol{w}, \boldsymbol{v} \wedge \boldsymbol{w})
$$

where $\{\boldsymbol{v}, \boldsymbol{w}\}$ is an orthonormal basis of $\Pi_{p}$
Remark 5.6. For (not necessarily orthonormal) basis $\{\boldsymbol{x}, \boldsymbol{y}\}$ of $\Pi_{p}$, the sectional curvature is expressed as

$$
K\left(\Pi_{p}\right)=\frac{\boldsymbol{K}(\boldsymbol{x} \wedge \boldsymbol{y}, \boldsymbol{x} \wedge \boldsymbol{y})}{\langle\boldsymbol{x} \wedge \boldsymbol{y}, \boldsymbol{x} \wedge \boldsymbol{y}\rangle}
$$

where $\langle$,$\rangle of the right-hand side is the inner product of \wedge^{2} T_{p} M$ induced from the Riemannian metric.
Remark 5.7. The sectional curvature is a scalar corresponding to a 2-plane in the tangent space $T_{p} M$. Hence it can be considered as a function of 2-Grassmannian bundle induced from the tangent bundle $T M$.

### 5.3 Curvature Tensor

Let $(M, g)$ be a Riemannian manifold and $\nabla$ the Levi-Civita connection. Define a trilinear map (5.5)

$$
R: \mathfrak{X}(M) \times \mathfrak{X}(M) \times \mathfrak{X}(M) \ni(X, Y, Z) \mapsto R(X, Y) Z=\nabla_{X} \nabla_{Y} Z-\nabla_{Y} \nabla_{X} Z-\nabla_{[X, Y]} Z \in \mathfrak{X}(M)
$$

By the properties Lemma 3.6 of the connection and the property (3.4) of the Lie bracket, the following Lemma is obvious.

Lemma 5.8. For any function $f \in \mathcal{F}(M)$ and vector fields $X, Y, Z \in \mathfrak{X}(M)$,

$$
R(f X, Y) Z=R(X, f Y) Z=R(X, Y)(f Z)=f R(X, Y) Z
$$

holds.
Corollary 5.9. Assume the vector fields $X, Y, Z$ and $\widetilde{X}, \widetilde{Y}, \widetilde{Z} \in \mathfrak{X}(M)$ satisfy $X_{p}=\widetilde{X}_{p}, Y_{p}=\widetilde{Y}_{p}$ and $Z_{p}=\widetilde{Z}_{p}$ for a point $p \in M$. Then

$$
(R(X, Y) Z)_{p}=(R(\widetilde{X}, \widetilde{Y}) \widetilde{Z})_{p}
$$

In other words, $R$ in (5.5) induces a trilinear map

$$
R_{p}: T_{p} M \times T_{p} M \times T_{p} M \rightarrow T_{p} M
$$

Definition 5.10. A trilinear map $R(X, Y) Z$ is called the curvature tensor of $(M, g)$. In addition, a quadrilinear map

$$
R(X, Y, Z, T)=\langle R(X, Y) Z, T\rangle: \mathfrak{X}(M)^{4} \rightarrow \mathcal{F}(M)
$$

is also called the curvature tensor. In fact, $R \in \Gamma\left(T^{*} M \otimes T^{*} M \otimes T^{*} M \otimes T^{*} M\right)$, that is $R$ is ( 0,4 )-tensor field, because $R$ induces a quadrilinear map

$$
R:\left(T_{p} M\right)^{4} \rightarrow \mathbb{R}
$$

for each $p \in M$.
Lemma 5.11. Let $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ be an orthonormal frame on a domain $U \subset M$, and $K=\left(\kappa_{i}^{j}\right)$ the curvature form with respect to the frame. Then it holds that

$$
\kappa_{i}^{j}(X, Y)=R\left(X, Y, \boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right)
$$

for each $(i, j)$.
So by (5.3), Proposition 5.2, Corollary 5.3 yield
Proposition 5.12. - $R(X, Y, Z, T)=-R(Y, X, Z, T)=-R(X, Y, T, Z)$,

- $R(X, Y, Z, T)+R(Y, Z, X, T)+R(Z, X, Y, T)=0$,
- $R(X, Y, Z, T)=R(Z, T, X, Y)$.

Moreover, the sectional curvature $K\left(\Pi_{p}\right)$ in Definition 5.5 is computed by

$$
\begin{equation*}
K\left(\Pi_{p}\right)=\frac{R(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{y}, \boldsymbol{x})}{\langle\boldsymbol{x}, \boldsymbol{x}\rangle\langle\boldsymbol{y}, \boldsymbol{y}\rangle-\langle\boldsymbol{x}, \boldsymbol{y}\rangle^{2}} \tag{5.6}
\end{equation*}
$$

## Exercises

5-1 Consider a Riemannian metric

$$
g=d r^{2}+\{\varphi(r)\}^{2} d \theta^{2} \quad \text { on } \quad U:=\left\{(r, \theta) ; 0<r<r_{0},-\pi<\theta<\pi\right\}
$$

where $r_{0} \in(0,+\infty]$ and $\varphi$ is a positive smooth function defined on $\left(0, r_{0}\right)$ with

$$
\lim _{r \rightarrow+0} \varphi(r)=0, \quad \lim _{r \rightarrow+0} \frac{\varphi(r)}{r}=1
$$

