# Convex extension of the Busemann-Hausdorff area integrand and the Plateau problem in arbitrary codimension 

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This is a revised second version. In case there are any typos or mistakes with regard to contents, please do not hesitate to contact me at pistre@instmath.rwth-aachen.de.
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## Zusammenfassung

Diese Masterarbeit verallgemeinert ein Result über die Minimierer des Busemann-Hausdorff Flächenfunktionals auf beliebige Kodimension. Im Jahr 2014 wurde das Plateau Problem im Finslerraum $\left(\mathbb{R}^{3}, F\right)$ von Overath und von der Mosel gelöst, [OvdM14]. Die vorliegende Arbeit zeigt, dass die Beweistechnik von Overath und von der Mosel erweitert werden kann, um die Existenz von flächenminimierenden Oberflächen im Finslerraum $\left(\mathbb{R}^{n}, F\right)$ nachzuweisen, falls die Finslermetrik reversibel ist. Um dieses Ziel zu erreichen, wird Burago and Ivanov's Arbeit [BI12] über die Konvexität der zwei-dimensionalen Busemann-Hausdorff Flächendichte ausführlich diskutiert. Abschließend wird ein Zusammenhang zwischen dem Busemann-Hausdorff Flächenintegranden und der Theorie über Cartan Integranden aufgezeigt.


#### Abstract

This thesis generalises a result on minimisers of the Busemann-Hausdorff area functional to arbitrary codimension. Recently, Overath and von der Mosel solved the Plateau problem for three-dimensional Finsler space $\left(\mathbb{R}^{3}, F\right)$, [OvdM14]. The present work shows that their proof technique can be extended to show the existence of area minimising surfaces in $n$-dimensional Finsler space $\left(\mathbb{R}^{n}, F\right)$ for reversible Finsler metrics. To achieve this goal, this thesis extensively discusses Burago and Ivanov's work [B112] on the convexity of the two-dimensional Busemann-Hausdorff area density. Finally, a connection of the Busemann-Hausdorff area integrand and the well-investigated theory of Cartan integrands is illustrated.


## About the second version

This is a revised second version, dating 31.12.2016. Apart from typing errors, the result of Lemma 1.1.14 was slightly improved. In case there are any further typos or mistakes with regard to contents, please do not hesitate to contact me at pistre@instmath.rwth-aachen.de.

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## Introduction

In this thesis we investigate the Plateau problem in $n$-dimensional Finsler space for reversible Finsler metrics $F$.

The Plateau problem (named in honour of the Belgian physicist J.A.F. Plateau (1801-1883)) is one example of a boundary value problem for minimal surfaces. In 1873, Plateau conducted a number of soap film experiments during which he noted that every single closed wire, however complicated its geometric form may be, bounds at least one soap film. By Johan Bernoulli's principle of virtual work, soap films in stable equilibrium correspond to surfaces of minimal surface area. Mathematically, a closed wire can be modelled by a closed rectifiable Jordan curve. In addition, one can prove that surfaces of minimal area are minimal surfaces (surfaces whose mean curvature vanishes). The question which the Plateau problem poses is the following (see [DHKW92, Chapter 4, pp. 221-226]):

Given a closed rectifiable Jordan curve $\Gamma$, is there a minimal surface spanned by $\Gamma$ ?

The first question one needs to answer is how to "measure" surface area. This depends on the geometric setting one studies the surfaces in. For the purpose of this thesis, we consider surfaces immersed into a Finsler space and consequently the notion of Busemann-Hausdorff area (or Finsler area) introduced by Busemann in [Bus47] (see Section 2.1). Secondly, one needs to discuss the class of permitted surfaces. For example, we do not consider fractal surfaces. Then we want to solve the variational problem

$$
\operatorname{minimise} \operatorname{area}_{B}^{F}(X)
$$

over a class of admissible surfaces $X \in \mathcal{C}(\Gamma)$ (which will be defined in Section 3.1).
In [OvdM14], Overath and von der Mosel showed the existence of Finsler area minimisers in codimension one and established higher regularity of solutions. However, Overath and von der Mosel considered a more general class of Finsler metrics whose " $m$-harmonic symmetrisation" is a Finsler metric as well. In their work they applied the theory of Cartan functionals to Finsler area. For this, assume $X: \mathcal{M} \rightarrow \mathbb{R}^{n}$ is a smooth immersion of a smooth m-manifold into $\mathbb{R}^{n}$ and
$I \in C^{0}\left(\mathbb{R}^{n} \times \mathbb{R}^{N}\right), N=\binom{n}{m}$ is positively homogeneous in its second argument, that is,

$$
I(x, t y)=t I(x, y)
$$

for all $(x, y) \in \mathbb{R}^{n} \times \mathbb{R}^{N}$ and $t>0$. Therefore, the $m$-form

$$
i(p)=I\left(X(p), d X_{p}\left(\left.\frac{\partial}{\partial u^{1}}\right|_{p}\right) \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{2}}\right|_{p}\right) \wedge \cdots \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{m}}\right|_{p}\right)\right) .
$$

for a local coordinate chart $\left(u^{\alpha}\right)_{\alpha=1}^{m}$ on $\mathcal{M}$ is globally well-defined. Then

$$
\mathcal{I}(X)=\int_{p \in \mathcal{M}} i(p)
$$

is called Cartan functional and $I$ is the corresponding Cartan integrand. A broad range of results in the theory of Cartan functionals has been established. Hildebrandt and von der Mosel, in particular, showed the existence and regularity of Cartan minimisers in a certain class of admissible surfaces ([HvdM03b, HvdM03c]). Their result holds for Cartan integrands which are positive definite, that is,

$$
M_{1}|z| \leq I(x, z) \leq M_{2}|z|
$$

for all $(x, z) \in \mathbb{R}^{n} \times \mathbb{G C}_{m}\left(\mathbb{R}^{n}\right)$ and semi-elliptic, that is,

$$
I\left(x, t z_{1}+(1-t) z_{2}\right) \leq t I\left(x, z_{1}\right)+(1-t) I\left(x, z_{2}\right)
$$

for all $x \in \Omega, z_{1}, z_{2} \in \mathbb{R}^{N}$ and $t \in[0,1]$ (see Section 3.2).
In codimension one $(n=m+1)$, Overath and von der Mosel identified the two-dimensional Busemann-Hausdorff area integrand as a positive definite, semi-elliptic Cartan integrand so that the existence result for Cartan minimisers yields a solution to the Plateau problem in Finsler space. The crucial part for the positive definiteness (see Theorem 3.3.5) is a representation of the BusemannHausdorff area integrand as a spherical integral (see Theorem 3.3.2) found by Overath in [Ove13] which holds in arbitrary codimension. The semi-ellipticity of the Busemann-Hausdorff area integrand for reversible Finsler metrics has been proved by Busemann in [Bus49] - but only in codimension one.

Recently, Burago and Ivanov [BI12] established this semi-ellipticity in arbitrary codimension (see Theorem 2.2.2). Their work uses concepts from multilinear algebra, convex geometry and convex analysis. Essentially, they showed that the convexity condition for the Busemann-Hausdorff area integrand is equivalent to the existence of so-called calibrators which support the area integrand "in
every direction" (see Lemma 2.2.4). Subsequently, they reformulate the notion of a calibrator into an inequality regarding the Euclidean area of centrally symmetric two-dimensional polygons on the plane (see inequality (2.2.22)). This inequality is proved with elementary results from convex geometry and convex analysis (see Theorem 2.2.7).

Therefore, in the last chapter we combine Overath and von der Mosel's proof technique with the result found by Burago and Ivanov and solve the Plateau problem in Finsler space in arbitrary codimension for a reversible Finsler metric

The thesis is outlined as follows. Chapter 1 begins with basic definitions and results on multilinear algebra, convex and differential geometry. These results are grouped into sections according to the respective topics. In particular, we introduce in Section 1.1 the algebraic notions of the tensor product and exterior power of a vector space. Subsequently, the Plücker embedding and some properties are presented. The section concludes with a Riesz-type isomorphism that arises in an inner product space. In Section 1.2 we illustrate basic principles of convex sets and polytopes such as the polarity of polytopes and polyhedral sets and the support function of a convex set. The next subsections cover properties of the mixed volumes of convex sets and we develop a well-known explicit formula for a certain mixed volume. Finally, a maximum principle for convex functions over convex domains is presented. Section 1.3 introduces basic notions of differential geometry and serves as a reference section for the subsequent chapters.

Chapter 2 covers an extensive treatment of Burago and Ivanov's work [BI12] on the convexity of the two-dimensional Busemann-Hausdorff area density. First, we formally introduce the BusemannHausdorff definition of volume on a Finsler manifold. This extends to the notion of the Finsler area functional $\operatorname{area}_{\Omega}^{F}$ of immersed $m$-dimensional submanifolds. Secondly, we define the convexity of an area density and present Burago and Ivanov's reformulation thereof. Using the concepts introduced in Chapter 1, we prove the central result on polygons described above (see Theorem 2.2.7).

In Chapter 3 we formulate the Plateau problem in Finsler space for arbitrary codimension (see Theorem 3.1.1). Subsequently, we give basic notions and results of Cartan functional theory. In Section 3.3 we identify the Busemann-Hausdorff area integrand as a Cartan integrand which leads to a solution of the Plateau problem.

## Chapter 1

## Preliminaries of multilinear algebra, convex geometry and differential geometry

In this chapter we want to give an introduction to different concepts that are needed in the subsequent chapters. The focus here lies on results of multilinear algebra and convex geometry, as they are used extensively in Burago and Ivanov's work on the convexity of the two-dimensional Busemann-Hausdorff area density [BI12]. We present their work in greater detail in the next chapter. Differential geometric concepts are stated merely for reference and we omit most of the proofs in the respective section.

### 1.1 Multilinear algebra

This section aims to introduce some results of multilinear algebra. We begin by algebraically defining the tensor product of a finite number of vector spaces. This leads to elementary results for the exterior power $\Lambda^{m}(V)$ of a vector space. In addition, the Plücker embedding arises, relating $m$-dimensional subspaces to a certain subset of the $m$ th exterior power. By imposing an inner product structure on the vector space $V$, we introduce the notion of a volume form and prove a Riesz-type theorem. The results and notation in this section have mainly been borrowed from the treatise of Lee [Lee13, Chapter 12]. The section on the Plücker embedding is based on the book of Harris [Har95].

Some words to clarify notation. The set $\mathbb{N}$ is the set of natural numbers $\{1,2,3, \ldots\}$ and $\mathbb{R}$ is the set of real numbers. All vector spaces are to be considered over the field of real numbers. In some parts, we commit the mild sin of identifying a linear map with its matrix representation with respect to a basis. In addition, we usually write $L v$ instead of $L(v)$ for the image of a vector $v$ under a linear map $L$. A matrix $A \in \mathbb{R}^{m \times n}$ is represented as $\left(A_{j}^{i}\right)=\left(A_{j}^{i}\right)_{i, j}=\left(A_{j}^{i}\right)_{i=1, \ldots, m, j=1, \ldots, n}$ or $\left(A_{j}^{i}\right)_{i, j=1}^{m}$ if $m=n$. Therein $i$ is the row and $j$ the column index.

### 1.1.1 The tensor product

Suppose $V_{1}, V_{2}, \ldots, V_{m}$ are finite-dimensional $\mathbb{F}$-vector spaces of dimensions $n_{1}, n_{2}, \ldots, n_{m}$, respectively. Recall that the cartesian product $V_{1} \times V_{2} \times \cdots \times V_{m}$ turns into a vector space if it is endowed with component-wise addition and scalar multiplication given by $a\left(v_{1}, v_{2}, \ldots, v_{m}\right)=\left(a v_{1}, a v_{2}, \ldots, a v_{m}\right)$. The cartesian product is of dimension $n_{1}+n_{2}+\ldots+n_{m}$. Our aim is to construct a new vector space $V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ of (finite) dimension $n_{1} \cdot n_{2} \cdot \ldots \cdot n_{m}$ which consists of linear combinations of objects of the form $v_{1} \otimes v_{2} \otimes \cdots \otimes v_{m}$ where $v_{i} \in V_{i}$ and in such a manner that $v_{1} \otimes v_{2} \otimes \cdots \otimes v_{m}$ depends linearly on each $v_{i}$ separately (as opposed to the cartesian product where the scalar multiplication is not homogeneous in each entry separately). A natural algebraic way to achieve this goal is to construct this new vector space as a certain quotient vector space. We begin by forming the free vector space with basis $V_{1} \times V_{2} \times \cdots \times V_{m}$, that is, define the set of formal sums

$$
\mathcal{F}(N):=\left\{\sum_{n \in N} a_{n} n \mid a_{n} \in K, a_{n}=0 \text { for all but finitely many } n \in N\right\}
$$

where $N:=V_{1} \times V_{2} \times \cdots \times V_{m}$. We can define an addition and a scalar multiplication on $\mathcal{F}(N)$ by

$$
\begin{aligned}
\left(\sum_{n \in N} a_{n} n\right)+\left(\sum_{n \in N} b_{n} n\right) & :=\sum_{n \in N}\left(a_{n}+b_{n}\right) n \\
\alpha\left(\sum_{n \in N} b_{n} n\right) & :=\sum_{n \in N}\left(\alpha a_{n}\right) n
\end{aligned}
$$

This turns $\mathcal{F}(N)$ into a vector space over the base field $\mathbb{F}$ of the vector spaces $V_{i}$. Take note that by definition any element of $N$ is a basis element. In fact, the free vector space above is of dimension $\prod_{i=1}^{m} \# \mathbb{F} \cdot \operatorname{dim}\left(V_{i}\right)=(\# \mathbb{F})^{m} \cdot \prod_{i=1}^{m} \operatorname{dim}\left(V_{i}\right)$. Here $\# S$ denotes the cardinality of a set $S$. Therefore, when $\mathbb{F}=\mathbb{R}$ this space is far too large for our intended purpose of creating a vector space of finite dimension $n_{1} \cdot n_{2} \cdot \ldots \cdot n_{m}$. For this reason, we proceed by factoring out the linearity relations we are looking for. Consider the subspace $U$ of $\mathcal{F}(N)$ given by

$$
\operatorname{span}\left\{\left(v_{1}, \ldots, a v_{i}+w_{i}, \ldots, v_{m}\right)-a\left(v_{1}, \ldots, v_{i}, \ldots, v_{m}\right)-\left(v_{1}, \ldots, w_{i}, \ldots, v_{m}\right) \mid a \in \mathbb{R}, v_{i}, w_{i} \in V_{i}\right\} .
$$

The tensor product of $\boldsymbol{V}_{\mathbf{1}}, \boldsymbol{V}_{\mathbf{2}}, \ldots, \boldsymbol{V}_{\boldsymbol{m}}$, denoted as $V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$, is defined as the quotient space

$$
V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}:=\mathcal{F}\left(V_{1} \times V_{2} \times \cdots \times V_{m}\right) / U
$$

The equivalence class of an element $\left(v_{1}, v_{2}, \ldots, v_{m}\right)$ in $V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ will be denoted by

$$
v_{1} \otimes v_{2} \otimes \cdots \otimes v_{m}:=\left(v_{1}, v_{2}, \ldots, v_{m}\right)+U
$$

and is called the tensor product of $\boldsymbol{v}_{1}, \boldsymbol{v}_{2}, \ldots, \boldsymbol{v}_{\boldsymbol{m}}$. By construction, the tensor product satisfies our linearity requirements

$$
v_{1} \otimes \cdots \otimes\left(a v_{i}+w_{i}\right) \otimes \cdots \otimes v_{m}=a\left(v_{1} \otimes \cdots \otimes v_{i} \otimes \cdots \otimes v_{m}\right)+v_{1} \otimes \cdots \otimes w_{i} \otimes \cdots \otimes v_{m}
$$

and any element of $V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ can be written as a linear combination of elements of the form $v_{1} \otimes v_{2} \otimes \cdots \otimes v_{m}$ where $v_{i} \in V_{i}$. As a real quotient space the tensor product is itself a real vector space (see e.g. [DF04, Theorem 7, p. 412]).

A map $F: V_{1} \times V_{2} \times \cdots \times V_{m} \rightarrow X$ into a vector space $X$ is called multilinear if it is linear as a function of each entry separately, that is, if for each $i=1, \ldots, m$

$$
F\left(v_{1}, \ldots, a v_{i}+w_{i}, \ldots, v_{m}\right)=a F\left(v_{1}, \ldots, v_{i}, \ldots, v_{m}\right)+F\left(v_{1}, \ldots, w_{i}, \ldots, v_{m}\right)
$$

The tensor product satisfies the following universal property. A proof can be found in [Lee13, Proposition 12.7, p. 309] but it will be omitted here.

Proposition 1.1.1 (Universal Property of the Tensor Product Space)
Let $V_{1}, V_{2}, \ldots, V_{m}$ be finite-dimensional real vector spaces. If $A: V_{1} \times V_{2} \times \cdots \times V_{m} \rightarrow X$ is a multilinear map into a vector space $X$, then there is a unique linear map $\widetilde{A}: V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m} \rightarrow X$ such that $A=\widetilde{A} \circ \pi$, where $\pi$ is the projection map given by $\pi\left(v_{1}, v_{2}, \ldots, v_{m}\right):=v_{1} \otimes v_{2} \otimes \cdots \otimes v_{m}$. That is, the following diagram commutes:


The strength of the universal property and the reason for its name lie in the fact that it uniquely determines the tensor product up to unique isomorphism. Some texts also define a tensor product by its universal property and subsequently use the construction shown above to prove the existence of a tensor product (which by the following uniqueness result can then be called the tensor product).

## Proposition 1.1.2

Let $V_{1}, V_{2}, \ldots, V_{m}$ be finite-dimensional real vector spaces and suppose $\pi^{\prime}: V_{1} \times V_{2} \times \ldots \times V_{m} \rightarrow Z$ is a multilinear map into a vector space $Z$ with the following universal property.
For any multilinear map $B: V_{1} \times V_{2} \times \ldots \times V_{m} \rightarrow Y$ there is a unique linear map $\widetilde{B}: Z \rightarrow Y$ such that $B=\widetilde{B} \circ \pi^{\prime}$. That is, the following diagram commutes:


Then there is a unique isomorphism $\Phi: V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m} \rightarrow Z$ such that $\pi^{\prime}=\Phi \circ \pi$ where $\pi: V_{1} \times V_{2} \times \cdots \times V_{m} \rightarrow V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ is the canonical projection.

## Remark:

The preceding proposition shows that the details of the construction of the tensor product space are irrelevant as long as the resulting space - in the notation of Proposition 1.1.2 this is $Z$ - satisfies the universal property.

Proof of Proposition 1.1.2: The proof will be given by diagram chasing. First, we use Proposition 1.1.1 for $X=Z$ and $A=\pi^{\prime}$. This gives a unique linear map $\widetilde{A}: V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m} \rightarrow Z$ such that $\widetilde{A} \circ \pi=\pi^{\prime}$. Similarly, the universal property for $Z$ from the hypothesis applied to the vector space $Y=V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ and the multilinear map $B=\pi$ yields a unique linear map $\widetilde{B}: Z \rightarrow V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ such that $\widetilde{B} \circ \pi^{\prime}=\pi$. In turn, the linear composite map $C:=\widetilde{A} \circ \widetilde{B}: Z \rightarrow Z$ satisfies $C \circ \pi^{\prime}=\widetilde{A} \circ \widetilde{B} \circ \pi^{\prime}=\widetilde{A} \circ \pi=\pi^{\prime}$. If we reapply the universal property for $Z$ from the hypothesis for $Y=Z$ and $B=\pi^{\prime}$ we see that $C$ is the unique linear map such that $C \circ \pi^{\prime}=\pi^{\prime}$. Of course, the identity map $\operatorname{Id}_{Z}: Z \rightarrow Z$ is also such a map and so by uniqueness we find that $C=\operatorname{Id}_{Z}$. Analogously (using Proposition 1.1.1 again for $X=V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ and $A=\pi$ ), one obtains $\widetilde{B} \circ \widetilde{A}=\operatorname{Id}_{V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}}$. Therefore, $\widetilde{A}$ and $\widetilde{B}$ are inverse to each other and setting $\Phi:=\widetilde{A}$ proves the claim. Since $\widetilde{A}$ is unique by the universal property, so is $\Phi$.

To establish the dimension of the tensor product $V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ we need the following result which involves a choice of basis for each component vector space $V_{i}$. In fact, in the above construction the vector spaces need not be finite-dimensional. The following will be the first explicit reference to a finite basis. We will omit the proof again because it only involves standard arguments from linear
algebra and usage of the universal property. Essentially, we can form a basis for the tensor product space by taking all possible tensor products of basis vectors of the component spaces.

Proposition 1.1.3 ([Lee13, Proposition 12.8, p. 309])
Let $V_{1}, V_{2}, \ldots, V_{m}$ be finite-dimensional real vector spaces of dimensions $n_{1}, n_{2}, \ldots, n_{m}$ respectively. Suppose $\left\{e_{1}^{(j)}, e_{2}^{(j)}, \ldots, e_{n_{j}}^{(j)}\right\}$ is a basis for $V_{j}$ for each $j=1, \ldots, m$. Then the set

$$
\mathcal{C}:=\left\{e_{i_{1}}^{(1)} \otimes e_{i_{2}}^{(2)} \otimes \cdots \otimes e_{i_{m}}^{(m)} \mid 1 \leq i_{j} \leq n_{j}, j=1, \ldots, m\right\}
$$

is a basis for $V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$. Therefore, the dimension of $V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ is

$$
\operatorname{dim}\left(V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}\right)=n_{1} n_{2} \cdots n_{m}
$$

Let us denote the set of multilinear functions $F: V_{1} \times V_{2} \times \cdots \times V_{m} \rightarrow \mathbb{R}$ by $L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)$. It turns into a vector space if endowed with the usual pointwise addition and scalar multiplication

$$
\begin{aligned}
\left(F+F^{\prime}\right)\left(v_{1}, \ldots, v_{i}, \ldots, v_{m}\right) & :=F\left(v_{1}, \ldots, v_{i}, \ldots, v_{m}\right)+F^{\prime}\left(v_{1}, \ldots, v_{i}, \ldots, v_{m}\right) \\
(a F)\left(v_{1}, \ldots, v_{i}, \ldots, v_{m}\right) & :=a\left(F\left(v_{1}, \ldots, v_{i}, \ldots, v_{m}\right)\right)
\end{aligned}
$$

The dual space of a vector space $V$ consists of all linear functionals $\omega: V \rightarrow \mathbb{R}$ and will be denoted by $V^{*}$. Its elements are called linear forms or 1-forms. When dealing with indexed sets of elements of $V$ and $V^{*}$ we will use lower indices to represent vectors and upper indices to represent linear forms. Let $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ be a basis for $V$. The set $\left\{\varepsilon^{1}, \varepsilon^{2}, \ldots, \varepsilon^{n}\right\} \subset V^{*}$ defined by $\varepsilon^{i}\left(e_{j}\right):=\delta_{j}^{i}$ forms a basis for $V^{*}$ and is called the dual basis to the initial basis for $V$. Here $\delta_{j}^{i}$ is the Kronecker delta, that is, $\delta_{j}^{i}=1$ if $i=j$ and 0 otherwise. To avoid cluttering of indices we will occasionally denote the dual basis by $\left\{e_{1}^{*}, e_{2}^{*}, \ldots, e_{n}^{*}\right\}$ instead.

For $\omega^{i} \in V_{i}^{*}, i=1, \ldots, m$ define a map $\omega^{1} \otimes \omega^{2} \otimes \cdots \otimes \omega^{m}: V_{1} \times V_{2} \times \cdots \times V_{m} \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
\left(\omega^{1} \otimes \omega^{2} \otimes \cdots \otimes \omega^{m}\right)\left(v_{1}, v_{2}, \ldots, v_{m}\right):=\omega^{1}\left(v_{1}\right) \omega^{2}\left(v_{2}\right) \cdots \omega^{m}\left(v_{m}\right) . \tag{1.1.1}
\end{equation*}
$$

We call this map the concrete tensor product of $\omega^{1}, \omega^{2}, \ldots, \omega^{m}$. Note that the concrete tensor product is well-defined. This function is multilinear because $\mathbb{R}$ is a field, that is,

$$
\omega^{1} \otimes \omega^{2} \otimes \cdots \otimes \omega^{m} \in L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)
$$

In fact, the space of multilinear functions is spanned be such elements.

Lemma 1.1.4 ([Lee13, Proposition 12.4, p. 306])
Let $V_{1}, V_{2}, \ldots, V_{m}$ be finite-dimensional real vector spaces of dimensions $n_{1}, n_{2}, \ldots, n_{m}$, respectively. Suppose $\left\{e_{1}^{(j)}, e_{2}^{(j)}, \ldots, e_{n_{j}}^{(j)}\right\}$ is a basis for $V_{j}$ and by $\left\{\varepsilon_{(j)}^{1}, \varepsilon_{(j)}^{2}, \ldots, \varepsilon_{(j)}^{n_{j}}\right\}$ denote the corresponding dual basis for $V_{j}^{*}$ for each $j=1, \ldots, m$. Then the set

$$
\mathcal{C}^{\prime}:=\left\{\varepsilon_{(1)}^{i_{1}} \otimes \varepsilon_{(2)}^{i_{2}} \otimes \cdots \otimes \varepsilon_{(m)}^{i_{m}} \mid 1 \leq i_{j} \leq n_{j}, j=1, \ldots, m\right\}
$$

is a basis for $L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)$.

Thus, we conclude that a multilinear function $F \in L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)$ is uniquely determined by its values $F_{i_{1} i_{2} \ldots i_{m}}:=F\left(e_{i_{1}}^{(1)}, e_{i_{2}}^{(2)}, \ldots, e_{i_{m}}^{(m)}\right)$ on a basis - analogously to linear functions. By Lemma 1.1.4 and Proposition 1.1.3 both vector spaces $L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)$ and $V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m}$ are of the same dimension and therefore isomorphic. Naturally, the question arises whether there is an even stronger connection between these spaces that respects the algebraic structure, that is, the multilinearity of the elements. Indeed, the following result holds.

Lemma 1.1.5 ([Lee13, Proposition 12.10, p. 311])
If $V_{1}, V_{2}, \ldots, V_{m}$ are finite-dimensional vector spaces, there is a canonical isomorphism such that

$$
V_{1}^{*} \otimes V_{2}^{*} \otimes \cdots \otimes V_{m}^{*} \cong L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)
$$

Proof: Define a map $\Phi: V_{1}^{*} \times V_{2}^{*} \times \cdots \times V_{m}^{*} \rightarrow L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)$ by

$$
\Phi\left(\omega^{1}, \omega^{2}, \ldots, \omega^{m}\right)\left(v_{1}, v_{2}, \ldots, v_{m}\right):=\omega^{1}\left(v_{1}\right) \omega^{2}\left(v_{2}\right) \cdots \omega^{m}\left(v_{m}\right)
$$

Because each $\omega^{i}$ is a linear form the expression on the right is indeed a multilinear function in the arguments $v_{1}, v_{2}, \ldots, v_{m}$. Further, $\Phi$ is a multilinear function in the arguments $\omega^{1}, \omega^{2}, \ldots, \omega^{m}$ because

$$
\begin{aligned}
\Phi\left(\omega^{1}, \omega^{2}, \ldots, a \omega^{i}+\widetilde{\omega}^{i}, \ldots, \omega^{m}\right)\left(v_{1}, v_{2}, \ldots, v_{m}\right)= & \omega^{1}\left(v_{1}\right) \omega^{2}\left(v_{2}\right) \cdots\left(a \omega^{i}+\widetilde{\omega}^{i}\right)\left(v_{i}\right) \cdots \omega^{m}\left(v_{m}\right) \\
= & \omega^{1}\left(v_{1}\right) \omega^{2}\left(v_{2}\right) \cdots\left(a \omega^{i}\left(v_{i}\right)+\widetilde{\omega}^{i}\left(v_{i}\right)\right) \cdots \omega^{m}\left(v_{m}\right) \\
= & a \omega^{1}\left(v_{1}\right) \omega^{2}\left(v_{2}\right) \cdots \omega^{i}\left(v_{i}\right) \cdots \omega^{m}\left(v_{m}\right) \\
& +\omega^{1}\left(v_{1}\right) \omega^{2}\left(v_{2}\right) \cdots \widetilde{\omega}^{i}\left(v_{i}\right) \cdots \omega^{m}\left(v_{m}\right) \\
= & a \Phi\left(\omega^{1}, \omega^{2}, \ldots, \omega^{i}, \ldots, \omega^{m}\right)\left(v_{1}, v_{2}, \ldots, v_{m}\right) \\
& +\Phi\left(\omega^{1}, \omega^{2}, \ldots, \widetilde{\omega}^{i}, \ldots, \omega^{m}\right)\left(v_{1}, v_{2}, \ldots, v_{m}\right)
\end{aligned}
$$

By the universal property of the tensor product (Proposition 1.1.1) $\Phi$ gives rise to a unique linear
$\operatorname{map} \widetilde{\Phi}: V_{1}^{*} \otimes V_{2}^{*} \otimes \cdots \otimes V_{m}^{*} \rightarrow L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)$ such that

$$
\widetilde{\Phi}\left(\omega^{1} \otimes \omega^{2} \otimes \cdots \otimes \omega^{m}\right)\left(v_{1}, v_{2}, \ldots, v_{m}\right)=\omega^{1}\left(v_{1}\right) \omega^{2}\left(v_{2}\right) \cdots \omega^{m}\left(v_{m}\right)
$$

From this property we can conclude that $\widetilde{\Phi}$ takes tensor products to concrete tensor products of linear forms as defined in (1.1.1). In addition, this implies that $\widetilde{\Phi}$ takes the basis of $V_{1}^{*} \otimes V_{2}^{*} \otimes \cdots \otimes V_{m}^{*}$ given by Proposition 1.1.3 to the basis of $L\left(V_{1}, V_{2}, \ldots, V_{m} ; \mathbb{R}\right)$ given by Lemma 1.1.4, so it is an isomorphism.

Through this identification, the elements of $V_{1}^{*} \otimes V_{2}^{*} \otimes \cdots \otimes V_{m}^{*}$ can either be regarded as elements of the abstract tensor product space or concretely as multilinear functions - whichever is more suitable. Furthermore, each vector space $V_{i}$ can be canonically identified with its bidual space $V_{i}^{* *}$ by the isomorphism $\Phi: V \rightarrow V^{* *}, \Phi(v)(f):=f(v)$. If we fix a basis $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ for $V$ and denote its dual basis by $\left\{e_{1}^{*}, e_{2}^{*}, \ldots, e_{n}^{*}\right\}$ and its bidual basis (the dual of the dual basis) by $\left\{e_{1}^{* *}, e_{2}^{* *}, \ldots, e_{n}^{* *}\right\}$ then it can be seen that $\Phi$ is an isomorphism. The mapping $\Phi$ sends a basis vector $e_{i}$ of $V$ to a basis vector $e_{i}^{* *}$ of $V^{* *}$ because

$$
\Phi\left(e_{i}\right)\left(e_{j}^{*}\right)=e_{j}^{*}\left(e_{i}\right)=\delta_{i}^{j}=e_{i}^{* *}\left(e_{j}^{*}\right)
$$

Hence, $\Phi$ is bijective. Therefore, we obtain another canonical identification for the spaces

$$
V_{1} \otimes V_{2} \otimes \cdots \otimes V_{m} \cong V_{1}^{* *} \otimes V_{2}^{* *} \otimes \cdots \otimes V_{m}^{* *} \cong L\left(V_{1}^{*}, V_{2}^{*}, \ldots, V_{m}^{*}, \mathbb{R}\right)
$$

### 1.1.2 THE EXTERIOR POWER

Henceforth, we only consider a single $n$-dimensional vector space $V$. For an integer $1 \leq m \leq n$, let us denote the $m$-fold tensor power $V^{\otimes m}:=V \otimes \cdots \otimes V$ by $T^{m}(V)$. An element $\alpha \in T^{m}(V)$ is called an $\boldsymbol{m}$-tensor of $\boldsymbol{V}$ of rank $\boldsymbol{m}$. We want to focus on elements of $T^{m}(V)$ that are alternating. An $m$-tensor $v_{1} \otimes v_{2} \otimes \cdots \otimes v_{m}$ of $V$ is called alternating if $v_{i}=v_{j}$ for some $i \neq j$. Denote the subspace spanned by such alternating tensors by $U^{\prime}$. Then we define the $\boldsymbol{m}$-th exterior power of $\boldsymbol{V}$ to be $\bigwedge^{m}(V):=T^{m}(V) / U^{\prime}$. The equivalence class of an element $v_{1} \otimes v_{2} \otimes \cdots \otimes v_{m}$ in $\Lambda^{m}(V)$ will be denoted

$$
v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}:=v_{1} \otimes v_{2} \otimes \cdots \otimes v_{m}+U^{\prime}
$$

The elements of $\bigwedge^{m}(V)$ and $\bigwedge^{m}\left(V^{*}\right)$ are called $\boldsymbol{m}$-vectors and (exterior) m-forms respectively.

Recall that we can think of both $m$-vectors and $m$-forms as multilinear functions on $V^{m}$ and $\left(V^{*}\right)^{m}$ respectively. Note that these multilinear functions are alternating (that is, they evaluate to zero on repeated arguments) by the definition of the exterior power. From this point, we will almost exclusively present properties of $m$-vectors. Similar results hold for $m$-forms by exchanging the roles of $V$ and $V^{*}$ and by using the isomorphism between $V$ and $V^{* *}$ wherever necessary.

An $m$-vector is said to be simple if it can be expressed in the form $\sigma=v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}$, where $v_{1}, v_{2}, \ldots v_{m} \in V$ are linearly independent. The subset of $\Lambda^{m}(V)$ of all simple $m$-vectors is called the Grassmannian cone $G C_{m}(V)$ - the reason for this nomenclature will be explained later. It is important to notice that not every $m$-vector is simple (as we will see shortly) and that the representation as $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}$ is, in general, not unique.

The significance of the simple $m$-vectors is that they provide a basis for $\bigwedge^{m} V$. Therefore, every non-simple $m$-vectors can be written uniquely as a linear combination of simple ones.

Lemma 1.1.6 ([Lee13, Proposition 14.8, p. 353])
Suppose $\left\{v_{1}, v_{2}, \ldots, v_{n}\right\}$ is a basis for $V$. Then for each $1 \leq m \leq n$ the set

$$
\mathcal{L}:=\left\{v_{i_{1}} \wedge v_{i_{2}} \wedge \cdots \wedge v_{i_{m}} \mid 1 \leq i_{1}<i_{2}<\ldots<i_{m} \leq n\right\}
$$

constitutes a basis for the exterior power $\bigwedge^{m}(V)$. Therefore, its dimension is

$$
\operatorname{dim}\left(\bigwedge^{m} V\right)=\binom{n}{m}
$$

Let $\left\{v_{1}, v_{2}, \ldots, v_{n}\right\}$ be a basis for $V$. Define the alternation mapping Alt: $T^{m} V \rightarrow \bigwedge^{m}(V)$ on a basis $m$-tensor as the signed average over all permutations of the components of this tensor, that is,

$$
\operatorname{Alt}\left(v_{i_{1}} \otimes v_{i_{2}} \otimes \cdots \otimes v_{i_{m}}\right):=\frac{1}{m!} \sum_{\sigma \in S_{m}} \operatorname{sgn}(\sigma) v_{\sigma\left(i_{1}\right)} \wedge v_{\sigma\left(i_{2}\right)} \wedge \cdots \wedge v_{\sigma\left(i_{m}\right)}
$$

Given an $m$-vector $v \in \bigwedge^{m}(V)$ and an $l$-vector $w \in \bigwedge^{l}(V)$ where $1 \leq m+l \leq n$, we define their wedge product to be

$$
v \wedge w:=\frac{(m+l)!}{m!l!} \operatorname{Alt}(v \otimes w)
$$

The wedge product is a mapping from $\bigwedge^{m}(V) \times \bigwedge^{l}(V)$ to $\bigwedge^{m+l}(V)$.

Lemma 1.1.7 (Basic Properties of the Exterior Power)
(i) Bilinearity: For $a \in \mathbb{R}, v \in \bigwedge^{m}(V)$ and $w, w^{\prime} \in \bigwedge^{l}(V)$

$$
\begin{aligned}
& \left(a w+w^{\prime}\right) \wedge v=a(w \wedge v)+\left(w^{\prime} \wedge v\right) \\
& v \wedge\left(a w+w^{\prime}\right)=a(v \wedge w)+\left(v \wedge w^{\prime}\right)
\end{aligned}
$$

(ii) Associativity: For $v \in \bigwedge^{k}(V), w \in \bigwedge^{l}(V)$ and $u \in \Lambda^{m}(V)$

$$
(v \wedge w) \wedge u=v \wedge(w \wedge u)
$$

(iii) Anticommutativity: Let $v_{1}, v_{2}, \ldots, v_{m} \in V$ and $\sigma \in S_{m}$ be a permutation on the integers $\{1,2, \ldots, m\}$. Then

$$
v_{\sigma(1)} \wedge v_{\sigma(2)} \wedge \cdots \wedge v_{\sigma(m)}=\operatorname{sgn}(\sigma) v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}
$$

Furthermore, for $v \in \bigwedge^{m}(V)$ and $w \in \bigwedge^{l}(V)$ and a simple $m$-vector $u \in G C_{m}(V)$

$$
\begin{aligned}
& v \wedge w=(-1)^{m l} w \wedge v \\
& u \wedge u=0
\end{aligned}
$$

(iv) The ordered m-tuple $\left(v_{1}, v_{2}, \ldots, v_{m}\right)$ is linearly independent in $V$ if and only if

$$
v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m} \neq 0
$$

(v) Suppose $W$ is a linear subspace of $V$ and $\mathcal{B}_{1}:=\left\{w_{1}, w_{2}, \ldots, w_{m}\right\}$ and $\mathcal{B}_{2}:=\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ are two bases of $W$. Then

$$
v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}=\lambda w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}
$$

for $\lambda=\operatorname{det}\left(T_{i}^{j}\right)_{i, j=1, \ldots, m} \neq 0$ where $T: W \rightarrow W$ is the linear map sending $\mathcal{B}_{1}$ to $\mathcal{B}_{2}$.
Proof: We omit the proofs of the parts (i) - (iii) (see [Lee13, Proposition 14.11, p. 356]). To prove part (iv) ([Lee13, Problem 14-4 (a), p. 376]) suppose $v_{1}, v_{2}, \ldots, v_{m}$ are linearly independent. We can extend $\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ to a basis $\left\{v_{1}, v_{2}, \ldots, v_{n}\right\}$ of $V$. The top exterior power $\bigwedge^{n}(V)$ is
one-dimensional by means of Lemma 1.1.6 and is spanned by the non-zero $n$-vector

$$
0 \neq v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}=\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right) \wedge\left(v_{m+1} \wedge v_{m+2} \wedge \cdots \wedge v_{n}\right)
$$

so $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m} \neq 0$. Conversely, suppose that $v_{1}, v_{2}, \ldots, v_{m}$ are linearly dependent and assume without loss of generality that $v_{1}=a^{2} v_{2}+\ldots+a^{m} v_{m}$. Then

$$
v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}=\left(a^{2} v_{2}+\ldots+a^{m} v_{m}\right) \wedge\left(v_{2} \wedge \cdots \wedge v_{m}\right)=0
$$

where the last equality is due to parts (i) and (iii).
For part (v) ([Lee13, Problem 14-4 (b), p. 376]) consider the linear map $T: W \rightarrow W$ given by $T w_{i}:=v_{i}$ for $i=1, \ldots, m$. In slight abuse of notation we identify the linear map $T$ with its matrix representation. Since both $m$-tuples span the same vector space we can express $v_{i}$ as a linear combination of the $w_{j}$, that is, $v_{i}=T w_{i}=\sum_{j=1}^{m} T_{i}^{j} w_{j}$ where $T_{i}^{j} \in \mathbb{R}$. Then by using part (i) we find that

$$
\begin{aligned}
v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m} & =T w_{1} \wedge T w_{2} \wedge \cdots \wedge T w_{m} \\
& =\left(\sum_{j_{1}=1}^{m} T_{1}^{j_{1}} w_{j_{1}}\right) \wedge\left(\sum_{j_{2}=1}^{m} T_{2}^{j_{2}} w_{j_{2}}\right) \wedge \cdots \wedge\left(\sum_{j_{m}=1}^{m} T_{m}^{j_{m}} w_{j_{m}}\right) \\
& =\sum_{j_{1}, j_{2}, \ldots, j_{m}=1}^{m}\left(\prod_{l=1}^{m} T_{l}^{j_{l}}\right) w_{j_{1}} \wedge w_{j_{2}} \wedge \cdots \wedge w_{j_{m}}
\end{aligned}
$$

By definition of the exterior power, any $m$-vector is alternating. Hence, any of the summands where any two of the indices $j_{l}$ coincide can be discarded. Thus, only summands where $\left(j_{1}, j_{2}, \ldots, j_{m}\right)$ are a permutation of $(1,2, \ldots, m)$ add up. By this reasoning and part (iii) we can write

$$
\begin{aligned}
v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m} & =\sum_{\sigma \in S_{m}}\left(\prod_{l=1}^{m} T_{l}^{\sigma(l)}\right) w_{\sigma(1)} \wedge w_{\sigma(2)} \wedge \cdots \wedge w_{\sigma(m)} \\
& =\sum_{\sigma \in S_{m}} \operatorname{sgn}(\sigma)\left(\prod_{l=1}^{m} T_{l}^{\sigma(l)}\right) w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}
\end{aligned}
$$

The last coefficient is the determinant of the matrix of the map $T$. This map is a bijection since $W$ is spanned by both $\mathcal{B}_{1}$ and $\mathcal{B}_{2}$. Thus, the determinant is non-zero.

Now we are equipped to show that the set $G C_{m}(V)$ of simple $m$-vectors is indeed a proper subset of $\bigwedge^{m}(V)$ as advertised earlier.

## Example:

Let $\sigma=e_{1} \wedge e_{2}+e_{3} \wedge e_{4} \in \bigwedge^{2}\left(\mathbb{R}^{4}\right)$, where $e_{i} \in \mathbb{R}^{4}$ are the standard basis vectors. Then $\sigma$ is not a simple 2 -vector. To show this, suppose $\sigma=v \wedge w$ for linearly independent $v, w \in \mathbb{R}^{4}$. Then $v=\sum_{i=1}^{4} a^{i} e_{i}$ and $w=\sum_{j=1}^{4} b^{j} e_{j}$. Using the basic properties (i) - (iii) in Lemma 1.1.7 we find

$$
\begin{aligned}
e_{1} \wedge e_{2}+e_{3} \wedge e_{4}=\sigma= & v \wedge w=\left(\sum_{i=1}^{4} a^{i} e_{i}\right) \wedge\left(\sum_{j=1}^{4} b^{j} e_{j}\right)=\sum_{i=1}^{4} \sum_{j=1}^{4} a^{i} b^{j} e_{i} \wedge e_{j} \\
= & a^{1} b^{1} e_{1} \wedge e_{1}+a^{1} b^{2} e_{1} \wedge e_{2}+a^{1} b^{3} e_{1} \wedge e_{3}+a^{1} b^{4} e_{1} \wedge e_{4} \\
& +a^{2} b^{1} e_{2} \wedge e_{1}+a^{2} b^{2} e_{2} \wedge e_{2}+a^{2} b^{3} e_{2} \wedge e_{3}+a^{2} b^{4} e_{2} \wedge e_{4} \\
& +a^{3} b^{1} e_{3} \wedge e_{1}+a^{3} b^{2} e_{3} \wedge e_{2}+a^{3} b^{3} e_{3} \wedge e_{3}+a^{3} b^{4} e_{3} \wedge e_{4} \\
& +a^{4} b^{1} e_{4} \wedge e_{1}+a^{4} b^{2} e_{4} \wedge e_{2}+a^{4} b^{3} e_{4} \wedge e_{3}+a^{4} b^{4} e_{4} \wedge e_{4} \\
= & \left(a^{1} b^{2}-a^{2} b^{1}\right) e_{1} \wedge e_{2}+\left(a^{1} b^{3}-a^{3} b^{1}\right) e_{1} \wedge e_{3}+\left(a^{1} b^{4}-a^{4} b^{1}\right) e_{1} \wedge e_{4} \\
& +\left(a^{2} b^{3}-a^{3} b^{2}\right) e_{2} \wedge e_{3}+\left(a^{2} b^{4}-a^{4} b^{2}\right) e_{2} \wedge e_{4}+\left(a^{3} b^{4}-a^{4} b^{3}\right) e_{3} \wedge e_{4} .
\end{aligned}
$$

By equating coefficients we find for example that

$$
\begin{align*}
& a^{1} b^{2}-a^{2} b^{1}=1 \\
& a^{3} b^{4}-a^{4} b^{3}=1 \\
& a^{1} b^{3}-a^{3} b^{1}=0  \tag{1.1.2}\\
& a^{2} b^{3}-a^{3} b^{2}=0 .
\end{align*}
$$

Assume $b^{3} \neq 0$ and substitute $a^{1}=\frac{a^{3}}{b^{3}} b^{1}$ into the first equation to get

$$
\frac{a^{3}}{b^{3}} b^{1} b^{2}-a^{2} b^{1}=1
$$

But substituting the last equation of (1.1.2) into the left hand side of the previous identity gives the contradiction

$$
\frac{a^{2}}{b^{3}} b^{1} b^{3}-a^{2} b^{1}=a^{2} b^{1}-a^{2} b^{1}=0 \neq 1
$$

Thus, $b^{3}=0$. But then we know from (1.1.2) that $a^{3} b^{2}=0=a^{3} b^{1}$. Suppose $a^{3}=0$. This contradicts the second equation of (1.1.2). So, we must have $b^{2}=0=b^{1}$ which contradicts the first equation of (1.1.2) and therefore $\sigma$ is not simple.

There is a useful characterisation of the exterior power of the dual space. Let us define a linear mapping $\Phi: \bigwedge^{m}\left(V^{*}\right) \rightarrow\left(\bigwedge^{m}(V)\right)^{*}$. First, define $\Phi$ on simple $m$-forms and $m$-vectors by

$$
\Phi\left(\omega^{i_{1}} \wedge \omega^{i_{2}} \wedge \cdots \wedge \omega^{i_{m}}\right)\left(v_{j_{1}} \wedge v_{j_{2}} \wedge \cdots \wedge v_{j_{m}}\right):=\operatorname{det}\left(\omega^{i_{l}}\left(v_{j_{k}}\right)\right)_{l, k=1, \ldots, m}
$$

Then extend the above definition of $\Phi$ bilinearly to all of $\bigwedge^{m}\left(V^{*}\right)$ and $\left(\bigwedge^{m}(V)\right)^{*}$.

## Lemma 1.1.8

The map $\Phi: \bigwedge^{m}\left(V^{*}\right) \rightarrow\left(\bigwedge^{m}(V)\right)^{*}$ is an isomorphism. Moreover, if $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ is a basis for $V$

$$
\Phi\left(e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*}\right)=\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right)^{*}
$$

where the asterisk denotes the dual basis vector in the corresponding dual space.
Proof: By definition, $\Phi$ is a linear map. Let $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ be a basis for $V$. Then $\left\{e_{1}^{*}, e_{2}^{*}, \ldots, e_{n}^{*}\right\}$ is the dual basis for $V^{*}$ (we break with our convention of upper indices for linear forms here) and the sets $\left\{e_{i_{1}} \wedge e_{i_{2}} \wedge \ldots \wedge e_{i_{m}} \mid 1 \leq i_{1}<i_{2}<\ldots<i_{m} \leq n\right\}$ and $\mathcal{B}:=\left\{e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*} \mid 1 \leq\right.$ $\left.i_{1}<i_{2}<\ldots<i_{m} \leq n\right\}$ are bases for $\bigwedge^{m}(V)$ and $\bigwedge^{m}\left(V^{*}\right)$, respectively. In addition, the set $\mathcal{C}:=\left\{\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right)^{*} \mid 1 \leq i_{1}<i_{2}<\ldots<i_{m} \leq n\right\}$ forms the dual basis of $\left(\bigwedge^{m}(V)\right)^{*}$. We will show that $\Phi$ maps the basis $\mathcal{B}$ to the basis $\mathcal{C}$ which proves that it is an isomorphism. Let $I=\left(i_{1}, i_{2}, \ldots, i_{m}\right)$ and $J=\left(j_{1}, j_{2}, \ldots, j_{m}\right)$ be two strictly increasing multi-indices. Calculate

$$
\Phi\left(e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*}\right)\left(e_{j_{1}} \wedge e_{j_{2}} \wedge \cdots \wedge e_{j_{m}}\right)=\operatorname{det}\left(e_{i_{l}}^{*}\left(e_{j_{k}}\right)\right)_{l, k=1, \ldots, m}=\operatorname{det}\left(\delta_{j_{k}}^{i_{l}}\right)_{l, k=1, \ldots, m}
$$

If $I=J$ then the right hand side is the determinant of the identity matrix and evaluates to 1 . Suppose $I \neq J$, then there is at least one $l \in\{1, \ldots, m\}$ such that $i_{l} \neq j_{k}$ for all $k \in\{1, \ldots, m\}$, that is, there is at least one row in the matrix $\left(\delta_{j_{k}}^{i_{l}}\right)_{l, k=1 \ldots m}$ which consists only of zero-valued entries. Therefore, the right hand side in the latter equation evaluates to zero. This proves that $\Phi\left(e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*}\right)=\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right)^{*}$ by definition of the dual basis.

Hereafter, we will make no distinction between these two spaces and consequently let $m$-forms $\omega \in \bigwedge^{m}\left(V^{*}\right)$ act on $m$-vectors $\sigma \in \bigwedge^{m}(V)$.