Classify the function $\varphi$ so that $g$ is of constant sectional curvature.
5-2 Let $M \subset \mathbb{R}^{n+1}$ be an embedded submanifold endowed with the Riemannian metric induced from the canonical Euclidean metric of $\mathbb{R}^{n+1}$. Then the position vector $\boldsymbol{x}(p)$ of $p \in M$ induces a smooth map

$$
\boldsymbol{x}: M \ni p \longmapsto \boldsymbol{x}(p) \in \mathbb{R}^{n+1}
$$

which is an $(n+1)$-tuple of $C^{\infty}$-functions. Let $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ be an orthonormal frame defined on a domain $U \subset M$. Since $T_{p} M \subset \mathbb{R}^{n+1}$, we can consider that $\boldsymbol{e}_{j}$ is a smooth map from $U \rightarrow \mathbb{R}^{n+1}$. Take a dual basis $\left(\omega^{j}\right)$ to $\left[\boldsymbol{e}_{j}\right]$. Prove that

$$
d \boldsymbol{x}=\sum_{j=1}^{n} \boldsymbol{e}_{j} \omega^{j}
$$

holds on $U$. Here, we regard that $d \boldsymbol{x}$ is an $(n+1)$-tuple of differential forms and $\boldsymbol{e}_{j}$ is an $\mathbb{R}^{n+1}$-valued function for each $j$.

## $6 \quad$ Space forms

### 6.1 Constant sectional curvature

Let $(M, g)$ be a Riemannian $n$-manifold, and let

$$
\begin{aligned}
\operatorname{Gr}_{2}(T M):= & \cup_{p} \operatorname{Gr}_{2}\left(T_{p} M\right) \\
& \operatorname{Gr}_{2}\left(T_{p} M\right):=2 \text {-Grassmannian of } T_{p} M=\left\{\Pi_{p} \subset T_{p} M ; \text { 2-dimensional subspace }\right\} .
\end{aligned}
$$

The sectional curvature defined in Definition 5.5 is a map $K: \operatorname{Gr}_{2}(T M) \rightarrow \mathbb{R}$ such that

$$
K\left(\Pi_{p}\right):=\boldsymbol{K}(\boldsymbol{v} \wedge \boldsymbol{w}, \boldsymbol{v} \wedge \boldsymbol{w})
$$

where $\{\boldsymbol{v}, \boldsymbol{w}\}$ is the orthonormal basis of $\Pi_{p}$.
Fix a point $p$, and take an orthornormal frame $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ defined on a neighborhood $U$ of $p$. Denote by $\left(\omega^{j}\right), \Omega=\left(\omega_{i}^{j}\right)$ and $K=\left(\kappa_{i}^{j}\right)$ the dual frame, the connection form and the curvature form with respect to the frame $\left[\boldsymbol{e}_{j}\right]$, respectively.
Theorem 6.1. Assume there exists a real number $k$ such that $K\left(\Pi_{p}\right)=k$ for all 2-dimensional subspace $\Pi_{p} \in T_{p} M$ for a fixed $p$. Then the curvature form is expressed as

$$
\kappa_{j}^{i}=k \omega^{i} \wedge \omega^{j}
$$

Conversely, the curvature form is written as above, the sectional curvature at $p$ is constant $k$.
Proof. By the assumption, $k=K\left(\operatorname{Span}\left\{\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right\}\right)=\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{j}, \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{j}\right)=\kappa_{j}^{i}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right)$. Let

$$
\boldsymbol{v}:=\cos \theta \boldsymbol{e}_{i}+\sin \theta \boldsymbol{e}_{j}, \quad \boldsymbol{w}:=\cos \varphi \boldsymbol{e}_{l}+\sin \varphi \boldsymbol{e}_{m}
$$

where $\{i, j\} \neq\{l, m\}$, and set $\Pi_{\theta, \varphi}:=\operatorname{Span}\{\boldsymbol{v}, \boldsymbol{w}\} \subset T_{p} M$. Then by biliniearity of the $\wedge$-product on $T_{p} M$, it holds that

$$
\boldsymbol{v} \wedge \boldsymbol{w}=\cos \theta \cos \varphi \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}+\cos \theta \sin \varphi \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}+\sin \theta \cos \varphi \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}+\sin \theta \sin \varphi \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}
$$

Since $\{\boldsymbol{v}, \boldsymbol{w}\}$ is an orthonormal basis of $\Pi_{\theta, \varphi}$, biliniearity and symmetricity of $\boldsymbol{K}$ implies

$$
\begin{align*}
k= & K\left(\Pi_{\theta, \varphi}\right)=\boldsymbol{K}(\boldsymbol{v} \wedge \boldsymbol{w}, \boldsymbol{v} \wedge \boldsymbol{w})  \tag{6.1}\\
= & \cos ^{2} \theta \cos ^{2} \varphi \boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}\right)+\cos ^{2} \theta \sin ^{2} \varphi \boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}, \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}\right) \\
& +\sin ^{2} \theta \cos ^{2} \varphi \boldsymbol{K}\left(\boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}\right)+\sin ^{2} \theta \sin ^{2} \varphi \boldsymbol{K}\left(\boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right) \\
& +2 \cos ^{2} \theta \cos \varphi \sin \varphi \boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}\right)+2 \cos \theta \sin \theta \cos ^{2} \varphi \boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}\right) \\
& +2 \cos \theta \sin \theta \cos \varphi \sin \varphi\left(\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right)+\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}\right)\right) \\
& +2 \cos \theta \sin \theta \sin ^{2} \varphi \boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right)+2 \sin ^{2} \theta \cos \varphi \sin \varphi \boldsymbol{K}\left(\boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right) \\
= & k+2\left(\cos ^{2} \theta \cos \varphi \sin \varphi \boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}\right)+\cos \theta \sin \theta \cos ^{2} \varphi \boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}\right)\right. \\
& +\cos \theta \sin \theta \cos \varphi \sin \varphi\left(\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right)+\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}\right)\right) \\
& \left.+\cos \theta \sin \theta \sin ^{2} \varphi \boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right)+\sin ^{2} \theta \cos \varphi \sin \varphi \boldsymbol{K}\left(\boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right)\right)
\end{align*}
$$