### 1.1.3 The Plücker embedding

Let us investigate the implications of Lemma 1.1.7 (v) in more detail. Consider an $N$-dimensional vector space $\mathcal{V}$ and define an equivalence relation $\sim$ on the set $\mathcal{V} \backslash\{0\}$ by proposing that two non-
zero elements $v, w \in \mathcal{V}$ are equivalent if and only if there is a non-zero number $\lambda \in \mathbb{R}$ such that $v=\lambda w$. The equivalence class $[v]_{\sim}$ is the one-dimensional subspace spanned by $v \in \mathcal{V}$. The quotient $\mathbb{P}(\mathcal{V}):=(\mathcal{V} \backslash\{0\}) / \sim$ is called projective space. Geometrically speaking, we can view the projective space $\mathbb{P}(\mathcal{V})$ as the set of all lines in $\mathcal{V}$ that pass through the origin, that is, the set of all one-dimensional subspaces of the $N$-dimensional vector space $\mathcal{V}$.

Generalising the preceding idea, let us define the set $G_{m}(\mathcal{V})$ of all $m$-dimensional subspaces of an $N$-dimensional vector space $\mathcal{V}$ for each $1 \leq m \leq N$. This set is called the Grassmannian of $\mathcal{V}$.

We introduce a mapping $\rho: G_{m}(V) \rightarrow \mathbb{P}\left(\bigwedge^{m}(V)\right)$. Let $W \in G_{m}(V)$ be an $m$-dimensional subspace of $V$ and $\left\{w_{1}, w_{2}, \ldots, w_{m}\right\}$ a basis of $W$. Define $\rho(W):=\left[w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right]_{\sim}$. By means of Lemma 1.1.7 (v) this mapping is well-defined because different choices of bases for $W$ yield the same image in the projective space $\mathbb{P}\left(\bigwedge^{m}(V)\right)$. The mapping $\rho$ is known as the Plücker embedding.

Proposition 1.1.9 ([Har95, p. 64])
The Plücker embedding $\rho: G_{m}(V) \rightarrow \mathbb{P}\left(\bigwedge^{m}(V)\right)$ is injective.
Proof: This proof is due to [Hud07, Proposition 2.3, p. 3]. Define $\varphi: \mathbb{P}\left(\bigwedge^{m}(V)\right) \rightarrow \bigcup_{s=1}^{n} G_{s}(V)$ by

$$
\varphi\left([w]_{\sim}\right):=\left\{v \in V \mid v \wedge w=0 \in \bigwedge^{m+1}(V)\right\}
$$

The set $\varphi\left([w]_{\sim}\right)$ is a subspace of $V$ due to the bilinearity of the wedge product (Lemma 1.1.7 (i)). Thus, the map $\varphi$ is well-defined. We will show that $\varphi$ is a left-inverse of $\rho$, that is, $\varphi \circ \rho=\operatorname{Id}_{G_{m}(V)}$. Let $W \in G_{m}(V)$ be arbitrary with basis $\left\{w_{1}, w_{2}, \ldots, w_{m}\right\}$, such that $\rho(W)=\left[w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right]_{\sim}$. For each $w \in W$, it is clear by Lemma 1.1.7 (iv) that $w \wedge w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}=0$ hence, $W \subset \varphi \circ \rho(W)$. Conversely, if $v \in \varphi \circ \rho(W)$ then $v \wedge w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}=0$. Extend $\left(w_{1}, w_{2}, \ldots, w_{m}\right)$ to a basis $\left\{w_{1}, w_{2}, \ldots, w_{n}\right\}$ of $V$ and write $v=\sum_{i=1}^{n} a^{i} w_{i}$. The anticommutativity of the wedge product (Lemma 1.1.7 (iii)) implies

$$
\begin{aligned}
0 & =\left(\sum_{i=1}^{n} a^{i} w_{i}\right) \wedge w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m} \\
& =(-1)^{m} \sum_{i=m+1}^{n} a^{i} w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m} \wedge w_{i}
\end{aligned}
$$

and because $w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m} \wedge w_{i}$ for $i=m+1, \ldots, n$ are linearly independent $(m+1)$-vectors by means of Lemma 1.1.6, all the coefficients $a^{i}$ for $i=m+1, \ldots, n$ vanish. Thus, $v=\sum_{i=1}^{m} a^{i} w_{i} \in W$ and $\varphi \circ \rho(W) \subset W$. This proves that $\rho$ is injective.

Many more results are known about the Plücker embedding. For example, the Plücker embedding
is, in fact, a topological embedding, that is, a homeomorphism onto its image and hence deserves its name. For more details, we refer to the survey article [BN91]. The Plücker embedding provides an important relationship between $m$-dimensional subspaces and simple $m$-vectors.

Lemma 1.1.10 ([Hud07, Lemma 3.8, p. 7])
An element $[w]_{\sim} \in \mathbb{P}\left(\bigwedge^{m}(V)\right)$ lies in the image of the Grassmannian under the Plücker embedding if and only if $w$ is a simple m-vector. That is, $[w]_{\sim} \in \operatorname{Im}(\rho)$ if and only if $w \in G C_{m}(V)$.
Proof: If $w$ can be written as $w=v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}$ for linearly independent vectors $v_{1}, v_{2}, \ldots, v_{m} \in V$ then the subspace of $V$ spanned by $\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ is $m$-dimensional, hence there is a $U \in G_{m}(V)$ with $U=\operatorname{span}\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ and $\rho(U)=[w]_{\sim}$. Conversely, suppose $[w]_{\sim}=\rho(U)$ for some $U \in G_{m}(V)$. Choose a basis $\left\{u_{1}, u_{2}, \ldots, u_{m}\right\}$ for $U$. Then by definition of the Plücker embedding $[w]_{\sim}=\left[u_{1} \wedge u_{2} \wedge \cdots \wedge u_{m}\right]_{\sim}$. Thereby, $w$ can be written as $\lambda u_{1} \wedge u_{2} \wedge \cdots \wedge u_{m}$ for some non-zero scalar $\lambda \in \mathbb{R}$, and so $\omega$ is a simple $m$-vector.

Generally, an algebraic cone is a subset $C$ of a vector space such that $\lambda c \in C$ whenever $c \in C$ and $\lambda>0$. The Plücker embedding provides us with a tool to treat the Grassmannian as a subset of the projectivisation of simple $m$-vectors. This is the reason why the set of simple vectors is called the Grassmannian cone $G C_{m}(V)$. Thompson [Tho96, p. 196] attributes this nomenclature to Busemann, Ewald and Shepard [BES63].

### 1.1.4 InNER PRODUCTS, VOLUME FORMS AND A RIESZ-TYPE ISOMORPHISM

Let us conclude the section on multilinear algebra by imposing an additional structure. Suppose $\left(V,\langle\cdot, \cdot\rangle_{V}\right)$ is an inner product space, which is a finite-dimensional vector space with a positive-definite, symmetric, bilinear mapping $\langle\cdot, \cdot\rangle_{V}: V \times V \rightarrow \mathbb{R}$, called an inner product. Notice that the inner product induces a norm

$$
\|\cdot\|_{V}: V \rightarrow \mathbb{R}_{+}, v \mapsto\langle v, v\rangle_{V}^{\frac{1}{2}}
$$

on $V$. The dual norm

$$
\|\omega\|_{V^{*}}:=\sup _{\|u\|_{V}=1}|\omega(u)|
$$

of a linear form $\omega \in V^{*}$ turns the dual space $V^{*}$ into a normed space. By the Riesz representation theorem [Wer00, Theorem V.3.6, p. 228] the map $J: V \rightarrow V^{*}$ given by $J(v)(u):=\langle v, u\rangle_{V}$ for $v, u \in V$ is an isomorphism. The dual space $V^{*}$ turns into an inner product space by defining
$\langle\omega, \nu\rangle_{V^{*}}:=\left\langle J^{-1}(\omega), J^{-1}(\nu)\right\rangle_{V}$. Since $J^{-1}$ is linear and $\langle\cdot, \cdot\rangle_{V}$ is bilinear, so is $\langle\cdot, \cdot\rangle_{V^{*}}$. The symmetry of this mapping follows immediately from the symmetry of $\langle\cdot, \cdot\rangle_{V}$. Because $J^{-1}$ is an isomorphism and the inner product on $V$ is positive-definite, the positive-definiteness of $\langle\cdot, \cdot\rangle_{V^{*}}$ follows.

With respect to this inner product on the dual space, $J$ is an isometry, which means it is an isomorphism that preserves the inner product. In fact, the dual norm as defined above is induced by this inner product because $J$ is an isometry. This can be seen by calculating that for some $\omega \in V^{*}$ and the corresponding $v \in V$ with $J(v)=\omega$ one obtains

$$
\|\omega\|_{V^{*}}=\sup _{\|u\|_{V}=1}|\omega(u)|=\sup _{\|u\|_{V}=1}|J(v)(u)|=\sup _{\|u\|_{V}=1}\left|\langle v, u\rangle_{V}\right| \leq \sup _{\|u\|_{V}=1}\|v\|_{V}\|u\|_{V}=\|v\|_{V}
$$

Further,

$$
\left|w\left(\frac{v}{\|v\|_{V}}\right)\right|=\left|\left\langle v, \frac{v}{\|v\|_{V}}\right\rangle_{V}\right|=\|v\|_{V}
$$

and thus,

$$
\|\omega\|_{V^{*}}^{2}=\|v\|_{V}^{2}=\langle v, v\rangle_{V}=\left\langle J^{-1}(\omega), J^{-1}(\omega)\right\rangle_{V}=\langle\omega, \omega\rangle_{V^{*}}
$$

Moreover, we can endow the exterior power $\bigwedge^{m}(V)$ with an inner product. Define it on simple $m$-vectors by

$$
\left\langle v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}, w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right\rangle_{\wedge^{m}(V)}:=\operatorname{det}\left(\left\langle v_{i}, w_{j}\right\rangle_{V}\right)_{i, j=1 \ldots m}
$$

where $v_{i}, w_{i} \in V$ for $i=1, \ldots, m$ and extend bilinearly (see e.g. [Tho96, p. 192]). The symmetry then follows from the properties of the determinant and the inner product on $V$. The positive-definiteness is established as follows. The matrix $G:=\left(\left\langle v_{i}, v_{j}\right\rangle_{V}\right)_{i, j=1 \ldots m} \in \mathbb{R}^{m \times m}$ - the so-called Gram matrix of $\langle\cdot, \cdot\rangle_{V}$ - is positive semi-definite because for $0 \neq \alpha \in \mathbb{R}^{m}$

$$
\begin{equation*}
\alpha^{t} G \alpha=\sum_{i=1}^{m} \sum_{j=1}^{m} \alpha_{i}\left\langle v_{i}, v_{j}\right\rangle_{V} \alpha_{j}=\left\langle\sum_{i=1}^{m} \alpha_{i} v_{i}, \sum_{j=1}^{m} \alpha_{j} v_{j}\right\rangle_{V}=\left\|\sum_{i=1}^{m} \alpha_{i} v_{i}\right\|_{V}^{2} \geq 0 \tag{1.1.3}
\end{equation*}
$$

where we simply renamed the index $j$ in the last equation. Suppose $\alpha \neq 0$ is an eigenvector of $G$ and $\lambda \in \mathbb{R}$ its eigenvalue. Then $\lambda=\left\|\sum_{i=1}^{m} \alpha_{i} v_{i}\right\|_{V}^{2} /\|\alpha\|_{\mathbb{R}^{m}}^{2} \geq 0$ by 1.1.3. Since the determinant is the product of the eigenvalues, we have $\operatorname{det}(G) \geq 0$. Further, $\operatorname{det}(G)=0$ if and only if at least one eigenvalue equals zero. This means $\operatorname{det}(G)=0$ if and only if there is $0 \neq \alpha \in \mathbb{R}^{m}$ such that $\sum_{i=1}^{m} \alpha_{i} v_{i}=0$. Equivalently, $\operatorname{det}(G)=0$ if and only if the set $\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ is linearly dependent.

Due to Lemma 1.1.7 (iv) this means $\operatorname{det}(G)=0$ if and only if $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}=0$. Hence, $\langle\cdot, \cdot\rangle_{\wedge^{m}(V)}$ does indeed define an inner product.

Again the map $\mathcal{J}: \bigwedge^{m}(V) \rightarrow\left(\bigwedge^{m}(V)\right)^{*}$ given by $\mathcal{J}(\sigma)(\tau):=\langle\sigma, \tau\rangle_{\bigwedge^{m}(V)}$ for $\sigma, \tau \in \bigwedge^{m}(V)$ is an isomorphism which is isometric with respect to the induced dual inner product on $\left(\bigwedge^{m}(V)\right)^{*}$

Recall from Lemma 1.1.8 that $\bigwedge^{m}\left(V^{*}\right)$ and $\left(\bigwedge^{m}(V)\right)^{*}$ are isomorphic through $\Phi$. Since now there are inner products on each of these spaces at our disposal, the question arises whether $\Phi$ is also an isomorphism of inner product spaces not only vector spaces, that is, whether $\Phi$ preserves the inner product.

Let $\mathcal{I}:=\left\{\left(i_{1}, i_{2}, \ldots, i_{m}\right) \mid 1 \leq i_{1}<i_{2}<\ldots<i_{m} \leq n\right\}$ and $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ be a basis for $V$. Further, let $F=\sum_{I \in \mathcal{I}} a^{I} e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*}$ and $G=\sum_{I \in \mathcal{I}} b^{I} e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*}$ be two elements of $\wedge^{m}\left(V^{*}\right)$. Recall that $\Phi\left(e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*}\right)=\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right)^{*}$ by Lemma 1.1.8. It follows from the definition of the induced inner product on the dual space that

$$
\begin{aligned}
\langle\Phi(F), \Phi(G)\rangle_{\left(\wedge^{m}(V)\right)^{*}} & =\sum_{I, J \in \mathcal{I}} a^{I} b^{J}\left\langle\Phi\left(e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*}\right), \Phi\left(e_{j_{1}}^{*} \wedge e_{j_{2}}^{*} \wedge \cdots \wedge e_{j_{m}}^{*}\right)\right\rangle_{\left(\wedge^{m}(V)\right)^{*}} \\
& =\sum_{I, J \in \mathcal{I}} a^{I} b^{J}\left\langle\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right)^{*},\left(e_{j_{1}} \wedge e_{j_{2}} \wedge \cdots \wedge e_{j_{m}}\right)^{*}\right\rangle_{\left(\wedge^{m}(V)\right)^{*}} \\
& =\sum_{I, J \in \mathcal{I}} a^{I} b^{J}\left\langle e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}, e_{j_{1}} \wedge e_{j_{2}} \wedge \cdots \wedge e_{j_{m}}\right\rangle_{\wedge^{m}(V)} \\
& =\sum_{I, J \in \mathcal{I}} a^{I} b^{J} \operatorname{det}\left(\left\langle e_{i_{l}}, e_{j_{k}}\right\rangle_{V}\right)_{l, k=1, \ldots, m} \\
& =\sum_{I, J \in \mathcal{I}} a^{I} b^{J} \operatorname{det}\left(\left\langle e_{i_{l}}^{*}, e_{j_{k}}^{*}\right\rangle_{V^{*}}\right)_{l, k=1, \ldots, m} \\
& =\sum_{I, J \in \mathcal{I}} a^{I} b^{J}\left\langle e_{i_{1}}^{*} \wedge e_{i_{2}}^{*} \wedge \cdots \wedge e_{i_{m}}^{*}, e_{j_{1}}^{*} \wedge e_{j_{2}}^{*} \wedge \cdots \wedge e_{j_{m}}^{*}\right\rangle_{\wedge^{m}\left(V^{*}\right)} \\
& =\langle F, G\rangle_{\wedge^{m}\left(V^{*}\right)}
\end{aligned}
$$

Thus, we have proved the following result.

## Lemma 1.1.11

The spaces $\left(\bigwedge^{m}\left(V^{*}\right),\langle\cdot, \cdot\rangle_{\bigwedge^{m}\left(V^{*}\right)}\right)$ and $\left(\left(\bigwedge^{m}(V)\right)^{*},\langle\cdot, \cdot\rangle_{\left(\bigwedge^{m}(V)\right)^{*}}\right)$ are isometrically isomorphic by means of the mapping $\Phi: \bigwedge^{m}\left(V^{*}\right) \rightarrow\left(\bigwedge^{m}(V)\right)^{*}$ in Lemma 1.1.8.

Similarly to Lemma 1.1.8 we will not distinguish these two spaces and use either of the inner products and norms - whichever is more convenient.

Two vectors $v, v^{\prime} \in \mathcal{V}$ of an inner product space $\mathcal{V}$ are said to be orthogonal if $\left\langle v, v^{\prime}\right\rangle_{\mathcal{V}}=0$ and orthonormal if additionally $\langle v, v\rangle_{\mathcal{V}}=1=\left\langle v^{\prime}, v^{\prime}\right\rangle_{\mathcal{V}}$. Let $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ be an orthonormal basis for
$V$. Then for two strictly increasing multi-indices $I$ and $J$

$$
\left\langle e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}, e_{j_{1}} \wedge e_{j_{2}} \wedge \cdots \wedge e_{j_{m}}\right\rangle_{\wedge^{m}(V)}=\operatorname{det}\left(\left\langle e_{i_{l}}, e_{j_{k}}\right\rangle_{V}\right)_{l, k=1 \ldots m}=\operatorname{det}\left(\delta_{j_{k}}^{i_{l}}\right)_{l, k=1 \ldots m}
$$

By the same argumentation used in the proof of Lemma 1.1.8 the right hand side is 0 if $I \neq J$ and 1 if $I=J$. The preceding calculation proves the next result.

## Corollary 1.1.12

An orthonormal basis $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ for $V$ yields an orthonormal basis $\left\{e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}} \mid 1 \leq\right.$ $\left.i_{1}<i_{2}<\ldots<i_{m} \leq n\right\}$ for the exterior power $\bigwedge^{m}(V)$.

## Corollary 1.1.13

The spaces $\left(\bigwedge^{m}\left(\mathbb{R}^{n}\right),\|\cdot\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)}\right)$ and $\left(\mathbb{R}^{N},\|\cdot\|_{\mathbb{R}^{N}}\right)$, where $N:=\binom{n}{m}$ and $\|\cdot\|_{\mathbb{R}^{N}}$ is the standard Euclidean norm on $\mathbb{R}^{N}$, are isometrically isomorphic.
Proof: Let $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ be the standard orthonormal basis for $\mathbb{R}^{n}$. Due to the preceding corollary the set $\mathcal{E}:=\left\{e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}} \mid 1 \leq i_{1}<i_{2}<\ldots<i_{m} \leq n\right\}$ is an orthonormal basis for $\bigwedge^{m}\left(\mathbb{R}^{n}\right)$. Consider the set $\mathcal{I}:=\left\{\left(i_{1}, i_{2}, \ldots, i_{m}\right) \mid 1 \leq i_{1}<i_{2}<\ldots<i_{m} \leq n\right\}$ of strictly increasing multi-indices. Let $\mathfrak{E}:=\left\{\widetilde{e}_{I} \mid I \in \mathcal{I}\right\}$ denote the standard orthonormal basis for $\mathbb{R}^{N}$. Since both spaces are of the same dimension $\binom{n}{m}$, let $\Phi: \bigwedge^{m}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}^{N}$ be the isomorphism that sends the basis $\mathcal{E}$ to the basis $\mathfrak{E}$. Then for $\sigma=\sum_{I \in \mathcal{I}} a^{I} e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}} \in \bigwedge^{m}\left(\mathbb{R}^{n}\right)$ we have that

$$
\|\sigma\|_{\bigwedge^{m}\left(\mathbb{R}^{n}\right)}^{2}=\left\|\sum_{I \in \mathcal{I}} a^{I} e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)}^{2}=\sum_{I \in \mathcal{I}}\left(a^{I}\right)^{2}
$$

because of the orthonormality of $\mathcal{E}$. On the other hand

$$
\|\Phi(\sigma)\|_{\mathbb{R}^{N}}^{2}=\left\|\sum_{I \in \mathcal{I}} a^{I} \Phi\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right)\right\|_{\mathbb{R}^{N}}^{2}=\left\|\sum_{I \in \mathcal{I}} a^{I} \widetilde{e}_{I}\right\|_{\mathbb{R}^{N}}^{2}=\sum_{I \in \mathcal{I}}\left(a^{I}\right)^{2}
$$

which proves the claim.

## Lemma 1.1.14

For any $\sigma=w_{1} \wedge w_{2} \wedge \cdots w_{m} \in G C_{m}(V)$ it holds that

$$
\begin{equation*}
\|\sigma\|_{\Lambda^{m}(V)} \leq \prod_{i=1}^{m}\left\|w_{i}\right\|_{V} \tag{1.1.4}
\end{equation*}
$$

with equality if and only if $w_{1}, w_{2}, \ldots, w_{m}$ are pairwise orthogonal.

Proof: We will take advantage of the alternating property of the wedge product. By means of the Gram-Schmidt process define the vectors

$$
\begin{array}{ll}
\widetilde{w}_{1}:=w_{1} & v_{1}:=\frac{\widetilde{w}_{1}}{\left\|\widetilde{w}_{1}\right\|_{V}} \\
\widetilde{w}_{i}:=w_{i}-\sum_{j=1}^{i-1}\left\langle w_{i}, v_{j}\right\rangle_{V} v_{j} & v_{i}:=\frac{\widetilde{w}_{i}}{\left\|\widetilde{w}_{i}\right\|_{V}},
\end{array} \quad i=2,3, \ldots m .
$$

Note that the vectors $\left\{\widetilde{w}_{1}, \widetilde{w}_{2}, \ldots, \widetilde{w}_{m}\right\}$ and $\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ are pairwise orthogonal respectively. By definition of the vectors $\widetilde{w}_{i}$ and $v_{i}$ and the multilinearity of the wedge product we see that

$$
\begin{aligned}
\widetilde{w}_{1} \wedge \widetilde{w}_{2} \wedge \cdots \wedge \widetilde{w}_{m} & =\widetilde{w}_{1} \wedge \widetilde{w}_{2} \wedge \cdots \wedge \widetilde{w}_{m-1} \wedge\left(w_{m}-\sum_{j=1}^{m-1}\left\langle w_{i}, v_{j}\right\rangle_{V} v_{j}\right) \\
& =\widetilde{w}_{1} \wedge \widetilde{w}_{2} \wedge \cdots \wedge \widetilde{w}_{m-1} \wedge w_{m}-\sum_{j=1}^{m-1}\left\langle w_{i}, v_{j}\right\rangle_{V} \widetilde{w}_{1} \wedge \widetilde{w}_{2} \wedge \cdots \wedge \widetilde{w}_{m-1} \wedge \frac{\widetilde{w}_{j}}{\left\|\widetilde{w}_{j}\right\|_{V}} \\
& =\widetilde{w}_{1} \wedge \widetilde{w}_{2} \wedge \cdots \wedge \widetilde{w}_{m-1} \wedge w_{m}
\end{aligned}
$$

where we used Lemma 1.1.7 (v) in the last step. Proceeding inductively, we have shown that

$$
\widetilde{w}_{1} \wedge \widetilde{w}_{2} \wedge \cdots \wedge \widetilde{w}_{m}=w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}
$$

Thus, it follows that

$$
w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}=\widetilde{w}_{1} \wedge \widetilde{w}_{2} \wedge \cdots \wedge \widetilde{w}_{m}=\left(\prod_{i=1}^{m}\left\|\widetilde{w}_{i}\right\|_{V}\right) v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}
$$

Use the definition of the induced norm to get

$$
\begin{align*}
\left\|\widetilde{w}_{i}\right\|_{V}^{2} & =\left\langle\widetilde{w}_{i}, \widetilde{w}_{i}\right\rangle_{V}=\left\langle w_{i}-\sum_{j=1}^{i-1}\left\langle w_{i}, v_{j}\right\rangle_{V} v_{j}, w_{i}-\sum_{k=1}^{i-1}\left\langle w_{i}, v_{k}\right\rangle_{V} v_{k}\right\rangle_{V} \\
& =\left\langle w_{i}, w_{i}\right\rangle_{V}-2 \sum_{j=1}^{i-1}\left\langle w_{i}, v_{j}\right\rangle_{V}^{2}+\sum_{j=1}^{i-1} \sum_{k=1}^{i-1}\left\langle w_{i}, v_{j}\right\rangle_{V}\left\langle w_{i}, v_{k}\right\rangle_{V} \underbrace{\left\langle v_{j}, v_{k}\right\rangle_{V}}_{=\delta_{j k}}  \tag{1.1.5}\\
& =\left\langle w_{i}, w_{i}\right\rangle_{V}-\sum_{j=1}^{i-1}\left\langle w_{i}, v_{j}\right\rangle_{V}^{2} \leq\left\langle w_{i}, w_{i}\right\rangle_{V}=\left\|w_{i}\right\|_{V}^{2}
\end{align*}
$$

Subsequently, we find that

$$
\begin{aligned}
\|\sigma\|_{\Lambda^{m}(V)} & =\left\|w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right\|_{\Lambda^{m}(V)} \\
& =\left(\prod_{i=1}^{m}\left\|\widetilde{w}_{i}\right\|_{V}\right)\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right\|_{\wedge^{m}(V)} \\
& \leq\left(\prod_{i=1}^{m}\left\|w_{i}\right\|_{V}\right)\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right\|_{\wedge^{m}(V)}=\left(\prod_{i=1}^{m}\left\|w_{i}\right\|_{V}\right)
\end{aligned}
$$

The last inequality follows from Corollary 1.1 .12 since $\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ is an orthonormal set of vectors. Assume equality holds in (1.1.4). We will show that $w_{i}=\widetilde{w}_{i}$ with $\widetilde{w}_{i}$ defined as above; for then $w_{1}, w_{2}, \ldots, w_{m}$ are pairwise orthogonal by the Gram-Schmidt process. By definition of $\widetilde{w}_{i}$,

$$
w_{i}=\widetilde{w}_{i}+\sum_{j=1}^{i-1}\left\langle w_{i}, v_{j}\right\rangle_{V} v_{j} .
$$

Equality in (1.1.4) implies equality in (1.1.5), so that $\sum_{j=1}^{i-1}\left\langle w_{i}, v_{j}\right\rangle_{V}^{2}=0$ for all $i=1,2, \ldots, m$. Therefore, for each $i=1,2, \ldots, m$ we know that $\left\langle w_{i}, v_{j}\right\rangle_{V}=0$ for $j=1,2, \ldots, i-1$. Hence, $w_{i}=\widetilde{w}_{i}$.