So, by letting $\theta=0$, we have

$$
\begin{equation*}
\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}\right)=0 \tag{6.2}
\end{equation*}
$$

Similarly, letting $\theta=\pi / 2, \varphi=0$ and $\varphi=\pi / 2$, we have $\boldsymbol{K}\left(\boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right)=\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}\right)=$ $\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right)=0$. Hence the equality (6.1) implies

$$
\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{l}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{m}\right)+\boldsymbol{K}\left(\boldsymbol{e}_{i} \wedge \boldsymbol{e}_{m}, \boldsymbol{e}_{j} \wedge \boldsymbol{e}_{l}\right)=0
$$

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By definition (5.4), this is equivalent to

$$
\kappa_{j}^{m}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{l}\right)+\kappa_{j}^{l}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{m}\right)=-\left(\kappa_{m}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{l}\right)+\kappa_{l}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{m}\right)\right) .
$$

Then by Proposition 5.2, we have

$$
0=\kappa_{m}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{l}\right)+\kappa_{l}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{m}\right)=\kappa_{m}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{l}\right)-\kappa_{i}^{j}\left(\boldsymbol{e}_{m}, \boldsymbol{e}_{l}\right)-\kappa_{m}^{j}\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{i}\right)=2 \kappa_{m}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{l}\right)-\kappa_{i}^{j}\left(\boldsymbol{e}_{m}, \boldsymbol{e}_{l}\right)
$$

Exchanging the roles of $i$ and $m$, it holds that $2 \kappa_{i}^{j}\left(\boldsymbol{e}_{m}, \boldsymbol{e}_{l}\right)-\kappa_{m}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{l}\right)=0$. So we have

$$
\begin{equation*}
\kappa_{i}^{j}\left(\boldsymbol{e}_{m}, \boldsymbol{e}_{l}\right)=0 \quad(\text { if }\{i, j\} \neq\{m, l\}) \tag{6.3}
\end{equation*}
$$

On the other hand, (6.2) means that $\kappa_{i}^{j}\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{l}\right)=\kappa_{i}^{j}\left(\boldsymbol{e}_{j}, \boldsymbol{e}_{l}\right)=0$ when $l \neq i, j$. Summing up, we have

$$
\kappa_{i}^{j}\left(\boldsymbol{e}_{k}, \boldsymbol{e}_{l}\right)= \begin{cases}k & (i, j)=(k, l) \\ 0 & \text { otherwise }\end{cases}
$$

proving the theorem.
We now consider the case that the assumption of Theorem 6.1 holds for each $p \in M$.
Theorem 6.2. Assume that for each $p$, there exists a real number $k(p)$ such that $K\left(\Pi_{p}\right)=k(p)$ for any $\Pi_{p} \in \operatorname{Gr}_{2}\left(T_{p} M\right)$. Then the function $k: M \ni p \rightarrow k(p) \in \mathbb{R}$ is constant provided that $M$ is connected.

Proof. By taking the exterior derivative of $\kappa_{i}^{j}=d \omega_{i}^{j}+\sum_{s} \omega_{s}^{j} \wedge \omega_{i}^{s}$, it holds that

$$
\begin{aligned}
d \kappa_{i}^{j} & =d\left(d \omega_{i}^{j}\right)+\sum_{s} \omega_{s}^{j} \wedge d \omega_{i}^{s}-\sum_{s} d \omega_{s}^{j} \wedge \omega_{i}^{s} \\
& =\sum_{s}\left(\kappa_{s}^{j}-\sum_{t} \omega_{t}^{j} \wedge \omega_{s}^{t}\right) \wedge \omega_{i}^{s}-\sum_{s} \omega_{s}^{j} \wedge\left(\kappa_{i}^{s}-\sum_{t} \omega_{t}^{s} \wedge \omega_{i}^{t}\right)
\end{aligned}
$$

and hence we have the identity

$$
\begin{equation*}
d \kappa_{i}^{j}=\sum_{s}\left(\kappa_{s}^{j} \wedge \omega_{i}^{s}-\omega_{s}^{j} \wedge \kappa_{i}^{s}\right), \tag{6.4}
\end{equation*}
$$

which is known as the second Bianchi identity. By our assumption, Theorem 6.1 implies that $\kappa_{i}^{j}=k \omega^{i} \wedge \omega^{j}$. Then by Lemma 3.17,

$$
\begin{aligned}
d \kappa_{i}^{j} & =d\left(k \omega^{i}\right) \wedge \omega^{j}-k \omega^{i} \wedge d \omega^{j}=d k \wedge \omega^{i} \wedge \omega^{j}+k d \omega^{i} \wedge \omega^{j}-k \omega^{i} \wedge d \omega^{j} \\
& =d k \wedge \omega^{i} \wedge \omega^{j}+\sum_{s} k \omega^{s} \wedge \omega_{s}^{i} \wedge \omega^{j}-\sum_{s} k \omega^{i} \wedge \omega^{s} \wedge \omega_{s}^{j}=d k \wedge \omega^{i} \wedge \omega^{j}+d \kappa_{i}^{j}
\end{aligned}
$$

holds for each $i$ and $j$. Thus, $d k \wedge \omega^{i} \wedge \omega^{j}=0$ for all $i$ and $j$, which implies $d k=0$. This equality is independent of choice of orthonormal frames. Since $M$ is connected, $k$ is constant.