Conversely, if $w_{1}, w_{2}, \ldots, w_{m}$ are pairwise orthogonal then by definition of the norm on the exterior power

$$
\begin{aligned}
\left\|w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right\|_{\wedge^{m}(V)}^{2} & =\operatorname{det}\left(\left\langle w_{i}, w_{j}\right\rangle_{V}\right)_{i, j=1 \ldots m}=\operatorname{det}\left(\left\langle w_{i}, w_{i}\right\rangle_{V} \delta_{i j}\right)_{i, j=1 \ldots m} \\
& =\prod_{i=1}^{m}\left\|w_{i}\right\|_{V}^{2}=\left(\prod_{i=1}^{m}\left\|w_{i}\right\|_{V}\right)^{2}
\end{aligned}
$$

Recall that the dimension of $\bigwedge^{m}(V)$ is the binomial coefficient $\binom{n}{m}$. Therefore, the top exterior power $\bigwedge^{n}(V)$ is one-dimensional, meaning it is just a real line. A non-zero $n$-form $\omega \in \bigwedge^{n}\left(V^{*}\right) \cong\left(\bigwedge^{n}(V)\right)^{*}$ is called a volume (or orientation) form for $V$. In the two-dimensional case we will call an element of $\bigwedge^{2}\left(V^{*}\right)$ an area form for $V$. A volume form provides the one-dimensional real line $\bigwedge^{n}(V)$ with an orientation in a natural way by proposing that $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}$ is positively oriented if

$$
\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)>0
$$

Two volume forms $\omega$ and $\omega^{\prime}$ establish the same orientation if and only if the factor by which they differ is positive. In particular, given a basis $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ for $V$, denote its dual basis $\left\{\varepsilon^{1}, \varepsilon^{2}, \ldots, \varepsilon^{n}\right\}$. Then the $n$-form $\omega_{e}=\varepsilon^{1} \wedge \varepsilon^{2} \wedge \cdots \wedge \varepsilon^{n}$ is a volume form such that $\omega_{e}\left(e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right)=1$. It is called the standard volume form. Since $\bigwedge^{n}\left(V^{*}\right)$ is one-dimensional, every other volume form $\omega$ is
a real multiple of the standard volume form $\omega_{e}$. Conversely, given an arbitrary volume form $\omega$ one can choose a basis $\left\{e_{i}^{\prime}\right\}_{i=1}^{m}$ for $V$ such that $\omega\left(e_{1}^{\prime} \wedge e_{2}^{\prime} \wedge \cdots \wedge e_{n}^{\prime}\right)=1$.

Let $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\} \subset V$ and $\left\{v_{1}, v_{2}, \ldots, v_{n}\right\} \subset V$ be two bases of $V$. Then the $n$-vectors $v_{1} \wedge v_{2} \wedge$ $\cdots \wedge v_{n}$ and $e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}$ are consistently oriented with respect to a volume form $\omega \in \wedge^{n}\left(V^{*}\right)$ if

$$
\operatorname{sgn}\left(\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)\right)=\operatorname{sgn}\left(\omega\left(e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right)\right)
$$

Note that this can always be achieved by switching $v_{1}$ and $v_{2}$ if necessary. Lemma 1.1.7 (v) and the properties of the norm on $\bigwedge^{n}(V)$ imply

$$
\frac{\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right\|_{\wedge^{n}(V)}}{\left\|e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right\|_{\wedge^{n}(V)}}=|\operatorname{det}(T)|=\frac{\left|\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)\right|}{\left|\omega\left(e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right)\right|}
$$

wherein $T: V \rightarrow V$ is the linear map that sends $e_{i}$ to $v_{i}$. If $e_{1} \wedge e_{2} \wedge \ldots \wedge e_{n}$ is a positively oriented unit $n$-vector then

$$
\begin{equation*}
\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right\|_{\wedge^{n}(V)}=\frac{\left|\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)\right|}{\omega\left(e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right)} \tag{1.1.6}
\end{equation*}
$$

and if, in addition, $\omega=\omega_{e}$ is the standard volume form corresponding to $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ then

$$
\begin{equation*}
\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right\|_{\wedge^{n}(V)}=\left|\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)\right| \tag{1.1.7}
\end{equation*}
$$

In Euclidean space the following result gives meaning to the name "volume form".

Proposition 1.1.15 ([For09, (5.3) Beispiel, p. 48])
Let $v_{1}, v_{2}, \ldots, v_{n} \in \mathbb{R}^{n}$ be linearly independent vectors and $P$ the parallelotope spanned by these vectors, that is, $P:=\left\{\sum_{i=1}^{n} t_{i} v_{i} \mid t_{i} \in[0,1]\right\}$. Then $\mathcal{L}^{n}(P)=\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right\|_{\wedge^{n}\left(\mathbb{R}^{n}\right)}$ where $\mathcal{L}^{n}$ is the $n$-dimensional Lebesgue measure on $\mathbb{R}^{n}$.

Proof: The standard basis $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\} \subset \mathbb{R}^{n}$ is orthonormal with respect to the standard Euclidean scalar product. Hence, $\left\|e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right\|_{\wedge^{n}\left(\mathbb{R}^{n}\right)}=1$ by Corollary 1.1.12. Consider the linear map $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ that sends $e_{i}$ to $v_{i}$. Then $P$ is the image of the unit cube $[0,1]^{n} \subset \mathbb{R}^{n}$ under this mapping. Recall that the $n$-dimensional Lebesgue measure of the unit cube is 1 . Further, the differential of a linear map is a linear map which has the same matrix representation. Now observe that by the change of variables theorem (see e.g. [For09, §13, Satz 2, p. 120]) the $n$-dimensional

Lebesgue measure of $P$ equals

$$
\begin{aligned}
\mathcal{L}^{n}(P) & =\int_{T\left([0,1]^{n}\right)} 1 \mathrm{~d} \mathcal{L}^{n}=\int_{[0,1]^{n}}|\operatorname{det} d T| \mathrm{d} \mathcal{L}^{n}=|\operatorname{det} T| \mathcal{L}^{n}\left([0,1]^{n}\right) \\
& =|\operatorname{det} T|=\frac{\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right\|_{\wedge^{n}\left(\mathbb{R}^{n}\right)}}{\left\|e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right\|_{\bigwedge^{n}\left(\mathbb{R}^{n}\right)}}=\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right\|_{\wedge^{n}\left(\mathbb{R}^{n}\right)}
\end{aligned}
$$

Hence, by using (1.1.6) and Proposition 1.1.15, we find that the volume form applied to some $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}$ gives the (signed) $n$-dimensional Lebesgue measure of the parallelotope spanned by the vectors $v_{1}, v_{2}, \ldots, v_{n}$ with respect to some normalisation. This normalisation corresponds to the assignment of a value for the volume of the "standard parallelotope". This geometric interpretation justifies the given terminology.

The following considerations are needed in the application of multilinear algebra in Chapter 2. Given a volume form $\omega \in \bigwedge^{n}\left(V^{*}\right)$, we want to to show that there is a Riesz-type isomorphism between $V$ and its dual space $V^{*}$ by using the volume form (instead of the inner product). Less vaguely, let us define a map $\tilde{\iota}_{\omega}: \bigwedge^{n-1}(V) \rightarrow V^{*}$ by

$$
\begin{equation*}
\tilde{\iota}_{\omega}(\sigma)(v):=\omega(\sigma \wedge v) \tag{1.1.8}
\end{equation*}
$$

where $\sigma \in \bigwedge^{n-1}(V), v \in V$.

## Proposition 1.1.16

The map $\tilde{\iota}_{\omega}: \bigwedge^{n-1}(V) \rightarrow V^{*}$ defined by (1.1.8) is an isomorphism.
Proof: First, note that both $\bigwedge^{n-1}(V)$ and $V^{*}$ are $n$-dimensional. Clearly $\tilde{l}_{\omega}$ is linear because for arbitrary but fixed $v \in V$ we know that

$$
\tilde{\iota}_{\omega}\left(a \sigma+\sigma^{\prime}\right)(v)=\omega\left(\left(a \sigma+\sigma^{\prime}\right) \wedge v\right)=a \omega(\sigma \wedge v)+\omega\left(\sigma^{\prime} \wedge v\right)=a \tilde{\iota}_{\omega}(\sigma)(v)+\tilde{\iota}_{\omega}\left(\sigma^{\prime}\right)(v)
$$

where we used the bilinearity of the wedge product (Lemma 1.1.7 (i)) and the linearity of a volume form. Let $\left\{v_{1}, v_{2}, \ldots v_{n}\right\}$ be a basis for $V$. Any $\sigma \in \Lambda^{n-1}(V)$ can be written as a linear combination of basis $(n-1)$-vectors $\sigma=\sum_{j=1}^{n} \alpha^{j} v_{1} \wedge v_{2} \wedge \cdots \wedge \widehat{v}_{j} \wedge \cdots \wedge v_{n}$ where the hat indicates that $v_{j}$ is to be omitted. Let $\sigma \in \Lambda^{n-1}(V)$ be such that $\omega(\sigma \wedge v)=0$ for all $v \in V$. Note that if $i \neq j$ then $v_{1} \wedge v_{2} \wedge \cdots \wedge \widehat{v}_{j} \wedge \cdots \wedge v_{n} \wedge v_{i}=0$. Thus, we have the following $n$ equations

$$
\begin{aligned}
0=\omega\left(\sigma \wedge v_{i}\right) & =\sum_{j=1}^{n} \alpha^{j} \omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge \widehat{v}_{j} \wedge \cdots \wedge v_{n} \wedge v_{i}\right) \\
& =\alpha^{i} \omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge \widehat{v}_{i} \wedge \cdots \wedge v_{n} \wedge v_{i}\right)
\end{aligned}
$$

$$
=(-1)^{n-i} \alpha^{i} \omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)
$$

for $i=1,2, \ldots, n$. Because $\omega$ is not the zero form and $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}$ spans $\bigwedge^{n}(V)$, we find that $\alpha^{i}=0$ for all $i=1,2, \ldots, n$, that is, $\sigma=0$. Hence, the kernel of $\tilde{\iota}_{\omega}$ is trivial and $\tilde{\iota}_{\omega}$ is injective.

Let $l \in V^{*}$ be arbitrary. For $j=1,2, \ldots, n$ define the coefficients

$$
\alpha^{j}:=(-1)^{n-j} \frac{l\left(v_{j}\right)}{\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)}
$$

and set $\sigma:=\sum_{j=1}^{n} \alpha^{j} v_{1} \wedge v_{2} \wedge \cdots \wedge \widehat{v}_{j} \wedge \cdots \wedge v_{n}$. Then $\tilde{\iota}_{\omega}(\sigma)=l$ since for arbitrary $v=\sum_{i=1}^{n} \beta^{i} v_{i} \in V$ it holds that $\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge \widehat{v}_{j} \wedge \cdots \wedge v_{n} \wedge v\right)=(-1)^{n-j} \beta^{j} \omega\left(v_{1} \wedge v_{2} \cdots \wedge v_{n}\right)$ and thus

$$
\begin{aligned}
\tilde{\iota}_{\omega}(\sigma)(v)=\omega(\sigma \wedge v) & =\sum_{j=1}^{n} \alpha^{j} \omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge \widehat{v}_{j} \wedge \cdots \wedge v_{n} \wedge v\right) \\
& =\sum_{j=1}^{n} \frac{(-1)^{n-j} l\left(v_{j}\right)}{\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)} \omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge \widehat{v}_{j} \wedge \cdots \wedge v_{n} \wedge v\right) \\
& =\frac{1}{\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)} \sum_{j=1}^{n} l\left(v_{j}\right) \beta^{j} \omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right) \\
& =\sum_{j=1}^{n} l\left(v_{j}\right) \beta^{j}=l\left(\sum_{j=1}^{n} \beta^{j} v_{j}\right)=l(v),
\end{aligned}
$$

that is, $\tilde{\iota}_{\omega}$ is also surjective and thereby an isomorphism.

Since $\operatorname{dim}\left(\bigwedge^{m}(V)\right)=\binom{n}{m}=\binom{n}{n-m}=\operatorname{dim}\left(\bigwedge^{n-m}(V)\right)$ we know these two spaces are isomorphic as vector spaces. This holds true generally for same-dimensional vector spaces, but the isomorphism is not necessarily canonical (that is, basis-independent). In this particular case however, the so-called Hodge dual $\star: \bigwedge^{m}(V) \rightarrow \bigwedge^{n-m}(V)$ indeed provides a canonical isomorphism between $\bigwedge^{m}(V)$ and $\bigwedge^{n-m}(V)$. Let us only state that the inner product structure and a fixed orientation is needed to show that the Hodge dual is an isomorphism. It is completely determined by the equation

$$
\begin{equation*}
(\star \lambda) \wedge \theta=\langle\lambda, \theta\rangle_{\wedge^{m}(V)} \varepsilon \tag{1.1.9}
\end{equation*}
$$

where $\lambda, \theta \in \bigwedge^{m}(V)$ are arbitrary and $\varepsilon \in \bigwedge^{n}(V)$ with $\|\varepsilon\|_{\Lambda^{n}(V)}=1$ is the positively oriented unit $n$-vector (see [BG68, §2.22, pp. 108-110] or [Fla63, §2.7, p. 15]).

Using the Hodge dual for $m=1$ together with Proposition 1.1.16 yields the isomorphism

$$
\begin{equation*}
\iota_{\omega}=\tilde{\iota}_{\omega} \circ \star: V \rightarrow V^{*} . \tag{1.1.10}
\end{equation*}
$$

Note that we only needed an inner product for the Hodge dual - Proposition 1.1.16 holds more generally in finite-dimensional normed vector spaces $V$. In particular, if the dimension $\operatorname{dim}(V)=n=2$ then the isomorphism $\tilde{\iota}_{\omega}$ suffices for our needs and no inner product on $V$ is required.

Let us define an $n$-vector $\omega^{*}:=\left(\iota_{\omega}\right)^{*} \omega \in \bigwedge^{n}(V) \cong \bigwedge^{n}\left(\left(V^{*}\right)^{*}\right)$ as the pullback of the volume form $\omega$ by $\iota_{\omega}$, that is,

$$
\begin{equation*}
\omega^{*}\left(\varphi^{1} \wedge \varphi^{2} \wedge \cdots \wedge \varphi^{n}\right):=\omega\left(\left(\iota_{\omega}\right)^{-1}\left(\varphi^{1}\right) \wedge\left(\iota_{\omega}\right)^{-1}\left(\varphi^{2}\right) \wedge \cdots \wedge\left(\iota_{\omega}\right)^{-1}\left(\varphi^{n}\right)\right) . \tag{1.1.11}
\end{equation*}
$$

for $\varphi^{1} \wedge \varphi^{2} \wedge \cdots \wedge \varphi^{n} \in G C_{n}\left(V^{*}\right)$ and extend linearly to $\wedge^{n}\left(V^{*}\right)$. This is well-defined because $\omega$ itself is an $n$-form and thus alternating. In addition, $\omega^{*}$ is non-zero because $\iota_{\omega}$ is an isomorphism. Therefore, $\omega^{*}$ is a volume form on the dual space $V^{*}$. It will be called the dual volume form with respect to $\omega$.

## Lemma 1.1.17

Let $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ be an orthonormal basis for $V$ and $\omega=\omega_{e}$ the standard volume form. Denote the dual basis on $V^{*}$ by $\left\{\varepsilon^{1}, \varepsilon^{2}, \ldots, \varepsilon^{n}\right\}$. Then $\iota_{\omega}\left(e_{i}\right)=\varepsilon^{i}$ for all $i=1,2, \ldots, n$ and

$$
\begin{aligned}
\omega^{*}\left(\varepsilon^{1} \wedge \varepsilon^{2} \wedge \cdots \wedge \varepsilon^{n}\right) & =1 \\
\left\|\varepsilon^{1} \wedge \varepsilon^{2} \wedge \cdots \wedge \varepsilon^{n}\right\|_{\wedge^{n}\left(V^{*}\right)} & =1
\end{aligned}
$$

Proof: The element $e_{1} \wedge e_{2} \wedge \ldots \wedge e_{n}$ is the positively oriented unit $n$-vector of $\wedge^{n}(V)$ since its components form an orthonormal basis of $V$. The first assertion then follows from (1.1.9) since

$$
\begin{aligned}
\iota_{\omega}\left(e_{i}\right)\left(e_{j}\right) & =\left(\tilde{\iota}_{\omega} \circ \star\right)\left(e_{i}\right)\left(e_{j}\right)=\tilde{\iota}_{\omega}\left(\star e_{i}\right)\left(e_{j}\right) \\
& =\omega\left(\star e_{i} \wedge e_{j}\right) \\
& =\omega\left(\left\langle e_{i}, e_{j}\right\rangle_{V} e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right) \\
& =\delta_{j}^{i} \omega\left(e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right)=\delta_{j}^{i}
\end{aligned}
$$

for all $i, j=1,2, \ldots, n$. Then

$$
\begin{aligned}
\omega^{*}\left(\varepsilon^{1} \wedge \varepsilon^{2} \wedge \cdots \wedge \varepsilon^{n}\right) & =\omega\left(\left(\iota_{\omega}\right)^{-1}\left(\varepsilon^{1}\right) \wedge\left(\iota_{\omega}\right)^{-1}\left(\varepsilon^{2}\right) \wedge \cdots \wedge\left(\iota_{\omega}\right)^{-1}\left(\varepsilon^{n}\right)\right) \\
& =\omega\left(e_{1} \wedge e_{2} \wedge \cdots \wedge e_{n}\right)=1
\end{aligned}
$$

If $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ is an orthonormal basis for $V$ then the dual basis $\left\{\varepsilon^{1}, \varepsilon^{2}, \ldots, \varepsilon^{n}\right\}$ is orthonormal in $V^{*}$ by virtue of the Riesz representation theorem. The definition of the norm on $\bigwedge^{n}\left(V^{*}\right)$ implies the
final statement.

## Proposition 1.1.18 (Isometric property of $\boldsymbol{\iota}_{\boldsymbol{\omega}}$ )

Let $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ be an orthonormal basis for $V$ and $\omega=\omega_{e}$ the standard volume form. The isomorphism $\iota_{\omega}$ given by (1.1.10) is isometric in the following sense:

$$
\left\|\iota_{\omega}\left(v_{1}\right) \wedge \iota_{\omega}\left(v_{2}\right) \wedge \cdots \wedge \iota_{\omega}\left(v_{n}\right)\right\|_{\wedge^{n}\left(V^{*}\right)}=\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right\|_{\wedge^{n}(V)}
$$

for all $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n} \in \bigwedge^{n}(V)$.
Proof: Recall that $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n} \neq 0 \in \bigwedge^{n}(V)$ if and only if $v_{1}, v_{2}, \ldots, v_{n} \in V$ are linearly independent. So, because $\iota_{\omega}$ is an isomorphism, the set $\left\{\iota_{\omega}\left(v_{1}\right), \iota_{\omega}\left(v_{2}\right), \ldots, \iota_{\omega}\left(v_{n}\right)\right\}$ is a basis for $V^{*}$ if $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n} \neq 0 \in \bigwedge^{n}(V)$. Let $T: V^{*} \rightarrow V^{*}$ be the map that sends $\varepsilon^{i}$ to $\iota_{\omega}\left(v_{i}\right)$ where $\varepsilon^{i}$ is the dual basis vector to $e_{i}$. Lemma 1.1.7 (v) and Lemma 1.1.17 yield

$$
\begin{aligned}
\left\|\iota_{\omega}\left(v_{1}\right) \wedge \iota_{\omega}\left(v_{2}\right) \wedge \cdots \wedge \iota_{\omega}\left(v_{n}\right)\right\|_{\wedge^{n}\left(V^{*}\right)} & =\frac{\left\|\iota_{\omega}\left(v_{1}\right) \wedge \iota_{\omega}\left(v_{2}\right) \wedge \cdots \wedge \iota_{\omega}\left(v_{n}\right)\right\|_{\wedge^{n}\left(V^{*}\right)}}{\left\|\varepsilon^{1} \wedge \varepsilon^{2} \wedge \cdots \wedge \varepsilon^{n}\right\|_{\wedge^{n}\left(V^{*}\right)}} \\
& =|\operatorname{det} T| \\
& =\left|\frac{\omega^{*}\left(\iota_{\omega}\left(v_{1}\right) \wedge \iota_{\omega}\left(v_{2}\right) \wedge \cdots \wedge \iota_{\omega}\left(v_{n}\right)\right)}{\omega^{*}\left(\varepsilon^{1} \wedge \varepsilon^{2} \wedge \cdots \wedge \varepsilon^{n}\right)}\right| \\
& =\left|\omega^{*}\left(\iota_{\omega}\left(v_{1}\right) \wedge \iota_{\omega}\left(v_{2}\right) \wedge \cdots \wedge \iota_{\omega}\left(v_{n}\right)\right)\right| \\
& =\left|\omega\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right)\right| \\
& =\left\|v_{1} \wedge v_{2} \wedge \cdots \wedge v_{n}\right\|_{\wedge^{n}(V)} .
\end{aligned}
$$

The last two equalities follow from the definition of $\omega^{*}$ and (1.1.7).

We state a technical result on the norms of $n$-vectors which is needed later. For the next lemma, suppose $V$ is an $n$-dimensional inner product space. Further, let $\omega \in \bigwedge^{n}\left(V^{*}\right)$ be a volume form and $e \in \bigwedge^{n}(V)$ the unit $n$-vector which is positively oriented with respect to the orientation induced by $\omega$ (that is, $\|e\|_{\Lambda^{n}(V)}=1$ and $\omega(e)>0$ ). Notice that the index set $\mathcal{I}$ in the following lemma is slightly different to our usage of it so far.