### 6.2 Space forms

Let $(M, g)$ be a Riemannian $n$-manifold. A path $\gamma:[0,+\infty) \rightarrow M$ is said to be a divergence path if for any compact subset $K \in M$, there exists $t_{0} \in(0,+\infty)$ such that $\gamma\left(\left[t_{0},+\infty\right)\right) \subset M \backslash K$. If any divergent path has infinite length, $(M, g)$ is said to be complete. ${ }^{9}$ In particular, a compact Riemannian manifold without boundary is automatically complete.

[^6]Definition 6.3. An $n$-dimensional space form is a complete Riemannian $n$-manifold of constant sectional curvature.

Example 6.4. The Euclidean $n$-space is a space form of constant sectional curvature 0 . In fact, let $\left(x^{1}, \ldots, x^{n}\right)$ be the canonical Cartesian coordinate system and set $\boldsymbol{e}_{j}=\partial / \partial x^{j}$. Then [ $\boldsymbol{e}_{j}$ ] is an orthornormal frame defined on the entire $\mathbb{R}^{n}$, and Propositions 4.1 and 4.2 implies that the connection form $\omega_{j}^{i}=0$. Hence the curvature forms vanish, and then the sectional curvature is identically zero.

So it is sufficient to show completeness. Let $\gamma:[0,+\infty) \rightarrow \mathbb{R}^{n}$ be a divergent path. Then for each $r>0$, there exists $t_{0}>0$ such that $|\gamma(t)|>r$ holds on $\left[t_{0},+\infty\right)$, equivalently, $|\gamma(t)| \rightarrow+\infty$ as $t \rightarrow+\infty$. So the length $L$ of the curve $\gamma$ is

$$
L=\lim _{t \rightarrow+\infty} \int_{0}^{t}|\dot{\gamma}(\tau)| d \tau \geqq \lim _{t \rightarrow+\infty}\left|\int_{0}^{t} \dot{\gamma}(\tau) d \tau\right|=\lim _{t \rightarrow+\infty}|\gamma(t)-\gamma(0)| \geqq \lim _{t \rightarrow+\infty}|\gamma(t)|-|\gamma(0)|=+\infty
$$

Here, we used the triangle inequality of integrals for vector-valued functions ${ }^{10}$.

### 6.3 The Hyperbolic spaces

Let $H^{n}\left(-c^{2}\right)$ be the hyperbolic $n$-space defined, where $c$ is a non-zero constant:

$$
H^{n}\left(-c^{2}\right):=\left\{\boldsymbol{x}=\left(x^{0}, \ldots, x^{n}\right) \in \mathbb{R}_{1}^{n+1} \left\lvert\,\langle\boldsymbol{x}, \boldsymbol{x}\rangle_{L}=-\frac{1}{c^{2}}\right., c x_{0}>0\right\}
$$

where $\left(\mathbb{R}_{1}^{n+1},\langle,\rangle_{L}\right)$ be the Lorentz-Minkowski $(n+1)$-space. The tangent space $T_{\boldsymbol{x}} H^{n}\left(-c^{2}\right)$ is the orthogonal complement $\boldsymbol{x}^{\perp}$ of $\boldsymbol{x}$, and the restriction $g_{H}$ of the inner product $\langle,\rangle_{L}$ to $T_{\boldsymbol{x}} H^{n}\left(-c^{2}\right)$ is positive definite. Thus, $\left(H^{n}\left(-c^{2}\right), g_{H}\right)$ is a Riemannian manifold, called the hyperbolic n-space.

Theorem 6.5. The hyperbolic space $\left(H^{n}\left(-c^{2}\right), g_{H}\right)$ is of constant sectional curvature $-c^{2}$.
Proof. Notice that $H^{n}\left(-c^{2}\right)$ can be expressed as a graph $x^{0}=\frac{1}{c} \sqrt{1+c^{2}\left(\left(x^{1}\right)^{2}+\cdots+\left(x^{n}\right)^{2}\right)}$ defined on the $\left(x^{1}, \ldots, x^{n}\right)$-hyperplane, that is, it is covered by single chart. Then there exists a orthonormal frame field $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ defined on entire $H^{n}\left(-c^{2}\right)$. Denote by $\left(\omega^{i}\right), \Omega=\left(\omega_{i}^{j}\right)$ and $K=\left(\kappa_{i}^{j}\right)$ the dual frame, the connection form and the curvature form with respect to $\left[\boldsymbol{e}_{j}\right]$, respectively.

Regarding $T_{\boldsymbol{x}} H^{n}\left(-c^{2}\right)$ as a linear subspace in $\mathbb{R}_{1}^{n+1}$, we can consider $\boldsymbol{e}_{j}$ as a vector-valued function. In addition the position vector $\boldsymbol{x} \in H^{n}\left(-c^{2}\right)$ can be also regarded as a vector-valued function. Since $T_{\boldsymbol{x}} H^{n}\left(-c^{2}\right)=\boldsymbol{x}^{\perp}$,

$$
\begin{equation*}
\mathcal{F}:=\left(\boldsymbol{e}_{0}, \boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right): H^{n}\left(-c^{2}\right) \rightarrow \mathrm{M}_{n+1}(\mathbb{R}) \quad \boldsymbol{e}_{0}=c \boldsymbol{x} \tag{6.5}
\end{equation*}
$$

gives a pseudo orthornormal frame along $H^{n}\left(-c^{2}\right)$, that is, $\mathcal{F}^{T} Y \mathcal{F}=Y(Y:=\operatorname{diag}(-1,1, \ldots, 1))$ holds.

As seen in Exercise 5-2, it holds that

$$
\begin{equation*}
d \boldsymbol{e}_{0}=c d \boldsymbol{x}=c \sum_{j=1}^{n} \omega^{j} \boldsymbol{e}_{j} \tag{6.6}
\end{equation*}
$$

On the other hand, for each $j=1, \ldots, n$, decompose the vector-valued one form $d \boldsymbol{e}_{j}$ as

$$
d \boldsymbol{e}_{j}=h_{j} \boldsymbol{e}_{0}+\sum_{s} \alpha_{j}^{s} \boldsymbol{e}_{s}
$$

[^7]where $h_{j}$ and $\alpha_{j}^{s}$ are one forms on $H^{n}\left(-c^{2}\right)$. Here,
$$
h_{j}=-\left\langle d \boldsymbol{e}_{j}, \boldsymbol{e}_{0}\right\rangle_{L}=-d\left\langle\boldsymbol{e}_{j}, \boldsymbol{e}_{0}\right\rangle_{L}+\left\langle\boldsymbol{e}_{j}, d \boldsymbol{e}_{0}\right\rangle_{L}=c \omega^{j}
$$
and
$$
\alpha_{j}^{s}=\left\langle d \boldsymbol{e}_{j}, \boldsymbol{e}_{s}\right\rangle_{L}=d\left\langle\boldsymbol{e}_{j}, \boldsymbol{e}_{s}\right\rangle_{L}-\left\langle\boldsymbol{e}_{j}, d \boldsymbol{e}_{s}\right\rangle_{L}=-\alpha_{s}^{j} .
$$

Differentiating (6.6), it holds that

$$
0=\frac{1}{c} d d \boldsymbol{e}_{0}=\sum_{j}\left(d \omega^{j} \boldsymbol{e}_{j}-\omega^{j} \wedge d \boldsymbol{e}_{j}\right)=\sum_{j, s} \omega^{s} \wedge \omega_{s}^{j} \boldsymbol{e}_{j}-\sum_{j, s} \omega^{j} \wedge \alpha_{j}^{s} \boldsymbol{e}_{s}=\sum_{j} \sum_{s} \omega^{s} \wedge\left(\omega_{s}^{j}-\alpha_{s}^{j}\right) \boldsymbol{e}_{j}
$$

because $\omega^{j} \wedge \omega^{j}=0$. Thus, we have $\sum_{s} \omega^{s} \wedge\left(\omega_{s}^{j}-\alpha_{s}^{j}\right)=0$, and then

$$
\begin{aligned}
& 0=\left(\sum_{s} \omega^{s} \wedge\left(\omega_{s}^{j}-\alpha_{s}^{j}\right)\right)\left(\boldsymbol{e}_{l}, \boldsymbol{e}_{m}\right)=\left(\omega_{l}^{j}\left(\boldsymbol{e}_{m}\right)-\alpha_{l}^{j}\left(\boldsymbol{e}_{m}\right)\right)-\left(\omega_{m}^{j}\left(\boldsymbol{e}_{l}\right)-\alpha_{m}^{j}\left(\boldsymbol{e}_{l}\right)\right), \\
& 0=\left(\omega_{j}^{m}\left(\boldsymbol{e}_{l}\right)-\alpha_{j}^{m}\left(\boldsymbol{e}_{l}\right)\right)-\left(\omega_{l}^{m}\left(\boldsymbol{e}_{j}\right)-\alpha_{l}^{m}\left(\boldsymbol{e}_{j}\right)\right)=-\left(\omega_{m}^{j}\left(\boldsymbol{e}_{l}\right)-\alpha_{m}^{j}\left(\boldsymbol{e}_{l}\right)\right)-\left(\omega_{l}^{m}\left(\boldsymbol{e}_{j}\right)-\alpha_{l}^{m}\left(\boldsymbol{e}_{j}\right)\right), \\
& 0=\left(\omega_{m}^{l}\left(\boldsymbol{e}_{j}\right)-\alpha_{m}^{l}\left(\boldsymbol{e}_{j}\right)\right)-\left(\omega_{j}^{l}\left(\boldsymbol{e}_{m}\right)-\alpha_{j}^{l}\left(\boldsymbol{e}_{m}\right)\right)=-\left(\omega_{l}^{m}\left(\boldsymbol{e}_{j}\right)-\alpha_{l}^{m}\left(\boldsymbol{e}_{j}\right)\right)+\left(\omega_{l}^{j}\left(\boldsymbol{e}_{m}\right)-\alpha_{l}^{j}\left(\boldsymbol{e}_{m}\right)\right),
\end{aligned}
$$

which conclude that $\omega_{l}^{j}=\alpha_{l}^{j}$. Summing up, we have

$$
\begin{equation*}
d \boldsymbol{e}_{j}=c \omega^{j} \boldsymbol{e}_{0}+\sum_{s} \omega_{j}^{s} \boldsymbol{e}_{s} \tag{6.7}
\end{equation*}
$$

Then the frame $\mathcal{F}$ in (6.5) satisfies

$$
d \mathcal{F}=\mathcal{F} \widetilde{\Omega}, \quad \text { where } \quad \widetilde{\Omega}=\left(\begin{array}{cc}
0 & c \boldsymbol{\omega}^{T}  \tag{6.8}\\
c \boldsymbol{\omega} & \Omega
\end{array}\right) \quad \text { and } \quad \boldsymbol{\omega}:=\left(\omega^{1}, \ldots, \omega^{n}\right)^{T} .
$$

The integrability condition of (6.8) is

$$
O=d \widetilde{\Omega}+\widetilde{\Omega} \wedge \widetilde{\Omega}=\left(\begin{array}{cc}
c^{2} \boldsymbol{\omega}^{T} \wedge \boldsymbol{\omega} & c\left(d \boldsymbol{\omega}^{T}+\omega^{T} \wedge \Omega\right) \\
c(d \boldsymbol{\omega}+\Omega \wedge \boldsymbol{\omega}) & d \Omega+\Omega \wedge \Omega+c^{2} \boldsymbol{\omega} \wedge \boldsymbol{\omega}^{T}
\end{array}\right)
$$

The lower-right components of the identity above yields

$$
\kappa_{i}^{j}+c^{2} \omega^{i} \wedge \omega^{j}=0
$$

Hence the sectional curvature of $\left(H^{n}\left(-c^{2}\right), g_{H}\right)=-c^{2}$.
Remark 6.6. One can show the completeness of $\left(H^{n}\left(-c^{2}\right), g_{H}\right)$ (cf. MTH.B505). Hence the hyperbolic space is a simply connected space form of constant negative sectional curvature.