## Lemma 1.1.19

Let $\left\{v_{1}, v_{2}, \ldots, v_{N}\right\} \subset V$ be a collection of vectors for arbitrary $N \in \mathbb{N}$. Suppose that for any strictly increasing multi-index $I \in \mathcal{I}:=\left\{\left(i_{1}, i_{2}, \ldots, i_{n}\right) \mid 1 \leq i_{1}<i_{2}<\ldots<i_{n} \leq N\right\}$ of length $n$ the $n$-vectors $v_{i_{1}} \wedge v_{i_{2}} \wedge \cdots \wedge v_{i_{n}}$ are pairwise consistently oriented. Let $m \in\{1,2, \ldots, n\}$ and fix a strictly increasing
multi-index $J=\left(j_{m+1}, j_{m+1}, \ldots, j_{n}\right)$ of length $n-m$ such that $m \leq j_{m+1}$ and $j_{n} \leq N$. Then

$$
\left\|\sum_{\substack{I \in \mathcal{I} \\ i_{m}<j_{m+1}}} v_{i_{1}} \wedge \cdots \wedge v_{i_{m}} \wedge v_{j_{m+1}} \wedge \cdots \wedge v_{j_{n}}\right\|_{\wedge^{n}(V)}=\sum_{\substack{I \in \mathcal{I} \\ i_{m}<j_{m+1}}}\left\|v_{i_{1}} \wedge \cdots \wedge v_{i_{m}} \wedge v_{j_{m+1}} \wedge \cdots \wedge v_{j_{n}}\right\|_{\wedge^{n}(V)} .
$$

In particular,

$$
\begin{equation*}
\left\|\sum_{I \in \mathcal{I}} v_{i_{1}} \wedge v_{i_{2}} \wedge \cdots \wedge v_{i_{n}}\right\|_{\wedge^{n}(V)}=\sum_{I \in \mathcal{I}}\left\|v_{i_{1}} \wedge v_{i_{2}} \wedge \cdots \wedge v_{i_{n}}\right\|_{\wedge^{n}(V)} \tag{1.1.12}
\end{equation*}
$$

Proof: The second assertion follows from the first for $m=n$. If $N<n$ then both sides of the equation are zero (because $\mathcal{I}=\varnothing$ ). Therefore, suppose $N \geq n$. Since $\bigwedge^{n}(V)$ is one-dimensional we can write $v_{i_{1}} \wedge v_{i_{2}} \wedge \cdots \wedge v_{i_{n}}=c_{i_{1} i_{2} \ldots i_{n}} e$ for some coefficents $c_{i_{1} i_{2} \ldots i_{n}} \in \mathbb{R}$. All of the coefficients $c_{i_{1} i_{2} \ldots i_{n}}=\omega\left(v_{i_{1}} \wedge v_{i_{2}} \wedge \cdots \wedge v_{i_{n}}\right) / \omega(e)$ are either positive or negative for all $I \in \mathcal{I}$ because the $n$ vectors $v_{i_{1}} \wedge v_{i_{2}} \wedge \cdots \wedge v_{i_{n}}$ are consistently oriented. In other words, $\operatorname{sgn}\left(c_{i_{1} i_{2} \ldots i_{n}}\right)$ is constant for any possible multi-index. Then

$$
\begin{aligned}
\sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}} \| v_{i_{1}} & \wedge \cdots \wedge v_{i_{m}} \wedge v_{j_{m+1}} \wedge \cdots \wedge v_{j_{n}} \|_{\wedge^{n}(V)} \\
& =\sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}}\left\|c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}} e\right\|_{\wedge^{n}(V)}=\sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}}\left|c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}}\right|\|e\|_{\wedge^{n}(V)} \\
& =\sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}}\left|c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}}\right|=\left|\sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}}\right| c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}}| | \\
& =|\sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}}^{\underbrace{\operatorname{sgn}\left(c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}}\right)}_{= \pm 1, \text { const }}\left|c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}}\right|\left|=\left|\sum_{i_{m \in \mathcal{I}}} c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}}\right|\right.}|
\end{aligned}
$$

and on the other hand

$$
\| \begin{aligned}
& \left\|\sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}} v_{i_{1}} \wedge \cdots \wedge v_{i_{m}} \wedge v_{j_{m+1}} \wedge \cdots \wedge v_{j_{n}}\right\|_{\wedge^{n}(V)} \\
& =\left\|\sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}} c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}} e\right\|_{\wedge^{n}(V)}=\mid \sum_{\substack{I \in \mathcal{I} \\
i_{m}<j_{m+1}}} c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}}\|e\|_{\Lambda^{n}(V)}
\end{aligned}
$$

$$
=\left|\sum_{\substack{I \in \mathcal{I} \\ i_{m}<j_{m+1}}} c_{i_{1} \ldots i_{m}, j_{m+1} \ldots j_{n}}\right|
$$

This section concludes with the special case $n=2$ for the previous result.

## Corollary 1.1.20

If $\operatorname{dim}(V)=n=2$ in the setting of Lemma 1.1.19 then the formulae simplify to

$$
\left\|\sum_{i=1}^{j-1} v_{i} \wedge v_{j}\right\|_{\wedge^{2}(V)}=\sum_{i=1}^{j-1}\left\|v_{i} \wedge v_{j}\right\|_{\Lambda^{2}(V)}
$$

for $j=2,3, \ldots, N$ and

$$
\left\|\sum_{1 \leq i<j \leq N} v_{i} \wedge v_{j}\right\|_{\Lambda^{2}(V)}=\sum_{1 \leq i<j \leq N}\left\|v_{i} \wedge v_{j}\right\|_{\Lambda^{2}(V)}
$$

### 1.2 Convex geometry

The scope of this section is to introduce basic concepts of convex geometry. In particular, we define polytopes - a special type of convex sets. Further, we describe the support function of a convex set and state some results for the special case of polytopes. Subsequently, the duality between the vertex and the half-space representation of polytopes is outlined. Further, we prove the Minkowski inequality for convex bodies and use a geometric argument to establish a well-known explicit formula for a specific mixed volume. Finally, we conclude this section with a result on the maxima of convex functions on polytopes. The results in this section are borrowed from the textbooks of Bonnesen and Fenchel [BF71], Brøndsted [Brø83], Grünbaum [Grü03], Matoušek [Mat02], Rockafellar [Roc70] and Thompson [Tho96, Chapter 0, Chapter 2]. Throughout the section ( $V,\|\cdot\|_{V}$ ) will denote an $n$-dimensional normed vector space.

### 1.2.1 Convex sets and polytopes

A non-empty subset $C$ of the vector space $V$ is said to be convex if $\lambda x+(1-\lambda) y \in C$ whenever $x, y \in C$ and $\lambda \in[0,1]$.

The segment joining $\boldsymbol{x}$ and $\boldsymbol{y}$ is defined as the set $\{\lambda x+(1-\lambda) y \mid \lambda \in[0,1]\}$ and denoted as $[x, y]$. Thus, as a more geometrical characterisation of the previous definition, a set $C$ is convex if and only if it contains all segments joining any two of its points.

Suppose $A$ and $B$ are convex sets. The Minkowski sum of $A$ and $B$ is

$$
A+B:=\{a+b \mid a \in A, b \in B\} .
$$

For any non-negative real number $\alpha$, define further

$$
\alpha A:=\{\alpha a \mid a \in A\}
$$

If $A=\{a\}$ is a singleton then we usually write $a+B$ instead of $\{a\}+B$. Note that the decomposition of a vector $v \in A+B$ as a sum of vectors in $A$ and $B$ is, in general, not unique.

## Lemma 1.2.1

Suppose $A$ and $B$ are convex sets and $\alpha \geq 0$ is a real number. Then $\alpha A+B$ is a convex set.
Proof: For $i=1,2$ let $x_{i}=\alpha a_{i}+b_{i}$ where $a_{i} \in A$ and $b_{i} \in B$ and $t \in[0,1]$. Then

$$
t x_{1}+(1-t) x_{2}=\alpha\left(t a_{1}+(1-t) a_{2}\right)+\left(t b_{1}+(1-t) b_{2}\right)
$$

which is an element of $\alpha A+B$ because $A$ and $B$ are convex.

Using induction we can easily generalise the previous result and show that arbitrary finite linear combinations of convex sets are convex. Furthermore, the associativity and commutativity of the vector operations in $V$ translate to the respective properties for linear combinations of convex sets. We can also show the following distributive law.

## Lemma 1.2.2

Let $\alpha, \beta_{1}, \beta_{2} \geq 0$ be non-negative real numbers and suppose $A_{1}, A_{2}$ and $B \subset V$ are convex sets. Then

$$
\begin{aligned}
\alpha\left(A_{1}+A_{2}\right) & =\alpha A_{1}+\alpha A_{2} \\
\left(\beta_{1}+\beta_{2}\right) B & =\beta_{1} B+\beta_{2} B
\end{aligned}
$$

Proof: The first assertion follows from the distributive law on $V$ since $x \in \alpha\left(A_{1}+A_{2}\right)$ if and only if there are $a_{i} \in A_{i}$ such that $x=\alpha\left(a_{1}+a_{2}\right)=\alpha a_{1}+\alpha a_{2}$ which is an element of $\alpha A_{1}+\alpha A_{2}$. Similarly, if $y \in\left(\beta_{1}+\beta_{2}\right) B$ then $y=\left(\beta_{1}+\beta_{2}\right) b=\beta_{1} b+\beta_{2} b$ and thus, $y \in \beta_{1} B+\beta_{2} B$. Conversely, let $y=\beta_{1} b_{1}+\beta_{2} b_{2} \in \beta_{1} B+\beta_{2} B$. Because $\beta_{1}$ and $\beta_{2}$ are non-negative, the numbers $\frac{\beta_{1}}{\beta_{1}+\beta_{2}}$ and $\frac{\beta_{2}}{\beta_{1}+\beta_{2}}$ are non-negative and add to 1 . Then the convexity of $B$ implies $\frac{\beta_{1}}{\beta_{1}+\beta_{2}} b_{1}+\frac{\beta_{2}}{\beta_{1}+\beta_{2}} b_{2} \in B$ and thus

$$
y=\beta_{1} b_{1}+\beta_{2} b_{2}=\left(\beta_{1}+\beta_{2}\right)\left(\frac{\beta_{1}}{\beta_{1}+\beta_{2}} b_{1}+\frac{\beta_{2}}{\beta_{1}+\beta_{2}} b_{2}\right) \in\left(\beta_{1}+\beta_{2}\right) B
$$

The reason we defined $\alpha A$ only for non-negative scalars is precisely because of this distributive law which only applies to non-negative scalars. Also, for many of the sets that we will consider in Chapter 2, multiplication by a negative scalar will be irrelevant because they will be symmetric in the following sense.

A set $A$ is symmetric with respect to the origin $0 \in \boldsymbol{V}$ if $(-1) A=\{-a \mid a \in A\}=A$. Henceforth, if we use the adjective "symmetric" it always means symmetric with respect to the origin.

We can see as follows that any intersection of convex sets is convex. If $x$ and $y$ are elements of all the members of the intersection, so is the segment joining $x$ and $y$ because each of the sets is convex. But then the line segment already is contained in the intersection itself. Therefore it is fitting to consider the intersection of the family of convex sets containing a given (not necessarily convex) set $A$. By construction, this is the smallest convex set containing $A$. It is called the convex hull of $A$ and denoted by conv $(A)$. Due to Carathéodory [DGK63, pp.103+115] the convex hull of a set $A \subset V$ is
characterised by

$$
\begin{equation*}
\operatorname{conv}(A)=\left\{\sum_{i=1}^{n+1} \lambda_{i} a_{i} \mid a_{i} \in A, \lambda_{i} \geq 0, \sum_{i=1}^{n+1} \lambda_{i}=1\right\} \tag{1.2.13}
\end{equation*}
$$

One can show that the convex hull operation preserves inclusions, in other words conv $(X) \subset \operatorname{conv}(Y)$ if $X \subset Y$.

A linear subspace $W \subset V$ is trivially convex. Clearly, a singleton $\{x\}$ is convex as well. Thus, any translate of a linear subspace, that is, $x+W$ is convex. Such a set is an affine subspace of $V$. Its dimension is defined as the dimension of the subspace of which it is a translate. The dimension of a convex set is the dimension of the affine subspace of smallest dimension containing the set. For example, the dimension of a cube is 3 , the dimension of a triangle 2 and the dimension of a point is 0 . Since any convex set is contained in the linear space $V$, the dimension of a convex set is bounded above by the dimension $n$ of the ambient space $V$.

A convex subset $E$ of a convex set $A$ is said to be an extreme subset of $A$ if, whenever $x \in E$ and $x=\lambda a_{1}+(1-\lambda) a_{2}$ for $a_{1}, a_{2} \in A$ and $\lambda \in(0,1)$, then $a_{1}, a_{2} \in E$ already.

A special type of convex sets, the "polytopes", will occupy much of our focus in this section and the application later. To avoid unnecessary clutter in explicit calculations and formulae, we make the following construction. Suppose $\left\{a_{1}, a_{2}, \ldots, a_{N}\right\}$ and $\left\{b_{1}, b_{2}, \ldots, b_{M}\right\}$ are two non-empty sets of finitely many points in $V$. Any such two sets are in polytopial relation to one another if and only if $\operatorname{conv}\left(\left\{a_{1}, a_{2}, \ldots, a_{N}\right\}\right)=\operatorname{conv}\left(\left\{b_{1}, b_{2}, \ldots, b_{M}\right\}\right)$. Clearly, this defines an equivalence relation because equality of sets is one such correspondence. The equivalence class $\left[\left\{a_{1}, a_{2}, \ldots, a_{N}\right\}\right]_{\text {poly }}$ can be well-ordered when ordering the number of points of each of its elements with the usual order $\leq$ of the natural numbers $\mathbb{N}$. For each equivalence class $\left[\left\{a_{1}, a_{2}, \ldots, a_{N}\right\}\right]_{p o l y}$ the well-ordering principle then yields the existence of an element with a minimal number of points.

A polytope is the convex hull of the minimal element in the equivalence class of a non-empty set of finitely many points in $V$, that is,

$$
A_{N}:=\left[a_{1} a_{2} \ldots a_{N}\right]:=\operatorname{conv}\left(\left\{a_{1}, a_{2}, \ldots, a_{N}\right\}\right),
$$

where $N$ is minimal as described above. The points $a_{1}, a_{2}, \ldots, a_{N}$ are the vertices of $A$. In the special case of dimension $n=2$, a polytope is called a polygon. If there is no need for the number of indices the index $N$ will be omitted.

If not explicitly stated otherwise a polytope in $V$ is considered as full dimensional, that is, its dimension is $n$. This can always be achieved by restricting the point of view to the affine subspace containing
the polytope. In doing so, the number of vertices $N$ is bounded below by $n+1$. This can be seen through (1.2.13).

The extreme subsets of a polytope are called the faces or $\boldsymbol{j}$-faces if their dimension is $j$. The ( $n-1$ )-dimensional faces of a polytope are called facets. Note that the 0-dimensional faces of a polytope are its vertices. One can show that the faces of a polytope are polytopes themselves ([Brø83, Theorem 7.3, p. 45]).

## Lemma 1.2.3

A polytope is a convex, closed and bounded (and thus compact) set.
Proof: Let $A=\left[a_{1} a_{2} \ldots a_{N}\right]$ be a polytope. By definition, a polytope is convex as the convex hull of its vertices. Note that a polytope $A$ is contained in the ball of radius $\max _{i=1,2, \ldots, N}\left\|a_{i}\right\|_{V}$ and thus bounded. Let $\left(x_{k}\right)_{k \in \mathbb{N}} \subset A$ be a sequence which converges to an arbitrary but fixed $x \in V$ in the norm topology of $V$. We need to show that $x \in A$. Any of the points $x_{k}$ can be expressed as a convex combination of the vertices, that is, $x_{k}=\sum_{i=1}^{N} \lambda_{i}^{(k)} a_{i}$ where $\sum_{i=1}^{N} \lambda_{i}^{(k)}=1$ and $\lambda_{i}^{(k)} \in[0,1]$. Because [0,1] is a compact set there is a subsequence, say $\left(k_{i_{1}}\right)_{i_{1} \in \mathbb{N}} \subset(k)_{k \in \mathbb{N}}$, such that $\left(\lambda_{1}^{\left(k_{i_{1}}\right)}\right)_{i_{1} \in \mathbb{N}}$ converges to some $\lambda_{1} \in[0,1]$. But then there is another subsequence of this subsequence, say $\left(k_{i_{2}}\right)_{i_{2} \in \mathbb{N}} \subset\left(k_{i_{1}}\right)_{i_{1} \in \mathbb{N}} \subset(k)_{k \in \mathbb{N}}$, such that both $\left(\lambda_{1}^{\left(k_{i_{2}}\right)}\right)_{i_{1} \in \mathbb{N}}$ and $\left(\lambda_{2}^{\left(k_{i_{2}}\right)}\right)_{i_{1} \in \mathbb{N}}$ converge to some $\lambda_{1}$ and $\lambda_{2} \in[0,1]$ respectively. Iterating this process $N$ times and renaming the final subsequence yields the convergence of $\left(\lambda_{i}^{(k)}\right)_{k \in \mathbb{N}}$ to $\lambda_{i} \in[0,1]$ for $i=1,2, \ldots, N$. Further, $\sum_{i=1}^{N} \lambda_{i}=\sum_{i=1}^{N} \lim _{k \rightarrow \infty} \lambda_{i}^{(k)}=\lim _{k \rightarrow \infty} \sum_{i=1}^{N} \lambda_{i}^{(k)}=1$. Define $y:=\sum_{i=1}^{N} \lambda_{i} a_{i} \in A$. Then the triangle inequality for the norm on $V$ implies

$$
\left\|y-x_{k}\right\|_{V} \leq \sum_{i=1}^{N}\left|\lambda_{i}^{(k)}-\lambda_{i}\right|\left\|a_{i}\right\|_{V} \underset{k \rightarrow \infty}{ } 0
$$

and

$$
\|y-x\|_{V} \leq\left\|y-x_{k}\right\|_{V}+\left\|x_{k}-x\right\|_{V} \underset{k \rightarrow \infty}{ } 0
$$

Thus, $x=\lim _{k \rightarrow \infty} x_{k}=y \in A$.

### 1.2.2 Support Functions

Consider an $(n-1)$-dimensional subspace $W$ of $V$ with basis $\left\{w_{1}, w_{2}, \ldots, w_{n-1}\right\}$. Extend this to a basis $\left\{w_{1}, w_{2}, \ldots, w_{n-1}, v\right\}$ of the entire space $V$ and define a linear form $f: V \rightarrow \mathbb{R}$ on this basis by $f\left(w_{i}\right)=0$ for $i=1,2, \ldots, n-1$ and $f(v)=1$. Then $\operatorname{ker}(f)=W$. Clearly, this particular linear
form $f$ is not unique - a different choice of basis yields a different linear form. Conversely, if $f$ is a non-zero element of the dual space $V^{*}$ then its image is all of $\mathbb{R}$, because for a vector $v \notin \operatorname{ker}(f)$ it holds that $f(\alpha v)=\alpha f(v)$ for all $\alpha \in \mathbb{R}$. The rank-nullity-theorem then states that the kernel of $f$ is of dimension $\operatorname{dim}(\operatorname{ker}(f))=\operatorname{dim}(V)-\operatorname{dim}(\operatorname{im}(f))=n-1$. Therefore, it is appropriate to introduce the following notion.

A translate of the kernel of a non-zero linear form $f \in V^{*}$ is called a hyperplane and the linear form that generates it is a normal to this hyperplane. If $\|f\|_{V^{*}}=1$ then $f$ is called a unit normal. Note that the normal to a hyperplane is, in general, not unique.

## Remark:

Suppose $V$ is equipped with an inner product. Then the notion of normality introduced above coincides with the usual geometrical understanding of normality. Consider a non-zero linear form $f \in V^{*}$. Through the Riesz representation theorem $f$ corresponds to a vector $\bar{n} \in V$ by $f(w)=\langle w, \bar{n}\rangle_{V}$ for all $w \in V . A$ vector is normal to a translate $x+W$ of a subspace in the usual sense if it is orthogonal to all vectors in the subspace $W$. But then

$$
W=\left\{w \in V \mid\langle w, \bar{n}\rangle_{V}=0\right\}=\{w \in V \mid f(w)=0\}=\operatorname{ker}(f)
$$

that is, $W+x$ is in fact a hyperplane as defined above. In this section however, we will not assume that $V$ is equipped with an inner product because the results hold in all finite-dimensional normed spaces - which is the reason why we introduced the notion of normality for linear forms and not for vectors.

A hyperplane $H$ can be described as a level set of its normal $f$ because

$$
\begin{aligned}
H=x+\operatorname{ker}(f) & =\left\{x+v^{\prime} \mid v^{\prime} \in \operatorname{ker}(f)\right\} \\
& =\left\{v=x+v^{\prime} \mid v^{\prime} \in \operatorname{ker}(f), f(v)=f\left(v^{\prime}\right)+f(x)=f(x)\right\} \\
& =\{v \in V \mid f(v)=\alpha\}=: H_{\alpha}(f)
\end{aligned}
$$

where $\alpha:=f(x)$. Given a non-zero linear form $f$ we denote the family of hyperplanes that $f$ is normal to by $H_{\alpha}(f)$ where $\alpha \in \mathbb{R}$. Note $\operatorname{ker}(f)=H_{0}(f)$. Each hyperplane $H_{\alpha}(f)$ bounds the two closed half-spaces $H_{\alpha}^{+}(f):=\{v \in V \mid f(v) \geq \alpha\}$ and $H_{\alpha}^{-}(f):=\{v \in V \mid f(v) \leq \alpha\}$. Both of these half-spaces are convex because $f$ is linear and therefore any intersection of closed half-spaces is a closed convex set. Using the Hahn-Banach theorem and its geometric variants, one can prove the following partial converse.

Proposition 1.2.4 ([Tho96, Theorem 2.1.8, p. 48])
Every closed convex set is the intersection of those closed half-spaces which contain it.

A hyperplane $H:=H_{\alpha}(f)$ is said to be a supporting hyperplane for a closed convex set $K$, if $K \cap H \neq \varnothing$ and either $K \subset H_{\alpha}^{+}(f)$ or $K \subset H_{\alpha}^{-}(f)$. In this case we say that the linear form $f$ supports the convex set $K$ at each point in $K \cap H$. Note that $H_{-\alpha}^{+}(-f)=H_{\alpha}^{-}(f)$ and thus, we may restrict our point of view to hyperplanes of the form $H_{\alpha}^{-}(f)$. If $H_{\alpha}^{-}(f)$ is a supporting half-space for $K$ then $\frac{1}{\|f\|_{V^{*}}} f$ is called an outer unit normal.

With all this notation at hand, we can introduce an important mapping corresponding to a given bounded, closed and convex set $K$. The function

$$
h_{K}: V^{*} \rightarrow \mathbb{R}, f \mapsto \sup _{x \in K} f(x)
$$

is called the support function of $\boldsymbol{K}$. This function has the following properties.

## Lemma 1.2.5

(i) The support function is sublinear on $V^{*}$, that is,

$$
\begin{aligned}
h_{K}(\alpha f) & =\alpha h_{K}(f), \quad \text { for } \alpha \geq 0 \\
h_{K}(f+g) & \leq h_{K}(f)+h_{K}(g) .
\end{aligned}
$$

(ii) The support function respects Minkowski sums, that is,

$$
\begin{aligned}
h_{\alpha K} & =\alpha h_{K}, \quad \text { for } \alpha \geq 0 \\
h_{K+K^{\prime}} & =h_{K}+h_{K^{\prime}}
\end{aligned}
$$

(iii) If $0 \in K$ then $h_{K}(f) \geq 0$ and if 0 is an interior point of $K$ then $h_{K}(f)>0$ for all $f \neq 0$.
(iv) If $K$ is symmetric then $h_{K}(f)=h_{K}(-f)$ and $h_{K}(f)=\sup _{x \in K}|f(x)|$.
(v) If $K$ is the unit ball of $V$ then $h_{K}=\|\cdot\|_{V^{*}}$.
(vi) If $K \subset K^{\prime}$ then $h_{K} \leq h_{K^{\prime}}$.

Proof: Part (i) follows from the definition of $h_{K}$. For part (ii) calculate

$$
h_{\alpha K}(f)=\sup _{x \in \alpha K} f(x)=\sup _{y \in K} f(\alpha y)=\alpha \sup _{y \in K} f(y)=\alpha h_{K}(f)
$$

and

$$
\begin{aligned}
h_{K+K^{\prime}}(f) & =\sup _{x \in K+K^{\prime}} f(x)=\sup _{y \in K, y^{\prime} \in K^{\prime}} f\left(y+y^{\prime}\right) \\
& =\sup _{y \in K, y^{\prime} \in K^{\prime}} f(y)+f\left(y^{\prime}\right)=\sup _{y \in K} f(y)+\sup _{y^{\prime} \in K^{\prime}} f\left(y^{\prime}\right) \\
& =h_{K}(f)+h_{K^{\prime}}(f)
\end{aligned}
$$

for $f \in V^{*}$ and $\alpha \geq 0$.
To prove part (iii), suppose that $h_{K}(f)=\sup _{y \in K} f(y)<0$ for some $f \in V^{*}$. Since $0 \in K$, we have the contradiction

$$
0=f(0) \leq \sup _{y \in K} f(y)<0
$$

Let 0 be an interior point of $K$ and suppose further there is $f \in V^{*}$ with $h_{K}(f)=\sup _{x \in K} f(x)=0$, that is, $\left.f\right|_{K} \equiv 0$. Since 0 is an interior point there is a radius $\varepsilon>0$ such that the norm ball $B_{\varepsilon}(0)$ centered at 0 is contained in $K$. Moreover, $\left.f\right|_{B_{\varepsilon}(0)} \equiv 0$. This already implies $f \equiv 0$ since any $0 \neq x \in V$ can be scaled by the factor $\frac{\varepsilon}{2\|x\|_{V}}$ such that $y:=\frac{\varepsilon}{2\|x\|_{V}} x \in B_{\varepsilon}(0)$. The linearity of $f$ then yields $f(x)=\frac{2\|x\|_{V}}{\varepsilon} f(y)=0$. For part (iv), note that if $K$ is symmetric (that is, $K=(-1) K$ ) then

$$
\begin{aligned}
h_{K}(-f) & =\sup _{x \in K}(-f)(x)=\sup _{x \in K} f(-x) \\
& =\sup _{-x \in K} f(x)=\sup _{x \in(-1) K} f(x)=\sup _{x \in K} f(x) \\
& =h_{K}(f) .
\end{aligned}
$$

For part $(v)$ just note that $0 \in K$. Part (iii) and the definition of the dual norm immediately yield the claim since $h_{K}(f)=\sup _{x \in K}|f(x)|=\|f\|_{V^{*}}$. Part (vi) is clear because $\sup _{x \in K} f(x) \leq \sup _{x \in K^{\prime}} f(x)$ if $K \subset K^{\prime}$.