### 6.4 Isometries

A $C^{\infty}$-map $f: M \rightarrow N$ between manifolds $M$ and $N$ induces a linear map

$$
(d f)_{p}: T_{p} M \ni X \longmapsto(d f)_{p}(X)=\left.\frac{d}{d t}\right|_{t=0} f \circ \gamma(t) \in T_{f(p)} N
$$

where $\gamma:(-\varepsilon, \varepsilon) \rightarrow M$ is a smooth curve with $\gamma(0)=p$ and $\dot{\gamma}(0)=X$, called the differential of $f$. Since $p \in M$ is arbitrary, this induces a bundle homomorphism $d f: T M \rightarrow T N$.

Definition 6.7. A vector field on $N$ along a smooth map $f: M \rightarrow N$ is a map $X: M \rightarrow T N$ satisfying $\pi \circ X=f$, where $\pi: T N \rightarrow N$ is the canonical projection.

Then for each vector field $X \in \mathfrak{X}(M), d f(X)$ is a vector field on $N$ along $f$.
Definition 6.8. A $C^{\infty}$-map $f: M \rightarrow N$ between Riemannian manifolds $(M, g)$ and $(N, h)$ is called a local isometry if $\operatorname{dim} M=\operatorname{dim} N$ and $f^{*} h=g$ hold, that is,

$$
f^{*} h(X, Y):=h(d f(X), d f(Y))=g(X, Y)
$$

holds for $X, Y \in T_{p} M$ and $p \in M$.
Lemma 6.9. A local isometry is an immersion.
Proof. Let $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ be a (local) orthonormal frame of $M$, where $n=\operatorname{dim} M$. Set $\boldsymbol{v}_{j}:=d f\left(\boldsymbol{e}_{j}\right)$ $(j=1, \ldots, n)$ for a smooth map $f:(M, g) \rightarrow(N, h)$. If $f$ is a local isometry, $\left[\boldsymbol{v}_{1}(p), \ldots, \boldsymbol{v}_{n}(p)\right]$ is an orthonormal system in $T_{f(p)} N$, because

$$
h\left(\boldsymbol{v}_{i}, \boldsymbol{v}_{j}\right)=h\left(d f\left(\boldsymbol{e}_{i}\right), d f\left(\boldsymbol{e}_{j}\right)\right)=f^{*} h\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right)=g\left(\boldsymbol{e}_{i}, \boldsymbol{e}_{j}\right)
$$

Hence the differential $(d f)_{p}$ is of rank $n$.
The proof of Lemma 6.9 suggests the following fact:
Corollary 6.10. A smooth map $f:(M, g) \rightarrow(N, h)$ is a local isometry if and only if for each $p \in M$,

$$
\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right]:=\left[d f\left(\boldsymbol{e}_{1}\right), \ldots, d f\left(\boldsymbol{e}_{n}\right)\right]
$$

is an orthonormal frame for some orthonormal frame $\left[\boldsymbol{e}_{j}\right]$ on a neighborhood of $p$.

### 6.5 Local uniqueness of space forms

Theorem 6.11. Let $U \subset \mathbb{R}^{n}$ be a simply connected domain and $g$ a Riemannian metric on $U$. If the sectional curvature of $(U, g)$ is constant $k$, there exists a local isometry $f: U \rightarrow N^{n}(k)$, where

$$
N^{n}(k)= \begin{cases}S^{n}(k) & (k>0) \\ \mathbb{R}^{n} & (k=0) \\ H^{n}(k) & (k<0)\end{cases}
$$

Proof. Take an orthonormal frame $\left[\boldsymbol{e}_{1}, \ldots, \boldsymbol{e}_{n}\right]$ on $U$, and let $\left(\omega^{j}\right), \Omega=\left(\omega_{i}^{j}\right)$ and $K=\left(\kappa_{i}^{j}\right)$ be the dual frame, the connection form, and the curvature form with respect to [ $\boldsymbol{e}_{j}$ ], respectively. Since the sectional curvature is constant $k, \kappa_{i}^{j}=k \omega^{i} \wedge \omega^{j}$ holds for each $(i, j)$, because of Theorem 6.1.