Since $h_{K}$ is a positively homogeneous mapping its image is completely determined by the image of linear forms of unit length. As mentioned earlier any closed convex set is contained in the intersection of its supporting half-spaces. Although very important this is merely a qualitative statement. However, the next result shows that the support function $h_{K}$ completely describes these supporting half-spaces and provides a geometrical meaning to its value at a unit length linear form.


Figure 1.1: Distance of $0 \in K$ to $H_{\alpha}(f)$

## Proposition 1.2.6

Let $f \in V^{*}$ be a non-zero linear form. Suppose $K$ is a closed bounded and convex set which contains the origin as an interior point. Define $\alpha:=h_{K}(f)>0$. Then the hyperplane $H_{\alpha}(f)$ is a supporting hyperplane for $K$ and $K$ is contained in the half-space $H_{\alpha}^{-}(f)=\left\{v \in V \mid f(v) \leq h_{K}(f)\right\}$ (that is, $f$ is an outer normal to $H_{\alpha}(f)$ ). In particular, $h_{K}\left(\frac{f}{\|f\|_{V^{*}}}\right)$ gives the distance of the origin to the hyperplane $H_{\alpha}(f)=\left\{v \in V \mid f(v)=h_{K}(f)\right\}$, that is,

$$
\begin{equation*}
h_{K}\left(\frac{f}{\|f\|_{V^{*}}}\right)=\operatorname{dist}\left(0, H_{\alpha}(f)\right) . \tag{1.2.14}
\end{equation*}
$$

Proof: Since $K$ is compact as a closed and bounded set in $V$ and $f$ is continuous, there is $x \in K$ such that $f(x)=\sup _{y \in K} f(y)=h_{K}(f)=\alpha$. Then $x \in K \cap H_{\alpha}(f)$. Further, $K \subset H_{\alpha}^{-}(f)$ because $f(z) \leq \sup _{y \in K} f(y)=h_{K}(f)$ for all $z \in K$. Hence, $H_{\alpha}(f)$ is a supporting hyperplane for $K$.

Note that by the linearity of $f, \frac{1}{t} H_{\alpha}(f)=H_{t \alpha}(f)$ for all $t \neq 0$. Therefore, the hyperplane $H_{\alpha}(f)$ parametrises the open half-space $H_{0}^{+}(f)$ which is bounded by the kernel of $f$. Similarly, the hyperplane $H_{-\alpha}(f)$ parametrises the open half-space $H_{0}^{-}(f)$.

Then we calculate

$$
\begin{aligned}
\operatorname{dist}\left(0, H_{\alpha}(f)\right) & =\inf _{z \in H_{\alpha}(f)}\|z\|_{V}=\alpha \inf _{z \in H_{\alpha}(f)}\left(\frac{\alpha}{\|z\|_{V}}\right)^{-1} \\
& =\alpha\left(\sup _{z \in H_{\alpha}(f)} \frac{\alpha}{\|z\|_{V}}\right)^{-1}=h_{K}(f)\left(\sup _{z \in H_{\alpha}(f)} \frac{f(z)}{\|z\|_{V}}\right)^{-1} \\
& =h_{K}(f)\left(\sup _{z \in H_{\alpha}(f)} f\left(\frac{z}{\|z\|_{V}}\right)\right)^{-1}
\end{aligned}
$$

The expression $\frac{z}{\|z\|_{V}}$ does not change its value if $z$ is scaled by an arbitrary positive factor $t>0$. Thus, we can rewrite the supremum,

$$
\sup _{\substack{z \in H_{\alpha}(f) \\ t>0}} f\left(\frac{\frac{z}{t}}{\left\|\frac{z}{t}\right\|_{V}}\right)=\sup _{\substack{z \in \frac{1}{t} H_{\alpha}(f) \\ t>0}} f\left(\frac{z}{\|z\|_{V}}\right)=\sup _{\substack{z \in H_{+\alpha}(f) \\ t>0}} f\left(\frac{z}{\|z\|_{V}}\right)=\sup _{z \in V, f(z)>0} f\left(\frac{z}{\|z\|_{V}}\right) .
$$

Notice that

$$
\sup _{z \in V, f(z)>0} f\left(\frac{z}{\|z\|_{V}}\right)=\sup _{z \in V, f(z)>0}-f\left(\frac{-z}{\|-z\|_{V}}\right)=\sup _{z \in V, f(z)<0}-f\left(\frac{z}{\|z\|_{V}}\right)
$$

which implies

$$
\sup _{z \in V, f(z)>0} f\left(\frac{z}{\|z\|_{V}}\right)=\sup _{z \in V, f(z) \neq 0}\left|f\left(\frac{z}{\|z\|_{V}}\right)\right| .
$$

This finally yields

$$
\begin{aligned}
\operatorname{dist}\left(0, H_{\alpha}(f)\right) & =h_{K}(f)\left(\sup _{z \in V, f(z) \neq 0}\left|f\left(\frac{z}{\|z\|_{V}}\right)\right|\right)^{-1}=h_{K}(f)\left(\sup _{z \in V, z \neq 0}\left|f\left(\frac{z}{\|z\|_{V}}\right)\right|\right)^{-1} \\
& =h_{K}(f) \frac{1}{\|f\|_{V^{*}}}=h_{K}\left(\frac{f}{\|f\|_{V^{*}}}\right)
\end{aligned}
$$

The previous result will be used repeatedly in the proofs of explicit calculations on polytopes carried out in Section 1.2.5.

### 1.2.3 Polyhedra and polarity of convex sets

Recall that if $0 \in K$ is an interior point then $h_{K}(f)>0$ if $f \neq 0$. Therefore,

$$
H_{\alpha}^{-}(f)=\left\{v \in V \mid f(v) \leq h_{K}(f)\right\}=\left\{v \in V \left\lvert\, \frac{1}{h_{K}(f)} f(v) \leq 1\right.\right\}=H_{1}^{-}(\tilde{f})
$$

where $\tilde{f}:=\frac{1}{h_{K}(f)} f$. Thus, we may restrict our point of view further to half-spaces described by $H_{1}^{-}(f)$. By the previous proposition, we can then describe a closed bounded and convex set which contains the origin as an interior point by

$$
\begin{equation*}
K=\bigcap_{f \in V^{*}} H_{h_{K}(f)}^{-}(f)=\bigcap_{f \in V^{*}} H_{1}^{-}\left(\frac{1}{h_{K}(f)} f\right) \tag{1.2.15}
\end{equation*}
$$

A polyhedral set $P$ is the intersection of a finite number of half-spaces, that is,

$$
P=\bigcap_{i=1}^{M} H_{\alpha_{i}}^{-}\left(f^{i}\right)
$$

for non-zero linear forms $f^{i} \in V^{*}$ and scalars $\alpha_{i} \in \mathbb{R}, i=1,2, \ldots, M$. We assume here that the number $M$ is minimal in the sense that none of the half-spaces $H_{\alpha_{j}}^{-}\left(f^{j}\right)$ may be omitted in the intersection above (Brøndsted calls this an "irreducible" representation, see [Brø83, p. 52]). One of the main theorems in convexity theory of polytopes is the following statement.

Theorem 1.2.7 ([Brø83, Chapter 9, Theorem 9.2])
A non-empty subset $P$ of $V$ is a polytope if and only if it is a bounded polyhedral set.

This theorem gives two representations of a polytope - the vertex representation on the one hand and the half-space representation on the other. An elegant proof, using only elementary arguments, can be found in [Mat02, pp. 84-85]. There is a wide range of results on the combinatorial properties of such sets - the text books [Brø83] and [Grü03] provide good references for this topic. Regarding combinatorics of polytopes, let us only note here that a symmetric polytope has an even number of vertices and an even number of facets, which can be deduced by using the vertex and half-space representation respectively.

For the application in the next chapter we need another concept that plays an important role in convexity theory. For each closed convex set $K$ we can define its polar set $K^{\circ}$ in the dual space $V^{*}$ by

$$
\begin{align*}
K^{\circ} & :=\left\{f \in V^{*} \mid f(x) \leq 1 \text { for all } x \in K\right\}  \tag{1.2.16}\\
& =\bigcap_{x \in K}\left\{f \in V^{*} \mid f(x) \leq 1\right\} .
\end{align*}
$$

It is clear that the polar set is non-empty since $0 \in K^{\circ}$. The polar operation reverses set inclusions, for if $K_{1} \subset K_{2}$ then $K_{2}^{\circ} \subset K_{1}^{\circ}$. Moreover, the next properties hold to be true.

## Lemma 1.2.8

The polar set $K^{\circ}$ can be described through the support function of $K$ as $K^{\circ}=\left\{f \in V^{*} \mid h_{K}(f) \leq 1\right\}$. If $B$ is the unit ball in $V$ then $B^{\circ}$ is the unit ball in $V^{*}$. The polar set is a convex and closed set. If $K$ is bounded then $0 \in V^{*}$ is an interior point of $K^{\circ}$. Furthermore, if $0 \in V$ is an interior point of $K$ then $K^{\circ}$ is bounded.

Proof: Let $f \in K^{\circ}$, that is, $f(x) \leq 1$ for all $x \in K$. But then this also holds for the least upper
bound, so $h_{K}(f)=\sup _{x \in K} f(x) \leq 1$. Conversely, if $f \in V^{*}$ such that $h_{K}(f)=\sup _{x \in K} f(x) \leq 1$, then $f(x) \leq 1$ for all $x \in K$, that is, $f \in K^{\circ}$. If $B$ is the unit ball of $V$ then part $(v)$ of Lemma 1.2.5 yields

$$
B^{\circ}=\left\{f \in V^{*} \mid h_{B}(f) \leq 1\right\}=\left\{f \in V^{*} \mid\|f\|_{V^{*}} \leq 1\right\}
$$

For any convex combination of elements of $K^{\circ}$

$$
(t f+(1-t) g)(x)=t f(x)+(1-t) g(x) \leq t+(1-t)=1
$$

for all $x \in K$, therefore $K^{\circ}$ is convex. Further, let $\left(f_{k}\right)_{k \in \mathbb{N}} \subset K^{\circ}$ and $f_{k} \underset{k \rightarrow \infty}{ } f \in V^{*}$. Then for all $x \in K$

$$
\begin{aligned}
f(x) & =f(x)-f_{k}(x)+f_{k}(x) \\
& \leq\left\|f-f_{k}\right\|_{V^{*}}+f_{k}(x) \\
& \leq\left\|f-f_{k}\right\|_{V^{*}}+1 \xrightarrow[k \rightarrow \infty]{ } 1 .
\end{aligned}
$$

Thus, $f \in K^{\circ}$ and $K^{\circ}$ is closed. If $B_{R}(0)=R B_{1}(0)=R B$ is the norm ball in $V$ of radius $R>0$ centered at the origin, note that by Lemma 1.2.5 ( $v$ )

$$
h_{B_{R}(0)}(f)=R h_{B}(f)=R\|f\|_{V^{*}}
$$

and because $B^{\circ}$ is the unit ball in $V^{*}$ it follows that

$$
\begin{aligned}
\left(B_{R}(0)\right)^{\circ} & =\left\{f \in V^{*} \mid h_{B_{R}(0)}(f) \leq 1\right\} \\
& =\left\{f \in V^{*} \mid R\|f\|_{V^{*}} \leq 1\right\} \\
& =\left\{f \in V^{*} \mid\|f\|_{V^{*}} \leq R^{-1}\right\} \\
& =R^{-1} B^{\circ}=: B_{R^{-1}}^{*}(0) .
\end{aligned}
$$

In other words, $\left(B_{R}(0)\right)^{\circ}$ is the norm ball in the dual space $V^{*}$ of radius $R^{-1}$ centered at the origin. If $K$ is bounded then there is a radius $R>0$ such that $K \subset B_{R}(0)$. As the polar operation reverses inclusions, we have $B_{R^{-1}}^{*}(0) \subset K^{\circ}$, that is, 0 is an interior point of $K^{\circ}$. Conversely, if $B_{R}(0) \subset K$ for some $R>0$ the same argumentation yields $K^{\circ} \subset B_{R^{-1}}^{*}(0)$, that is, $K^{\circ}$ is bounded.

## Proposition 1.2.9

Let $P$ be a polytope which contains the origin as an interior point. Then $P^{\circ}$ is a polytope.

Proof: Suppose $P=\operatorname{conv}\left(a_{1}, a_{2}, \ldots, a_{N}\right)$. Then $P$ is closed bounded and convex due to Lemma 1.2.3. By Lemma 1.2.8 we know that $P^{\circ}$ is also closed, convex and contains $0 \in V^{*}$ as an interior point. In addition,

$$
P^{\circ}=\bigcap_{x \in P}\left\{f \in V^{*} \mid f(x) \leq 1\right\}=\bigcap_{i=1}^{N}\left\{f \in V^{*} \mid f\left(a_{i}\right) \leq 1\right\}
$$

The second equality follows from the fact that $a_{i} \in P$, so $f\left(a_{i}\right) \leq 1$ for $f \in P^{\circ}$. Conversely, suppose $f \in V^{*}$ such that $f\left(a_{i}\right) \leq 1$ for all $i=1,2, \ldots, N$. We then have $f(x)=f\left(\sum_{i=1}^{N} \lambda_{i} a_{i}\right)=$ $\sum_{i=1}^{N} \lambda_{i} f\left(a_{i}\right) \leq \sum_{i=1}^{N} \lambda_{i}=1$ for all $x \in P$, that is, $f \in P^{\circ}$.

For any normed space the canonical mapping $\Phi: V \rightarrow V^{* *}, \Phi(v)(f):=f(v)$ is linear and preserves the norm, that is, $\|\Phi(v)\|_{V^{* *}}=\|v\|_{V}$ (see [Wer00, Satz III.3.1, p. 105]). This implies that $\Phi$ is injective, for if $\Phi(v)=0$ then $v$ must be zero. Since $V$ and $V^{*}$ and $V^{* *}=\left(V^{*}\right)^{*}$ are all vector spaces of the same dimension $n$, the map $\Phi$ is already surjective by the rank-nullity theorem because $\operatorname{dim}(\operatorname{im}(\Phi))=\operatorname{dim}(V)-\operatorname{dim}(\operatorname{ker}(\Phi))=n-0=n$. Thus,

$$
\left\{f \in V^{*} \mid f\left(a_{i}\right) \leq 1\right\}=\left\{f \in V^{*} \mid \Phi\left(a_{i}\right)(f) \leq 1\right\}=H_{1}^{-}\left(\Phi\left(a_{i}\right)\right)
$$

is a half-space in $V^{*}$, whence $P^{\circ}$ is a polyhedral set as the intersection of those finitely many half-spaces. Furthermore, because $P$ contains the origin as an interior point, the polar set $P^{\circ}$ is bounded due to Lemma 1.2.8. By Theorem 1.2.7, $P^{\circ}$ then is a polytope.

More precisely (or rather "conversely" in terms of vertex and half-space representations) the following statement holds.

Lemma 1.2.10 ([Grü03, Exercise 3.4.7, p. 49] or [Mat02, Exercise 5.1.6 (b), p. 81])
If the convex polytope $P$ (containing the origin as an interior point) is given via its half-space representation as

$$
P=\bigcap_{i=1}^{M}\left\{v \in V \mid f^{i}(v) \leq 1\right\}
$$

where $f^{i} \in V^{*}$ for $i=1,2, \ldots, M$ then the polar set is given via the vertex representation

$$
P^{\circ}=\operatorname{conv}\left(f^{1}, f^{2}, \ldots, f^{M}\right)
$$

Proof: Define $Q:=\operatorname{conv}\left(f^{1}, f^{2}, \ldots, f^{M}\right) \subset V^{*}$. The set $Q^{\circ}=\left\{\tilde{v} \in V^{* *} \mid \tilde{v}(f) \leq 1\right.$ for all $\left.f \in Q\right\} \subset$ $V^{* *}$ is its polar set. We will show $Q^{\circ}=P$. Referring to the previous proof, recall the canonical
isomorphism $\Phi: V \rightarrow V^{* *}$ where $\Phi(v)(f):=f(v)$. Using this identification, we consider $Q^{\circ}$ to be a subset of $V$ and

$$
Q^{\circ}=\{v \in V \mid f(v) \leq 1 \text { for all } f \in Q\}=\bigcap_{f \in Q}\{v \in V \mid f(v) \leq 1\}
$$

Suppose $v \in Q^{\circ}$, then we have in particular $f^{i} \in Q$ for $i=1,2, \ldots, M$ so that $f^{i}(v) \leq 1$. Conversely, suppose $v \in V$ so that $f^{i}(v) \leq 1$ for all $i=1,2, \ldots, M$. Then for all $f \in Q$ we have $f(v)=$ $\left(\sum_{i=1}^{N} \lambda_{i} f^{i}\right)(v)=\sum_{i=1}^{N} \lambda_{i} f^{i}(v) \leq \sum_{i=1}^{N} \lambda_{i}=1$. Thus, the last identity can be rewritten as

$$
Q^{\circ}=\bigcap_{i=1}^{M}\left\{v \in V \mid f^{i}(v) \leq 1\right\}=P
$$

By taking the polar yet again, this implies $P^{\circ}=Q^{\circ \circ}$. (It is understood that this equality sign refers to the identification of $V^{*}$ with $\left(V^{*}\right)^{* *}$ via the mapping $\Phi^{*}: V^{*} \rightarrow\left(V^{*}\right)^{* *}, \Phi^{*}(v)(f)=f(v)$.) Furthermore, the polytope $Q$ is bounded (Lemma 1.2.3) so that by Lemma 1.2.8 the set $Q^{\circ}$ contains the origin as an interior point and is closed. Then $Q^{\circ \circ}=Q$ (see [Tho96, Theorem 2.2.6, p. 50]) which proves the claim - again it is understood that the equality sign refers to the identification of $V^{*}$ with $V^{* * *}$ via the mapping $\Phi^{*}$.

Note that the number of facets and vertices of a polytope and its polar set interchange their roles. This is one example of duality in the theory of convex polytopes and gives rise to the so-called "face-lattice" (see [Brø83, Chapter 1, §5, p. 29f]).

### 1.2.4 Mixed volume and the Minkowski inequality

We introduce a quantity that relates a collection of convex sets with each other - the mixed volume. For the application in Chapter 2 we prove the Minkowski inequality for convex bodies.

A convex body is a compact, convex subset of $V$ with non-empty interior. The collection of all convex bodies in $V$ is denoted by $\mathcal{C}_{b}=\mathcal{C}_{b}(V)$. In Lemma 1.2.3 we have seen that any full-dimensional polytope is a convex body. We need full-dimensionality here to ascertain a non-empty interior.

To motivate the notion of mixed volume, the following observation is useful. Using the homogeneity of the Hausdorff measure on $V$, one sees that the measure of a scalar multiple $\alpha=\alpha_{1}+\ldots+\alpha_{k}$ of a set $K$ is

$$
\mathcal{H}^{n}\left(\alpha_{1} K+\ldots+\alpha_{k} K\right)=\left(\alpha_{1}+\ldots+\alpha_{k}\right)^{n} \mathcal{H}^{n}(K)=\left(\sum_{i_{1}, \ldots, i_{n}=1}^{k} \alpha_{i_{1}} \ldots \alpha_{i_{n}}\right) \mathcal{H}^{n}(K)
$$

that is, $\mathcal{H}^{n}\left(\alpha_{1} K+\ldots+\alpha_{k} K\right)$ is a homogeneous polynomial of degree $n$ in the coefficients $\alpha_{i}$. In greater generality, one can show that this holds for a non-negative linear combination of convex bodies. In particular, for any collection of convex bodies $\left\{K_{1}, \ldots, K_{k}\right\} \subset V$ and non-negative coefficients $\alpha_{i} \geq 0$ the map $\left(\alpha_{1}, \ldots, \alpha_{k}\right) \mapsto \mathcal{H}^{n}\left(\alpha_{1} K_{1}+\ldots+\alpha_{k} K_{k}\right)$ is a homogeneous polynomial of degree $n$ (see [Tho96, Theorem 2.3.6, p. 55]). Hence, it can be expressed as

$$
\begin{equation*}
\mathcal{H}^{n}\left(\alpha_{1} K_{1}+\ldots+\alpha_{k} K_{k}\right)=\sum_{i_{1}, \ldots, i_{n}=1}^{k} \mathrm{~V}\left(K_{i_{1}}, \ldots, K_{i_{n}}\right) \alpha_{i_{1}} \ldots \alpha_{i_{n}} \tag{1.2.17}
\end{equation*}
$$

where the coefficient term $\mathrm{V}\left(K_{i_{1}}, \ldots, K_{i_{n}}\right)$ is symmetric in its arguments (because the Minkowski sum in the left hand side is commutative) and is called the mixed volume of the convex bodies $K_{i_{1}}, \ldots, K_{i_{n}}$. Note that the mixed volume is always a function of $n$ arguments - the dimension of the ambient space $V$ - independent of the length of the linear combination. By inspection, we can see $\mathrm{V}(K, \ldots, K)=\mathcal{H}^{n}(K)$. Furthermore, if any of the sets $K_{i}$ is translated linearly by some $x \in V$ the invariance of Hausdorff measure (on the left hand side of (1.2.17)) under translations implies that the mixed volume (on the right hand side of (1.2.17)) is invariant under translations of any of its arguments. The symmetry of mixed volume allows to reorder and collect like terms in the above polynomial expression. There are $\binom{n}{j_{1}, \ldots, j_{k}}:=\binom{n}{j_{1}}\binom{n-j_{1}}{j_{2}} \cdots\binom{n-j_{1}-\ldots-j_{k-1}}{j_{k}}$ terms such that $j_{1}$ of the $n$ indices $i_{1}, \ldots, i_{n}$ are equal to $1, j_{2}$ of the remaining $n-j_{1}$ indices are equal to 2 , etc. The mixed volume then only depends on the new indices $j_{1}, \ldots, j_{k}$ and will be denoted by

$$
\mathrm{V}\left(K_{1}\left[j_{1}\right], \ldots, K_{k}\left[j_{k}\right]\right):=\mathrm{V}\left(K_{1}, K_{1}, \ldots, K_{1}, K_{2}, K_{2}, \ldots, K_{2}, \ldots, K_{k}, K_{k}, \ldots, K_{k}\right)
$$

where $K_{l}$ appears $j_{l}$ times. Thus, the polynomial expression can be rewritten as

$$
\mathcal{H}^{n}\left(\alpha_{1} K_{1}+\ldots+\alpha_{k} K_{k}\right)=\sum_{j_{1}+\ldots+j_{k}=n}\binom{n}{j_{1}, \ldots, j_{k}} \mathrm{~V}\left(K_{1}\left[j_{1}\right], \ldots, K_{k}\left[j_{k}\right]\right) \alpha_{1}^{j_{1}} \ldots \alpha_{k}^{j_{k}} .
$$

For the Minkowski sum $s K+t K^{\prime}$ of two convex bodies $K$ and $K^{\prime}$ this formula becomes

$$
\begin{equation*}
\mathcal{H}^{n}\left(s K+t K^{\prime}\right)=\sum_{j=0}^{n}\binom{n}{j} \mathrm{~V}\left(K[n-j], K^{\prime}[j]\right) s^{n-j} t^{j} \tag{1.2.18}
\end{equation*}
$$

If $s=1$ in the last expression, we conclude the following formula.

## Corollary 1.2.11

The mixed volume of $K[n-1]$ and $K^{\prime}=K^{\prime}[1]$ can be calculated by

$$
\begin{equation*}
\mathrm{V}\left(K[n-1], K^{\prime}\right)=\frac{1}{n} \lim _{t \rightarrow 0} \frac{\mathcal{H}^{n}\left(K+t K^{\prime}\right)-\mathcal{H}^{n}(K)}{t} \tag{1.2.19}
\end{equation*}
$$

There are several inequalities that relate the different mixed volumes of a collection of convex bodies. An important one is the Brunn-Minkowski inequality (see [Tho96, Theorem 2.3.5, p. 55]) which will be included in the proof below. We will need the following consequence of this inequality in the application in Chapter 2.

Proposition 1.2.12 (Minkowski inequality for convex bodies, [BF71, §49, p. 91])
If $K$ and $K^{\prime}$ are convex bodies then

$$
\left(\mathrm{V}\left(K[n-1], K^{\prime}\right)\right)^{n} \geq\left(\mathcal{H}^{n}(K)\right)^{n-1} \mathcal{H}^{n}\left(K^{\prime}\right)
$$

Proof: The Brunn-Minkowski inequality states that

$$
\left(\mathcal{H}^{n}\left((1-\alpha) K+\alpha K^{\prime}\right)\right)^{\frac{1}{n}} \geq(1-\alpha)\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}}+\alpha\left(\mathcal{H}^{n}\left(K^{\prime}\right)\right)^{\frac{1}{n}} .
$$

In other words, the map $\mathcal{C}_{b} \rightarrow \mathbb{R}, K \mapsto\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}}$ is concave. Then for fixed convex bodies $K, K^{\prime} \subset V$ the map $[0,1] \rightarrow \mathbb{R}, \alpha \mapsto\left(\mathcal{H}^{n}\left((1-\alpha) K+\alpha K^{\prime}\right)\right)^{\frac{1}{n}}$ is a concave function. Define the function

$$
\Phi:[0,1] \rightarrow \mathbb{R}, \alpha \mapsto\left(\mathcal{H}^{n}\left((1-\alpha) K+\alpha K^{\prime}\right)\right)^{\frac{1}{n}}-(1-\alpha)\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}}-\alpha\left(\mathcal{H}^{n}\left(K^{\prime}\right)\right)^{\frac{1}{n}}
$$

which is also concave since the additional terms are linear in $\alpha$. Note that $\Phi(0)=\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}}-$ $\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}}=0$ and similarly $\Phi(1)=0$. Note further that $\Phi$ is a polynomial function in the variable $\alpha$ due to expression (1.2.18) and is thus differentiable on $(0,1)$ and the one-sided limits for the difference quotient at $\alpha=0$ and $\alpha=1$ exist. The concavity and continuity of $\Phi$ imply that its (right hand) derivative at $\alpha=0$ must be non-negative. We will shortly see that this implies the claim but first let us calculate the derivative of $\mathcal{H}^{n}\left((1-\alpha) K+\alpha K^{\prime}\right)$ at $\alpha=0$. Equation (1.2.18) yields

$$
\begin{aligned}
& \left.\frac{\mathrm{d}}{\mathrm{~d} \alpha}\right|_{\alpha=0}\left(\mathcal{H}^{n}\left((1-\alpha) K+\alpha K^{\prime}\right)\right) \\
= & \left.\frac{\mathrm{d}}{\mathrm{~d} \alpha}\right|_{\alpha=0}\left(\sum_{j=0}^{n}\binom{n}{j} \mathrm{~V}\left(K[n-j], K^{\prime}[j]\right)(1-\alpha)^{n-j} \alpha^{j}\right) \\
= & \left.\sum_{j=0}^{n}\binom{n}{j} \mathrm{~V}\left(K[n-j], K^{\prime}[j]\right) \frac{\mathrm{d}}{\mathrm{~d} \alpha}\right|_{\alpha=0}\left((1-\alpha)^{n-j} \alpha^{j}\right)
\end{aligned}
$$

$$
\begin{align*}
= & \left.\mathrm{V}\left(K[n], K^{\prime}[0]\right) \frac{\mathrm{d}}{\mathrm{~d} \alpha}\right|_{\alpha=0}(1-\alpha)^{n}+\left.n \mathrm{~V}\left(K[n-1], K^{\prime}[1]\right) \frac{\mathrm{d}}{\mathrm{~d} \alpha}\right|_{\alpha=0}\left((1-\alpha)^{n-1} \alpha\right) \\
& +\left.\sum_{j=2}^{n}\binom{n}{j} \mathrm{~V}\left(K[n-j], K^{\prime}[j]\right) \frac{\mathrm{d}}{\mathrm{~d} \alpha}\right|_{\alpha=0}\left((1-\alpha)^{n-j} \alpha^{j}\right) \\
= & -n \mathcal{H}^{n}(K)+n \mathrm{~V}\left(K[n-1], K^{\prime}[1]\right) \\
& +\left.\sum_{j=2}^{n}\binom{n}{j} \mathrm{~V}\left(K[n-j], K^{\prime}[j]\right)\left(j \alpha^{j-1}(1-\alpha)^{n-j}-(n-j) \alpha^{j}(1-\alpha)^{n-j-1}\right)\right|_{\alpha=0} \\
= & n\left(\mathrm{~V}\left(K[n-1], K^{\prime}[1]\right)-\mathcal{H}^{n}(K)\right) . \tag{1.2.20}
\end{align*}
$$

The last equality uses the fact that all summands with index $j \geq 2$ contain a power of $\alpha$ which vanishes at 0 . Proceeding, one gets

$$
\begin{aligned}
\left.\frac{\mathrm{d} \Phi}{\mathrm{~d} \alpha}\right|_{\alpha=0} & =\left.\frac{\mathrm{d}}{\mathrm{~d} \alpha}\right|_{\alpha=0}\left[\left(\mathcal{H}^{n}\left((1-\alpha) K+\alpha K^{\prime}\right)\right)^{\frac{1}{n}}-(1-\alpha)\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}}-\alpha\left(\mathcal{H}^{n}\left(K^{\prime}\right)\right)^{\frac{1}{n}}\right] \\
& =\left.\frac{1}{n}\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}-1} \frac{\mathrm{~d}}{\mathrm{~d} \alpha}\right|_{\alpha=0}\left[\mathcal{H}^{n}\left((1-\alpha) K+\alpha K^{\prime}\right)\right]+\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}}-\left(\mathcal{H}^{n}\left(K^{\prime}\right)\right)^{\frac{1}{n}} .
\end{aligned}
$$

Using (1.2.20), the last statement can be rewritten as

$$
\begin{aligned}
\left.\frac{\mathrm{d} \Phi}{\mathrm{~d} \alpha}\right|_{\alpha=0} & =\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}-1}\left[\mathrm{~V}\left(K[n-1], K^{\prime}[1]\right)-\mathcal{H}^{n}(K)\right]+\left(\mathcal{H}^{n}(K)\right)^{\frac{1}{n}}-\left(\mathcal{H}^{n}\left(K^{\prime}\right)\right)^{\frac{1}{n}} \\
& =\left(\mathcal{H}^{n}(K)\right)^{-\frac{n-1}{n}} \mathrm{~V}\left(K[n-1], K^{\prime}[1]\right)-\left(\mathcal{H}^{n}\left(K^{\prime}\right)\right)^{\frac{1}{n}} \\
& =\left(\mathcal{H}^{n}(K)\right)^{-\frac{n-1}{n}}\left\{\mathrm{~V}\left(K[n-1], K^{\prime}[1]\right)-\left(\mathcal{H}^{n}(K)\right)^{\frac{n-1}{n}}\left(\mathcal{H}^{n}\left(K^{\prime}\right)\right)^{\frac{1}{n}}\right\} .
\end{aligned}
$$

Thus, the non-negativity of $\left.\frac{\mathrm{d} \Phi}{\mathrm{d} \alpha}\right|_{\alpha=0}$ implies the claim.

### 1.2.5 An explicit formula for the mixed volume of a polytope and a convex set

There is an intuitive geometrical way of calculating the Minkowski sum of a polytope $A$ and a convex body $K$. First, consider the set $K^{\prime}:=K-k$ where $k$ is an interior point of $K$. Secondly, attach the set $K^{\prime}$ to each vertex of $A$ and take the convex hull (see Figure 1.2 (a)). Lastly, translate by the vector $k$, that is, $A+K=(A+(K-k))+k$.

## Lemma 1.2.13

Suppose $A_{N}=A$ is a polytope and $K$ a convex body in $V$ and $k$ an interior point of $K$. Their

Minkowski sum is given by

$$
A+K=\operatorname{conv}\left(\bigcup_{i=1}^{N}\left(a_{i}+K^{\prime}\right)\right)+k
$$

where $K^{\prime}:=K-k$.
Proof: Without loss of generality we may assume the origin is an interior point of $K$ and $k=0$. For the sake of brevity let $C:=\operatorname{conv}\left(\bigcup_{i=1}^{N}\left(a_{i}+K\right)\right)$. Clearly, $C \subset A+K$, for if $c \in C$ then $c=\sum_{l=1}^{n+1} \lambda_{l}\left(a_{i_{l}}+k_{l}\right)$ where $\sum_{l=1}^{n+1} \lambda_{l}=1, k_{l} \in K$ and $a_{i_{l}} \in\left\{a_{1}, a_{2}, \ldots, a_{N}\right\}$ for $l=1,2, \ldots, N$. Since both $A$ and $K$ are convex $c=\sum_{l=1}^{n+1} \lambda_{l} a_{i_{l}}+\sum_{l=1}^{n+1} \lambda_{l} k_{l} \in A+K$.

Conversely, let $c \in A+K$. Since each point in the polytope $A$ is a convex combination of the vertices $a_{1}, a_{2}, \ldots, a_{N}$, it follows that

$$
c=\sum_{i=1}^{N} \lambda_{i} a_{i}+k=\sum_{i=1}^{N} \lambda_{i}\left(a_{i}+k\right) \in C
$$

where $\sum_{i=1}^{N} \lambda_{i}=1$.

As a motivation for the next result consider a polytope $A \subset V$ with $N$ facets $F_{i}$. Recall that facets are the $(n-1)$-dimensional extreme subsets of $A$. Without loss of generality, suppose $A$ contains the origin $0 \in V$ as an interior point. An intuitive way of calculating its volume is the following geometric sequence of steps. First, one triangulates the polytope by connecting each of its facets with the origin which creates $N$ pyramids $P_{i}=\operatorname{conv}\left(F_{i} \cup\{0\}\right)$. Secondly, one adds up the volume of each pyramid. Lastly, one recalls that the volume of a pyramid is simply the $(n-1)$-dimensional volume of the "base" times its "height" divided by the dimension $n$. (This can be proven by using Cavalieri's principle, see e.g. [For09, $\S 5$, Satz 3, p. 49].) To put this into more mathematical terms, intuitively, the volume of $A$ is

$$
\mathcal{H}^{n}(A)=\sum_{i=1}^{N} \mathcal{H}^{n}\left(P_{i}\right)=\frac{1}{n} \sum_{i=1}^{N} \mathcal{H}^{n-1}\left(F_{i}\right) h_{A}\left(f^{i}\right)
$$

where $f^{i}$ denotes the outer unit normal (see p. 36) of the facet $F_{i}$ and $h_{A}$ is the support function of $A$. Due to Proposition 1.2.6 the value $h_{A}\left(f^{i}\right)$ is the distance of the affine subspace containing the facet $F_{i}$ to the origin (and so the height of the pyramid $P_{i}$ ).

With the help of Corollary 1.2.11 the mixed volume of a polytope and a convex body can be calculated geometrically. This leads to a formula very similar to the one which we just motivated.

Lemma 1.2.14 ([BF71, §29, Formula (3), p. 41ff])
Let $A$ be a polytope with facets $F_{1}, F_{2}, \ldots, F_{N}, K$ a convex body in $V$ and $k$ an arbitrary interior point of $K$. Suppose $f^{i}$ is the outer unit normal to the facet $F_{i}$ and $h_{K}$ is the support function of $K$. The mixed volume of $\underbrace{A, A, \ldots, A}_{(n-1)-\mathrm{times}}$ and $K$ is given by

$$
\begin{equation*}
\mathrm{V}(A[n-1], K)=\frac{1}{n} \sum_{i=1}^{N} \mathcal{H}^{n-1}\left(F_{i}\right) h_{K-k}\left(f^{i}\right) \tag{1.2.21}
\end{equation*}
$$

Proof: Since mixed volume is invariant under translation of any of its arguments, the left hand side $\mathrm{V}(A[n-1], K)$ equals $\mathrm{V}(A[n-1], K-k)$. Without loss of generality, we thereby may assume $K$ contains the origin of $V$ and $k=0$. Further suppose $n \geq 2$, for if $n=1$ then $A$ and $K=\left[k_{1}, k_{2}\right]$ are segments in the one-dimensional vector space $V$. In addition, $N=2$ because the number of vertices is minimal by definition of a polytope and, in fact, the facets of $A$ coincide with the vertices which are the endpoints of the segment. The value $h_{K}\left(f^{i}\right)$ gives the distance of $0 \in K$ to its endpoints (Proposition 1.2.6) and because the 0-dimensional Hausdorff measure counts the number of points in a set, $\mathcal{H}^{0}\left(F_{i}\right)=1$. Thus, $\frac{1}{1} \sum_{i=1}^{2} \mathcal{H}^{0}\left(F_{i}\right) h_{K}\left(f^{i}\right)=\sum_{i=1}^{2} h_{K}\left(f^{i}\right)=\left\|k_{1}\right\|_{V}+\left\|k_{2}\right\|_{V}=\left\|k_{1}-k_{2}\right\|_{V}=$ $\mathcal{H}^{1}(K)=\mathrm{V}(A[0], K)$. The third equality holds because the segment $\left[k_{1}, k_{2}\right]$ contains the origin, so $k_{1}=-\left|k_{1}\right| e_{1}$ and $k_{2}=\left|k_{2}\right| e_{1}$ point in opposite directions.
By Corollary 1.2.11 the mixed volume in the assertation can be calculated through

$$
\mathrm{V}(A[n-1], K)=\frac{1}{n} \lim _{t \rightarrow 0} \frac{\mathcal{H}^{n}(A+t K)-\mathcal{H}^{n}(A)}{t}
$$

So, a further examination of $\mathcal{H}^{n}(A+t K)$ is needed. We will show that the set $A+t K$ can be subdivided into several smaller pieces whose volumes are easier to calculate. According to Proposition 1.2.6 the hyperplane $H_{i}:=\left\{y \in V \mid f^{i}(y)=h_{K}\left(f^{i}\right)\right\}$ supports $K$, whereby there is $x_{i} \in H_{i} \cap K$ for each $i=1,2, \ldots, N$. Define the set $P_{i}:=\left\{y+s x_{i} \mid y \in F_{i}, 0 \leq s \leq t\right\}=F_{i}+\left[0, t x_{i}\right]$. This set describes a prism with $(n-1)$-dimensional base $F_{i}$ which is translated into direction $x_{i}$ (see Figure 1.2 (b)). Its "thickness" or height is

$$
h:=h_{A+t K}\left(f^{i}\right)-h_{A}\left(f^{i}\right)=h_{A}\left(f^{i}\right)+t h_{K}\left(f^{i}\right)-h_{A}\left(f^{i}\right)=t h_{K}\left(f^{i}\right)
$$

where we used Proposition 1.2.6 twice to calculate the distances of the origin to the supporting hyperplanes for $A+t K$ and $A$ defined by the normal $f^{i}$. The hyperplanes $H_{s}\left(f^{i}\right)=\left\{y \in V \mid f^{i}(y)=s\right\}$ are parallel to the facet $F_{i}$. Further, the intersection of a prism $P_{i}$ with these hyperplanes is either


Figure 1.2: Calculation of the mixed volume V $(A[n-1], K)$
empty or a translation of the facet $F_{i}$ into direction $x_{i}$. In particular,

$$
P_{i} \cap H_{s}\left(f^{i}\right)= \begin{cases}F_{i}+v_{s} & , s \in\left[h_{A}\left(f^{i}\right), h_{A+t K}\left(f^{i}\right)\right] \\ \varnothing & , s \notin\left[h_{A}\left(f^{i}\right), h_{A+t K}\left(f^{i}\right)\right] .\end{cases}
$$

where $v_{s} \in\left[0, t x_{i}\right]$ is a vector in direction $x_{i}$ (of length $\frac{s-h_{A}\left(f^{i}\right)}{h_{K}\left(f^{i}\right)}\left\|x_{i}\right\|_{V}$ ). Then by Cavalieri's principle ([For09, §5, Satz 3, p. 49]) we have

$$
\begin{aligned}
\mathcal{H}^{n}\left(P_{i}\right) & =\int_{\mathbb{R}} \mathcal{H}^{n-1}\left(P_{i} \cap H_{s}\left(f^{i}\right)\right) \mathrm{d} s=\int_{h_{A}\left(f^{i}\right)}^{h_{A+t K}\left(f^{i}\right)} \mathcal{H}^{n-1}\left(F_{i}+v_{s}\right) \mathrm{d} s \\
& =\int_{h_{A}\left(f^{i}\right)}^{h_{A}\left(f^{i}\right)+h} \mathcal{H}^{n-1}\left(F_{i}\right) \mathrm{d} s=h \mathcal{H}^{n-1}\left(F_{i}\right) \\
& =t \mathcal{H}^{n-1}\left(F_{i}\right) h_{K}\left(f^{i}\right) .
\end{aligned}
$$

Each of these $N$ prisms $P_{i}$ is contained in the sum $A+t K$ (by definition of the Minkowski sum) but not in $A$ (by construction). Therefore,

$$
A+t K=A \cup \bigcup_{i=1}^{N} P_{i} \cup S
$$

where $S:=(A+t K) \backslash\left(A \cup \bigcup_{i=1}^{N} P_{i}\right)$ is the remaining part of the Minkowski sum "close to" the ( $n-2$ )-dimensional faces of $A$ (dark gray in Figure $1.2(\mathrm{~b})$ ). Note that this decomposition of $A+t K$
is not pairwise disjoint but any intersection of the sets $A, P_{i}$ and $S$ is lower-dimensional and thus a set of $n$-dimensional Hausdorff measure zero. Therefore, we conclude

$$
\begin{aligned}
\mathrm{V}(A[n-1], K) & =\frac{1}{n} \lim _{t \rightarrow 0} \frac{\mathcal{H}^{n}(A+t K)-\mathcal{H}^{n}(A)}{t} \\
& =\frac{1}{n} \lim _{t \rightarrow 0} \frac{\mathcal{H}^{n}(A)+\sum_{i=1}^{N} \mathcal{H}^{n}\left(P_{i}\right)+\mathcal{H}^{n}(S)-\mathcal{H}^{n}(A)}{t} \\
& =\frac{1}{n} \lim _{t \rightarrow 0} \frac{\sum_{i=1}^{N} \mathcal{H}^{n-1}\left(F_{i}\right) t h_{K}\left(f^{i}\right)+\mathcal{H}^{n}(S)}{t} \\
& =\frac{1}{n} \sum_{i=1}^{N} \mathcal{H}^{n-1}\left(F_{i}\right) h_{K}\left(f^{i}\right)+\frac{1}{n} \lim _{t \rightarrow 0} \frac{\mathcal{H}^{n}(S)}{t} .
\end{aligned}
$$

We need to further investigate the term $\lim _{t \rightarrow 0} \frac{\mathcal{H}^{n}(S)}{t}$.

The intersection of a facet $F_{i}$ with another facet $F_{j}$ is either empty or a face of $A$ of dimension $j \leq n-2$ ([Brø83, Theorem 5.9, pp. 30+33 and Theorem 10.4, p. 65]). Recall that a face of a polytope again is a polytope and therefore compact by Lemma 1.2.3. Then the set

$$
\mathcal{F}:=\mathcal{F}^{n-2}(A):=\bigcup_{\substack{i, j=1 \\ i \neq j}}^{N}\left(F_{i} \cap F_{j}\right)
$$

is the union of all those faces of $A$ which are of dimension $j \leq n-2$. (In the two-dimensional drawing Figure 1.2 the set $\mathcal{F}$ is represented by the "vertices" which actually consist of $(n-2)$-dimensional faces of $A$.) The set $\mathcal{F}$ is compact since it is the finite union of compact sets.

Define the radius $R:=\sup _{\|f\|_{V^{*}}=1} h_{K}(f)$ and recall from Proposition 1.2.6 that for each unit length $f \in V^{*}$ the hyperplane $H_{h_{K}(f)}(f)=\left\{v \in V \mid f(v)=h_{K}(f)\right\}$ supports $K$, that is, $H_{h_{K}(f)}(f) \cap K \neq \varnothing$ and the value $h_{K}(f)$ is the distance of $0 \in K$ to this hyperplane. Geometrically, we may think of

$$
R=\sup _{\|f\|_{V^{*}}=1} h_{K}(f)=\sup \left\{\operatorname{dist}\left(0, H_{h_{K}(f)}(f)\right) \mid f \in V^{*},\|f\|_{V^{*}}=1\right\}
$$

as the supremal distance of $0 \in K$ to the boundary of $K$ or, in other words, the radius of the smallest ball which contains $K$ (Figure 1.2). That is,

$$
A \cup \bigcup_{i=1}^{N} P_{i} \cup S=A+t K \subset A+t B_{R}(0)=\left(\operatorname{int}(A) \cup \bigcup_{i=1}^{N} F_{i} \cup \mathcal{F}\right)+t B_{R}(0)
$$

However, since by definition $S \not \subset\left(A \cup \bigcup_{i=1}^{N} P_{i}\right)$, we know that actually

$$
S \subset \mathcal{F}+t B_{R}(0)=\bigcup_{c \in \mathcal{F}} B_{t R}(c)
$$

Since the closure $\bar{S}$ is compact (because its bounded as a subset of the bounded convex body $A+t K$ ) we may take a subcover of balls of radius $t R$ centered at finitely many, say $M=M(K)$, points $c_{i} \in \mathcal{F}$. Therefore, the volume of $S$ calculates to

$$
\mathcal{H}^{n}(S)=\mathcal{H}^{n}(\bar{S}) \leq \sum_{i=1}^{M} \mathcal{H}^{n}\left(B_{t R}\left(c_{i}\right)\right)=M \mathcal{H}^{n}\left(B_{1}(0)\right) R^{n} t^{n}=C t^{n}
$$

where $C=C(A, K, n)>0$ is a constant. Since $n \geq 2$, we conclude

$$
0 \leq \lim _{t \rightarrow 0} \frac{\mathcal{H}^{n}(S)}{t} \leq C \lim _{t \rightarrow 0} t^{n-1}=0
$$

which finishes the proof.

Note that if $K=A$ the previous result produces the formula for the volume of a polytope which was motivated earlier. Thus, a posteriori, we have also given a (rather complicated) proof of the usual volume formula for polytopes.

For the next chapter we need to prove the following result about polytopes in a two-dimensional space which intuitively holds true. Indeed, the proof is quite straightforward if one unravels all the tedious definitions that will be made.

## Lemma 1.2.15

Let $P=\left[a_{1} a_{2} \ldots a_{2 N}\right]$ be a symmetric polygon in a two-dimensional normed space $V$ with $2 N$ vertices. Further, let $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{N}$ be arbitrary positive real numbers. Then there exists a symmetric polygon $Q=\left[b_{1} b_{2} \ldots b_{2 N}\right]$ with $2 N$ vertices such that

$$
b_{i+1}-b_{i}=\lambda_{i}\left(a_{i+1}-a_{i}\right)
$$

for all $i=1,2, \ldots, N$, that is, the ith edge of $P$ is scaled by the factor $\lambda_{i}$ and therefore the $i$ th edges of $P$ and $Q$ are parallel.
Proof: The polygon $Q$ can be explicitly constructed from $P$. We will prove this lemma by induction on $N$. When $N=1$, the statement is trivial because we can simply choose $b_{1}:=\lambda_{1} a_{1}$ and $b_{2}:=-b_{1}$. Suppose the assertion holds for an $N \in \mathbb{N}$ and let $P=\left[a_{1} a_{2} \ldots a_{2(N+1)}\right]$ be a symmetric polygon with


Figure 1.3: Constructing the scaled polygon $Q$ from $P$ for $2(N+1)$ vertices
$2(N+1)$ vertices. For $i=1,2, \ldots, N+1$ define the vectors $v_{i}:=a_{i+1}-a_{i}$ corresponding to the edges of $P$ and $v_{i}^{\prime}:=\lambda_{i} v_{i}$ which will become the edges of the polygon $Q$. To be able to apply the inductive hypothesis we need to construct a symmetric polygon with $2 N$ vertices (Figure 1.3 (a)). We will do this by removing the vertices $a_{1}$ and $a_{N+2}=-a_{1}$, that is, define the symmetric $2 N$-gon $\widetilde{P}=\left[\widetilde{a}_{1} \widetilde{a}_{2} \ldots \widetilde{a}_{2 N}\right]$ where

$$
\widetilde{a}_{i}:=\left\{\begin{array}{ll}
a_{i+1} & i=1,2, \ldots, N \\
a_{i+2} & i=N+1, N+2, \ldots, 2 N
\end{array} .\right.
$$

Further, define the new scaling coefficients

$$
\widetilde{\lambda}_{i}:= \begin{cases}\lambda_{i+1} & i=1,2, \ldots, N-1 \\ 1 & i=N\end{cases}
$$

that is, the edge $\widetilde{v}_{N}:=\widetilde{a}_{N+1}-\widetilde{a}_{N}$ will not be scaled (Figure $1.3(\mathrm{~b})$ ). By the inductive hypothesis there exists a symmetric $2 N$-gon $\widetilde{Q}=\left[\widetilde{b}_{1} \widetilde{b}_{2} \ldots \widetilde{b}_{2 N}\right]$ such that

$$
\widetilde{b}_{i+1}-\widetilde{b}_{i}=\widetilde{\lambda}_{i}\left(\widetilde{a}_{i+1}-\widetilde{a}_{i}\right)
$$

for all $i=1,2, \ldots, N$ (Figure $1.3(\mathrm{c})$ ). Consider only the set $\left\{\widetilde{b}_{1}, \widetilde{b}_{2}, \ldots, \widetilde{b}_{N}\right\}$ of the first $N$ vertices. For notational convenience let us forget about the remaining vertices and (re-)define two new vertices by going along the initial scaled edges, namely, $\widetilde{b}_{0}:=\widetilde{b}_{1}-v_{1}^{\prime}$ and $\widetilde{b}_{N+1}:=\widetilde{b}_{N}+v_{N+1}^{\prime}$ (Figure 1.3 (d)). Translate the extended vertex set $\left\{\widetilde{b}_{0}, \widetilde{b}_{0}, \ldots, \widetilde{b}_{N+1}\right\}$ affinely such that $\widetilde{b}_{N+1}=-\widetilde{b}_{0}$ and denote the points by the same name (Figure 1.3 (e)). In fact, the map $T: V \rightarrow V, v \mapsto v-\frac{1}{2}\left(\widetilde{b}_{N+1}+\widetilde{b}_{0}\right)$ is such an affine translation. Finally, we define the vertices

$$
b_{i}:=\left\{\begin{array}{ll}
\widetilde{b}_{i-1} & i=1,2, \ldots, N+1 \\
-\widetilde{b}_{i-1-(N+1)} & i=N+2, N+3, \ldots, 2 N+2
\end{array} .\right.
$$

The $2(N+1)$-gon $Q=\left[b_{1} b_{2} \ldots b_{2 N+2}\right]$ is symmetric by construction (Figure 1.3 (f)). Unwinding the definitions, we compute

$$
b_{i+1}-b_{i}=\widetilde{b}_{i}-\widetilde{b}_{i-1}=\widetilde{\lambda}_{i-1}\left(\widetilde{a}_{i}-\widetilde{a}_{i-1}\right)=\lambda_{i}\left(a_{i+1}-a_{i}\right)
$$

for $i=2,3, \ldots, N+1$. Moreover, for $i=1$

$$
b_{2}-b_{1}=\widetilde{b}_{1}-\widetilde{b}_{0}=v_{1}^{\prime}=\lambda_{1} v_{1}=\lambda_{1}\left(a_{2}-a_{1}\right)
$$

### 1.2.6 Convex and linear functions on convex sets

This section ends with some results on the action of linear and convex functions on convex sets.