First, consider the case $k=0$ : In this case, $K=d \Omega+\Omega \wedge \Omega=O$, and then by Theorem 4.5, there exists the unique matrix valued function $\mathcal{F}: U \rightarrow \mathrm{SO}(n)$ satisfying

$$
d \mathcal{F}=\mathcal{F} \Omega, \quad \mathcal{F}\left(p_{0}\right)=\mathrm{id}
$$

where $p_{0} \in U$ is a fixed point. Decompose the matrix $\mathcal{F}$ into column vectors as $\mathcal{F}=\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right]$, and define an $\mathbb{R}^{n}$-valued one form

$$
\boldsymbol{\alpha}:=\sum_{j=1}^{n} \omega^{j} \boldsymbol{v}_{j}
$$

Then

$$
d \boldsymbol{\alpha}=\sum_{j=1}^{n}\left(d \omega^{j} \boldsymbol{v}_{j}-\omega^{j} \wedge d \boldsymbol{v}_{j}\right)=\sum_{j, s}\left(\omega^{s} \wedge \omega_{s}^{j}\right) \boldsymbol{v}_{j}-\sum_{j, s}\left(\omega^{j} \wedge \omega_{j}^{s}\right) \boldsymbol{v}_{s}=\mathbf{0}
$$

Hence by the Poincaré lemma (Theorem 4.8), there exists a smooth map $f: U \rightarrow \mathbb{R}^{n}$ satisfying $d f=\boldsymbol{\alpha}$. For such an $f$, it holds that

$$
d f\left(\boldsymbol{e}_{s}\right)=\alpha\left(\boldsymbol{e}_{s}\right)=\sum_{j=1}^{n} \omega^{j}\left(\boldsymbol{e}_{s}\right) \boldsymbol{v}_{j}=\boldsymbol{v}_{s}
$$

for $s=1, \ldots, n$. Hence $\left[d f\left(\boldsymbol{e}_{1}\right), \ldots, d f\left(\boldsymbol{e}_{n}\right)\right]=\left[\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right]$ is an orthonormal frame, and then $f$ is a local isometry because Corollary 6.10.

Next, consider the case $k=-c^{2}<0$. We set

$$
\widetilde{\Omega}:=\left(\begin{array}{cc}
0 & c \boldsymbol{\omega}^{T} \\
c \boldsymbol{\omega} & \Omega
\end{array}\right), \quad \text { where } \quad \boldsymbol{\omega}=\left(\begin{array}{c}
\omega^{1} \\
\vdots \\
\omega^{n}
\end{array}\right)
$$

as in (6.8) in Section ??. Since $\kappa_{i}^{j}=k \omega^{i} \wedge \omega^{j}=-c^{2} \omega^{i} \wedge \omega^{j}, d \widetilde{\Omega}+\widetilde{\Omega} \wedge \widetilde{\Omega}=O$ holds as seen in Section ??. Hence there exists an matrix valued function $\mathcal{F}: U \rightarrow \mathrm{M}_{n+1}(\mathbb{R})$ satisfying

$$
\begin{equation*}
d \mathcal{F}=\mathcal{F} \widetilde{\Omega}, \quad \mathcal{F}\left(p_{0}\right)=\mathrm{id} \tag{6.9}
\end{equation*}
$$

where $p_{0} \in U$ is a fixed point. Notice that

$$
\widetilde{\Omega}^{T} Y+Y \widetilde{\Omega}=O \quad Y=\left(\begin{array}{cccc}
-1 & 0 & \ldots & 0 \\
0 & 1 & \ldots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \ldots & 1
\end{array}\right)
$$

holds,

$$
d\left(\mathcal{F} Y \mathcal{F}^{T}\right)=\mathcal{F} \widetilde{\Omega} Y \mathcal{F}^{T}+\mathcal{F} Y \widetilde{\Omega}^{T} \mathcal{F}^{T}=\mathcal{F}\left(\widetilde{\Omega} Y+Y \widetilde{\Omega}^{T}\right) \mathcal{F}^{T}=O
$$

Hence, by the initial condition,

$$
\mathcal{F} Y \mathcal{F}^{T}=Y, \quad \text { that is, } \quad(\mathcal{F} Y)^{-1}=\mathcal{F}^{T} Y
$$

Thus, we have

$$
\begin{equation*}
\mathcal{F}^{T} Y \mathcal{F}=(\mathcal{F} Y)^{-1} \mathcal{F}=Y \mathcal{F}^{-1} \mathcal{F}=Y \tag{6.10}
\end{equation*}
$$

Decompose $\mathcal{F}=\left[\boldsymbol{v}_{0}, \boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{n}\right]$. Then (6.10) is equivalent to

$$
\begin{equation*}
-\left\langle\boldsymbol{v}_{0}, \boldsymbol{v}_{0}\right\rangle_{L}=\left\langle\boldsymbol{v}_{1}, \boldsymbol{v}_{1}\right\rangle_{L}=\cdots=\left\langle\boldsymbol{v}_{n}, \boldsymbol{v}_{n}\right\rangle_{L}=1, \quad\left\langle\boldsymbol{v}_{i}, \boldsymbol{v}_{j}\right\rangle=0 \quad(\text { if } i \neq j) \tag{6.11}
\end{equation*}
$$

In particular, the 0 -th component of $\boldsymbol{v}_{0}$ never vanishes, since

$$
-1=\left\langle\boldsymbol{v}_{0}, \boldsymbol{v}_{0}\right\rangle_{L}=-\left(v_{0}^{0}\right)^{2}+\left(v_{0}^{1}\right)^{2}+\cdots+\left(v_{0}^{n}\right)^{2} \quad \boldsymbol{v}_{0}=\left(v_{0}^{0}, v_{0}^{1}, \ldots, v_{0}^{n}\right)^{T} .
$$

Moreover, by the initial condition $\boldsymbol{v}_{0}\left(p_{0}\right)=(1,0, \ldots, 0)^{T}$,

$$
\begin{equation*}
v_{0}^{0}>0 \tag{6.12}
\end{equation*}
$$

holds.
Set $f:=\frac{1}{c} \boldsymbol{v}_{0}$. Then $f: U \rightarrow \mathbb{R}_{1}^{n+1}$ is the desired map. In fact, by (6.11) and (6.12),

$$
f \in H^{n}\left(-c^{2}\right)=\left\{\boldsymbol{x}=\left(x^{0}, \ldots, x^{n}\right)^{T} \in \mathbb{R}_{1}^{n+1} \left\lvert\,\langle\boldsymbol{x}, \boldsymbol{x}\rangle=-\frac{1}{c^{2}}\right., c x^{0}>0\right\}
$$

and

$$
d f\left(\boldsymbol{e}_{j}\right)=\frac{1}{c} d \boldsymbol{v}_{0}\left(\boldsymbol{e}_{j}\right)=\sum_{s=1}^{n} \omega^{s}\left(\boldsymbol{e}_{j}\right) \boldsymbol{v}_{s}=\boldsymbol{v}_{j}
$$

Hence $\left[\boldsymbol{v}_{j}\right]=\left[\boldsymbol{e}_{j}\right]$ is an orthonormal frame because (6.11).
The case $k>0$ is left as an exercise.

## Exercises

6-1 Prove that the sphere

$$
S^{3}(1)=\left\{\boldsymbol{x} \in \mathbb{R}^{4} ;\langle\boldsymbol{x}, \boldsymbol{x}\rangle=1\right\}
$$

of radius 1 in the Euclidean 4-space is of constant sectional curvature 1.
6-2 Prove Theorem 6.11 for $k=1$ and $n=2$, assuming Exercise 6-1.

## Bibliography

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[Tu17] L. W. Tu, Differential geometry, Springer Verlag, 2017.
[UY17] Masaaki Umehara and Kotaro Yamada, Differential geometry of curves and surfaces, World Scientific, 2017.

## Glossary

1－form 1－形式， 1 次微分形式， 14
affirm connection アファイン接続， 16
arc－length parameter 弧長径数， 7
bilinear 双線形， 15
Cauchy－Riemann equations コーシー・リーマン方程式， 12
column vector 列ベクトル， 3
compatibility condition 適合条件， 9
conjugate 共役， 13
covariant tensor 共変テンソル， 14
covariant 共変， 14
curvature tensor 曲率テンソル， 26
curvature 曲率， 7
dual space 双対空間， 14
eigenvalue 固有値， 3
exterior derivative 外微分， 16
exterior product 外積， 23
flat 平坦， 22
form（微分）形式， 15
Frenet frame フルネ枠， 7
gauge transformation ゲージ変換， 17
general linear group $(\mathrm{GL}(n, \mathbb{R}))$ 一般線形群， 3
harmonic function 調和関数， 12
holomorphic 正則（複素関数が）， 12
initial value problem 初期値問題， 1
integrability condition 可積分条件， 9
Laplacian ラプラシアン， 12
Levi－Suavity connection レビ・チビタ接続， 16
Lie algebra リー代数， 14
Lie bracket リー括弧積， 14
linear connection 線形接続， 16
linear function 1 次関数， 2
linear ordinary differential equation 線形常微分方程式， 2
ordinary differential equation 常微分方程式， 1 orthogonal group $(\mathrm{O}(n))$ 直交群， 4
orthonormal frame 直交枠， 16
partial differential equation 偏微分方程式， 9
regular curve 正則曲線， 7
regular matrix 正則行列， 3
Riemannian connection リーマン接続， 16
second Bianchi identity 第二ビアンキ恒等式， 28
sectional curvature 断面曲率， 25
simply connected 単連結， 10,20
skew－symmeetric matrix 交代行列，歪対称行列， 4
skew－symmetric 交代的，反対称， 15
solution 解， 1
space curve 空間曲線， 7
space form 空間形， 29
special linear group $(\mathrm{SL}(n, \mathbb{R}))$ 特殊線形群， 4 special orthogonal group $(\mathrm{SO}(n))$ 特殊直交群， 4
tensor テンソル， 14
torsion 据率， 7
trilinear 三重線形， 15
unknown function 未知関数， 1


[^0]:    ${ }^{1}|X|_{\mathrm{M}}>0$ whenever $X \neq O,|\alpha X|_{\mathrm{M}}=|\alpha||X|_{\mathrm{M}}$, and the triangle inequality.

[^1]:    ${ }^{2} \mathrm{GL}(n, \mathbb{R})=\left\{A \in \mathrm{M}_{n}(\mathbb{R}) ; \operatorname{det} A \neq 0\right\}$ : the general linear group.
    ${ }^{3}$ In this lecture, id denotes the identity matrix.

[^2]:    ${ }^{4} \mathrm{SL}(n, \mathbb{R})=\left\{A \in \mathrm{M}_{n}(\mathbb{R}) ; \operatorname{det} A=1\right\}$; the special lienar group.
    ${ }^{5} \mathrm{O}(n)=\left\{A \in \mathrm{M}_{n}(\mathbb{R}) ; A^{T} A=A A^{T}=\mathrm{id}\right\}:$ the orthogonal group; $\mathrm{SO}(n)=\{A \in \mathrm{O}(n) ; \operatorname{det} A=1\}$ : the special orthogonal group.

[^3]:    ${ }^{6}$ The theorem holds under the assumption of $C^{2}$-differentiablity.

[^4]:    4. July, 2023. Revised: 11. July, 2023
    ${ }^{7} \mathrm{GL}(n, \mathbb{R})$ denotes the set of $n \times n$-regular matrices.
[^5]:    ${ }^{8}$ Theorem 2.6 in Advanced Topics in Geometry E (MTH.B501).

[^6]:    ${ }^{9}$ Usually, completeness is defined in terms of geodesics: A Riemannian manifold $(M, g)$ is complete if any geodesics are defined on entire $\mathbb{R}$. The definition here is one of the equivalent conditions of completeness, expressed in the Hopf-Rinow theorem. cf. MTH.B505.

[^7]:    ${ }^{10}$ See, for example, Theorem A.1.4 in [UY17] for $n=2$. The idea of the proof works for general $n$.