## Lemma 1.2.16

Let $V_{1}$ and $V_{2}$ be two finite-dimensional normed spaces and $L: V_{1} \rightarrow V_{2}$ a linear map.
(i) If $K_{1} \subset V_{1}$ is a convex body then $L\left(K_{1}\right) \subset V_{2}$ is a convex body.
(ii) If $K_{2} \subset \operatorname{im}(L) \subset V_{2}$ is a convex body then the preimage $L^{-1}\left(K_{2}\right) \subset V_{1}$ is a convex set. If further $L$ is injective then $L^{-1}\left(K_{2}\right)$ is a convex body.
(iii) Suppose $P_{2} \subset \operatorname{im}(L) \subset V_{2}$ is a polytope containing the origin given by its half-space representation (Theorem 1.2.7)

$$
P_{2}=\bigcap_{i=1}^{M}\left\{y \in V_{2} \mid F^{i}(y) \leq 1\right\}
$$

for $F^{i} \in V_{2}^{*}$. Then the preimage $P_{1}:=L^{-1}\left(P_{2}\right) \subset V_{1}$ is a polytope and its half-space representation is given by

$$
P_{1}=\bigcap_{i=1}^{M}\left\{x \in V_{1} \mid f^{i}(x) \leq 1\right\}
$$

where $f^{i}:=F^{i} \circ L \in V_{1}^{*}$.
(iv) If $K_{2} \subset V_{2}$ is symmetric then the preimage $L^{-1}\left(K_{2}\right) \subset V_{1}$ is a symmetric set.

Proof: For $i=1,2$ let $x_{i} \in K_{1}$ and $y_{i} \in K_{2}$ such that $y_{i}=L x_{i}$ and let $t \in[0,1]$. Then by the linearity of $L$

$$
t y_{1}+(1-t) y_{2}=t L x_{1}+(1-t) L x_{2}=L\left(t x_{1}+(1-t) x_{2}\right) .
$$

If $K_{1}$ is convex this equality implies the convexity of $L\left(K_{1}\right)$ if read from right to left and conversely, if read from left to right it implies the convexity of $L^{-1}\left(K_{2}\right)$ if $K_{2}$ is convex. If $x_{0}$ is an interior point of $K_{1}$ there is $\varepsilon>0$ such that $B_{1}:=B_{\varepsilon\|L\|_{\mathcal{L}\left(V_{1}, V_{2}\right)}^{-1}}\left(x_{0}\right) \subset K_{1}$. Let $x$ be an arbitrary point in $B_{1}$. Then
$L x \in L\left(K_{1}\right)$ but also $\left\|L x-L x_{0}\right\|_{V_{2}} \leq\|L\|_{\mathcal{L}\left(V_{1}, V_{2}\right)}\left\|x-x_{0}\right\|_{V_{2}}<\|L\|_{\mathcal{L}\left(V_{1}, V_{2}\right)} \varepsilon\|L\|_{\mathcal{L}\left(V_{1}, V_{2}\right)}^{-1}=\varepsilon$, that is, $B_{\varepsilon}\left(L x_{0}\right) \subset L\left(K_{1}\right)$. In other words, $L\left(K_{1}\right)$ has non-empty interior. Since a linear map is continuous, the image of a compact set is compact. Therefore, $L\left(K_{1}\right)$ is a convex body if $K_{1}$ is.

For part (ii) note that $\widetilde{L}: V_{1} \rightarrow \operatorname{im}(L), v \mapsto L v$ is bijective and $L^{-1}\left(K_{2}\right)=\widetilde{L}^{-1}\left(K_{2} \cap \operatorname{im}(L)\right)$. Thus, we just interchange the roles of $K_{1}, V_{1}$ and $K_{2}, V_{2}$, respectively, and apply ( $i$ ) to the linear map $\widetilde{L}^{-1}: V_{2} \rightarrow V_{1}$.

For part (iii) denote $Q:=\bigcap_{i=1}^{M}\left\{x \in V_{1} \mid f^{i}(x) \leq 1\right\}$ and let $v \in P_{1}$. Then $L v \in P_{2}$, that is, $1 \geq F^{i}(L v)=f^{i}(v)$ for $i=1,2, \ldots, M$. Thus, $v \in Q$. Conversely, if $v \in Q$ then $f^{i}(v)=F^{i}(L v) \leq 1$ for $i=1,2, \ldots, M$, that is, $L v \in P_{2}$. Hence, $v \in P_{1}$.

For part (iv) suppose $K_{2}$ is symmetric and let $x \in L^{-1}\left(K_{2}\right)$, that is, $L x \in K_{2}$. Since $K_{2}$ is symmetric, we know that $L(-x)=-(L x) \in K_{2}$ or $-x \in L^{-1}\left(K_{2}\right)$. So, $L^{-1}\left(K_{2}\right)$ is symmetric.

Let us conclude this section with a result on the maxima of convex functions on polytopes. Suppose $A \subset V$ is a convex set. A map $F: A \rightarrow \mathbb{R} \cup\{-\infty,+\infty\}$ is called a convex function if

$$
F(\lambda v+(1-\lambda) w) \leq \lambda F(v)+(1-\lambda) F(w)
$$

whenever $\lambda \in[0,1]$ and $v, w \in A$. We call a convex function $F$ proper if it nowhere takes the value $-\infty$ and is not identically $+\infty$.

An easy induction proves that the inequality above holds for arbitrary finite convex combinations.

## Lemma 1.2.17

Let $A \in V$ be a convex set, $v_{1}, v_{2}, \ldots, v_{N} \in A, \lambda_{1}, \lambda_{2}, \ldots, \lambda_{N} \in[0,1]$ with $\sum_{i=1}^{N} \lambda_{i}=1$ and $F: A \rightarrow \mathbb{R}$ a convex function. Then

$$
F\left(\sum_{i=1}^{N} \lambda_{i} v_{i}\right) \leq \sum_{i=1}^{N} \lambda_{i} F\left(v_{i}\right) .
$$

Proof: For $N=1$ the claim is trivial and for $N=2$ it follows from the definition of a convex function. Assume the assertion holds for $N-1$. Note that $\sum_{i=1}^{N-1} \lambda_{i}=1-\lambda_{N}$ or $\sum_{i=1}^{N-1} \frac{\lambda_{i}}{1-\lambda_{N}}=1$. Define $w:=\sum_{i=1}^{N-1} \frac{\lambda_{i}}{1-\lambda_{N}} v_{i} \in V$. Then

$$
\begin{aligned}
F\left(\sum_{i=1}^{N} \lambda_{i} v_{i}\right) & =F\left(\sum_{i=1}^{N-1} \lambda_{i} v_{i}+\lambda_{N}\right)=F\left(\left(1-\lambda_{N}\right) w+\lambda v_{N} v_{N}\right) \\
& \leq\left(1-\lambda_{N}\right) F(w)+\lambda F\left(v_{N}\right) \\
& =\left(1-\lambda_{N}\right) F\left(\sum_{i=1}^{N-1} \frac{\lambda_{i}}{1-\lambda_{N}} v_{i}\right)+\lambda F\left(v_{N}\right)
\end{aligned}
$$

$$
\begin{aligned}
& \leq\left(1-\lambda_{N}\right) \sum_{i=1}^{N-1} \frac{\lambda_{i}}{1-\lambda_{N}} F\left(v_{i}\right)+\lambda F\left(v_{N}\right) \\
& =\sum_{i=1}^{N} \lambda_{i} F\left(v_{i}\right)
\end{aligned}
$$

where we used the definition of a convex function for the first inequality and the inductive hypothesis for the second.

The following important maximum principle for convex functions holds.

## Proposition 1.2.18

Let $P=\left[a_{1} a_{2} \ldots a_{N}\right]$ be a polytope in $V$ and $F: V \rightarrow \mathbb{R}$ a convex function which is bounded above on $P$. Then the supremum of $F$ on $P$ is attained at one of its vertices, that is,

$$
\sup _{x \in P} F(x)=\sup _{i=1,2, \ldots, N} F\left(a_{i}\right)=\max _{i=1,2, \ldots, N} F\left(a_{i}\right) .
$$

Proof: The second identity follows from the fact that $\left\{a_{1}, a_{2}, \ldots, a_{N}\right\}$ is a finite set. It is immediate that $\sup _{i=1,2, \ldots, N} F\left(a_{i}\right) \leq \sup _{x \in P} F(x)$ since $\left\{a_{1}, a_{2}, \ldots, a_{N}\right\} \subset \operatorname{conv}\left(a_{1}, a_{2}, \ldots, a_{N}\right)=P$. Conversely, let $v=\sum_{i=1}^{N} \lambda_{i} a_{i} \in P$ where $\sum_{i=1}^{N} \lambda_{i}=1$. Then the convexity of $F$ yields

$$
F(v)=F\left(\sum_{i=1}^{N} \lambda_{i} a_{i}\right) \leq \sum_{i=1}^{N} \lambda_{i} F\left(a_{i}\right) \leq \sum_{i=1}^{N} \lambda_{i} \max _{k=1,2, \ldots, N} F\left(a_{k}\right)=\max _{k=1,2, \ldots, N} F\left(a_{k}\right)
$$

since $v$ was arbitrary we get $\sup _{x \in P} F(x) \leq \max _{k=1,2, \ldots, N} F\left(a_{k}\right)$.

In fact, this result holds true if $P$ is merely convex and not necessarily a polytope but then the supremum is attained on the set of extreme points (the extreme subsets that consist only of one point), see [Roc70, Corollary 32.3.1, p. 344]. For the application in the next chapter only the special case above is needed.

### 1.3 Differential geometry

In this section we introduce basic concepts from differential geometry, such as the notion of a smooth manifold and its tangent space and tangent bundle. We describe the differential of a smooth map and examine smooth immersions and images of such maps. Finally, we consider Finsler metrics. In this section and for the rest of this thesis we will use the Einstein summation convention, which means that - if not explicitly stated otherwise - we sum over repeated upper and lower indices. We use the convention that Latin indices run from 1 to $n$ and Greek indices run from 1 to $m$. The presentation given here is borrowed from [Ove13, Chapters 1.3-1.4, pp. 28-41, Chapter 2] which in turn is based on the book of Lee [Lee13]. The source for the concepts of Finsler geometry is the textbook [BCS00] by Bao et al.

### 1.3.1 Smooth manifolds

Suppose $\mathcal{M}$ is a topological space. We call $\mathcal{M}$ a topological manifold of dimension $\boldsymbol{m}$ or a topological m-manifold if it fulfils the following three properties:

- $\mathcal{M}$ is a Hausdorff space: for every pair of distinct points $p, q \in \mathcal{M}$, there are disjoint open subsets $U, V \subset \mathcal{M}$ such that $p \in U$ and $q \in V$.
- $\mathcal{M}$ is second-countable: there exists a countable basis for the topology of $\mathcal{M}$.
- $\mathcal{M}$ is locally Euclidean of dimension $\boldsymbol{m}$ : for each point $p \in \mathcal{M}$ there is an open subset $U \subset \mathcal{M}$ containing $p$, an open subset $\widehat{U} \subset \mathbb{R}^{m}$ and a homeomorphism $\varphi: U \rightarrow \widehat{U}$.

A (local) coordinate chart is a pair $(U, \varphi)$ where $U$ is an open subset of $\mathcal{M}$ and $\varphi: U \rightarrow \widehat{U}$ is a homeomorphism onto an open subset $\widehat{U}=\varphi(U) \subset \mathbb{R}^{m}$. The set $U$ is called a coordinate neighbourhood and the map $\varphi$ a (local) coordinate map. The component functions $\left(u^{1}, u^{2}, \ldots, u^{m}\right)$ of $\varphi$, defined by $\varphi(p)=\left(u^{1}(p), u^{2}(p), \ldots, u^{m}(p)\right)$ ), are called (local) coordinates. Sometimes we write $\widehat{p}=\left(u^{1}(p), u^{2}(p), \ldots, u^{m}(p)\right)$ for the local coordinates of a point $p \in \mathcal{M}$ and we use any of the notations $(U, \varphi)=\left(U,\left(u^{1}, u^{2}, \ldots, u^{m}\right)\right)=\left(U,\left(u^{\alpha}\right)\right)$ for a coordinate chart (see [Lee13, pp. 2-4]). Another way to emphasise the dimension $\operatorname{dim}(\mathcal{M})$ of a manifold is to write it as an upper index, e.g. $\mathcal{M}^{m}$.

To make sense of derivatives of functions on a manifold or maps between manifolds we need to introduce an additional structure which turns a topological into a smooth manifold. Consider two coordinate charts $(U, \varphi)$ and $(V, \psi)$ of a topological manifold $\mathcal{M}$ where $U \cap V \neq \varnothing$. The transition
map between $\varphi$ and $\psi$ is the mapping

$$
\psi \circ \varphi^{-1}: \varphi(U \cap V) \subset \mathbb{R}^{m} \rightarrow \psi(U \cap V) \subset \mathbb{R}^{m}
$$

A smooth atlas $\mathcal{A}$ is the union of all those coordinate charts of $\mathcal{M}$ such that the union of the coordinate neighbourhoods cover $\mathcal{M}$ and each transition map is a smooth map from $\mathbb{R}^{m}$ to $\mathbb{R}^{m}$, that is, $\mathcal{A}$ fulfils the following two properties:

- $\mathcal{M}=\bigcup_{(U, \varphi) \in \mathcal{A}} U$,
- $\psi \circ \varphi^{-1}: \varphi(U \cap V) \rightarrow \psi(U \cap V)$ is of class $C^{\infty}\left(\mathbb{R}^{m}, \mathbb{R}^{m}\right)$ for all coordinate charts $(U, \varphi),(V, \psi) \in \mathcal{A}$ with $U \cap V \neq \varnothing$.

Every smooth atlas $\mathcal{A}$ of $\mathcal{M}$ is contained in a unique maximal smooth atlas ([Lee13, Proposition 1.17, p. 13]) - an atlas containing $\mathcal{A}$ which is not properly contained in any larger smooth atlas itself. A smooth structure on a topological manifold $\mathcal{M}$ is a maximal smooth atlas. A smooth manifold is a topological manifold endowed with a smooth structure. A coordinate chart of a smooth manifold is called a smooth coordinate chart.

Let $\mathcal{M}^{m}$ and $\mathcal{N}^{n}$ be smooth manifolds. A map $F: \mathcal{M} \rightarrow \mathcal{N}$ is called a smooth map if for every $p \in \mathcal{M}$ there are smooth coordinate charts $(U, \varphi)$ containing $p$ and $(V, \psi)$ containing $F(p)$ such that $F(U) \subset V$ and

$$
\widehat{F}:=\psi \circ F \circ \varphi^{-1} \in C^{\infty}(\varphi(U), \psi(V))
$$

where $\varphi(U)$ is an open subset of $\mathbb{R}^{m}$ and $\psi(V)$ is an open subset of $\mathbb{R}^{n}$. The function $\widehat{F}$ is called a coordinate representation of $\boldsymbol{F}$. Note that this definition does not depend on the choice of coordinate charts because the transition maps are smooth. The set of smooth maps between manifolds will be denoted by $C^{\infty}(\mathcal{M}, \mathcal{N})$. In the special case when $\mathcal{N}=\mathbb{R}$ we set

$$
C^{\infty}(\mathcal{M}):=C^{\infty}(\mathcal{M}, \mathbb{R})
$$

A homeomorphism from $\mathcal{M}$ to $\mathcal{N}$ is a continuous bijective map $F: \mathcal{M} \rightarrow \mathcal{N}$ that has a continuous inverse. A smooth bijective map $F: \mathcal{M} \rightarrow \mathcal{N}$ that has a smooth inverse is a diffeomorphism from

## $\mathcal{M}$ to $\boldsymbol{N}$.

Now we will introduce a "local linear approximation" of a given smooth $m$-dimensional manifold $\mathcal{M}$. Let $p$ be a point of $\mathcal{M}$. A linear map $v: C^{\infty}(\mathcal{M}) \rightarrow \mathbb{R}$ is called a derivation at $\boldsymbol{p}$ if it satisfies the
following product rule

$$
\begin{equation*}
v(f g)=f(p) v(g)+g(p) v(f) \tag{1.3.22}
\end{equation*}
$$

for all $f, g \in C^{\infty}(\mathcal{M})$. The set of all derivations at $p$ is called the tangent space at $\boldsymbol{p}$ and is denoted by $T_{p} \mathcal{M}$. An element of $T_{p} \mathcal{M}$ is called a tangent vector at $\boldsymbol{p}$. It can be shown that $T_{p} \mathcal{M}$ is an $m$-dimensional real vector space. A basis of $T_{p} \mathcal{M}$ can be given through a smooth coordinate chart $(U, \varphi)=\left(U,\left(u^{\alpha}\right)\right)$ containing $p$ by defining the derivations $\left.\frac{\partial}{\partial u^{\alpha}}\right|_{p}$ as

$$
\left.\frac{\partial}{\partial u^{\alpha}}\right|_{p}(f):=\left.\frac{\partial}{\partial u^{\alpha}}\right|_{\varphi(p)}\left(f \circ \varphi^{-1}\right)
$$

for $f \in C^{\infty}(U)$ and $\alpha=1,2, \ldots, m$ (see [Lee13, Proposition 3.15, p. 61]). Note that on the right-hand side the symbol $\frac{\partial}{\partial u^{\alpha}}$ means the partial derivative with respect to the coordinate $u^{\alpha}$ of a real-valued function, whereas on the left-hand side it denotes a derivation. These basis tangent vectors are called coordinate vectors and form the coordinate basis for $\boldsymbol{T}_{\boldsymbol{p}} \mathcal{M}$. Thus, any tangent vector $v \in T_{p} \mathcal{M}$ can be written uniquely as

$$
v=\left.v^{\alpha} \frac{\partial}{\partial u^{\alpha}}\right|_{p}
$$

where $v^{\alpha} \in \mathbb{R}$ and the summation convention is used if an upper index "in the denominator" is understood as a lower index. That is, the right-hand side is understood to be summed over $\alpha=$ $1,2, \ldots, m$.

The tangent bundle of $\mathcal{M}$ is the disjoint union of all tangent spaces $T_{p} \mathcal{M}$ to $\mathcal{M}$, that is,

$$
T \mathcal{M}:=\bigsqcup_{p \in \mathcal{M}} T_{p} \mathcal{M}
$$

Here the disjoint union $\bigsqcup_{p \in \mathcal{M}} T_{p} \mathcal{M}$ is the set $\bigcup_{p \in \mathcal{M}}\{p\} \times T_{p} \mathcal{M}$. Thus, an element of $T \mathcal{M}$ is of the form $(p, v)$ where $v \in T_{p} \mathcal{M}$. The tangent bundle of a smooth $m$-dimensional manifold is itself a smooth manifold of dimension $2 m$. A local smooth coordinate chart $(U, \varphi)=\left(U,\left(u^{\alpha}\right)\right)$ of $\mathcal{M}$ extends to a smooth coordinate chart of $T \mathcal{M}$ in the following way. Let $\pi: T \mathcal{M} \rightarrow \mathcal{M}$ be the natural projection map defined by $\pi(p, v):=p$. The set $\pi^{-1}(U) \subset T \mathcal{M}$ is the set of all tangent vectors of $\mathcal{M}$ at points of $U$. Define $\widetilde{\varphi}: \pi^{-1}(U) \rightarrow \mathbb{R}^{2 m}$ by

$$
\begin{equation*}
\widetilde{\varphi}\left(p,\left.v^{\alpha} \frac{\partial}{\partial u^{\alpha}}\right|_{p}\right):=\left(u^{1}(p), u^{2}(p), \ldots, u^{m}(p), v^{1}, v^{2}, \ldots, v^{m}\right) . \tag{1.3.23}
\end{equation*}
$$

Its image set is $\varphi\left(U_{p}\right) \times \mathbb{R}^{m}$. It can be shown that this map is a smooth coordinate chart and all such maps form a smooth atlas for $T \mathcal{M}$, turning it into a smooth manifold (see [Lee13, Proposition 3.18, p. 66]). The local coordinates given by (1.3.23) are called natural coordinates on $\boldsymbol{T} \boldsymbol{\mathcal { M }}$.

A smooth map $F: \mathcal{M} \rightarrow \mathcal{N}$ between two smooth manifolds gives rise to a linear map between their respective tangent spaces. The differential of $\boldsymbol{F}$ at $\boldsymbol{p}$ is the linear map

$$
d F_{p}: T_{p} \mathcal{M} \rightarrow T_{F(p)} \mathcal{N}
$$

which for a fixed $v \in T_{p} \mathcal{M}$ is given by its action on $f \in C^{\infty}(\mathcal{M})$ by

$$
d F_{p}(v)(f):=v(f \circ F)
$$

If one considers the differential at each point of $\mathcal{M}$ this map extends to a map between the tangent bundles, called the global differential of $\boldsymbol{F}$

$$
d F: T \mathcal{M} \rightarrow T \mathcal{N},(p, v) \mapsto\left(F(p), d F_{p}(v)\right)
$$

If $(U, \varphi)=\left(U,\left(u^{\alpha}\right)\right)$ is a smooth coordinate chart of $\mathcal{M}$ and $(V, \psi)=\left(V,\left(x^{i}\right)\right)$ a smooth coordinate chart of $\mathcal{N}$ then the differential of $F$ at $p$ can be locally expressed by

$$
\begin{equation*}
d F_{p}\left(\left.v^{\alpha} \frac{\partial}{\partial u^{\alpha}}\right|_{p}\right)=\left.v^{\alpha} \frac{\partial \widehat{F}^{i}}{\partial u^{\alpha}}(\widehat{p}) \frac{\partial}{\partial x^{i}}\right|_{F(p)} \tag{1.3.24}
\end{equation*}
$$

where $\widehat{F}:=\psi \circ F \circ \varphi^{-1}$ and $\widehat{p}=\varphi(p)$ are the coordinate representations of $F$ and $p$ (see [Lee13, Equation (3.10), p. 63]). Here latin indices are summed over $i=1,2, \ldots, n$ and greek indices are summed over $\alpha=1,2, \ldots, m$.

A smooth immersion is a smooth map $X: \mathcal{M} \rightarrow \mathcal{N}$ whose differential $d X_{p}: T_{p} \mathcal{M} \rightarrow T_{X(p)} \mathcal{N}$ is injective at each point $p$ in $\mathcal{M}$. Equivalently, a smooth map $X: \mathcal{M} \rightarrow \mathcal{N}$ is a smooth immersion if and only if $\operatorname{dim}\left(\operatorname{im}\left(d X_{p}\right)\right)=m$ for all $p \in \mathcal{M}$. A smooth embedding of $\boldsymbol{\mathcal { M }}$ into $\boldsymbol{\mathcal { N }}$ is a smooth immersion $X: \mathcal{M} \rightarrow \mathcal{N}$ which also is a homeomorphism onto its image $X(\mathcal{M})$.

Suppose $X: \mathcal{M}^{m} \rightarrow \mathcal{N}^{n}$ is a smooth immersion. One can show that locally $X$ is an embedding, that is, for each point $p \in \mathcal{M}$ there is a neighbourhood $U$ of $p$ in $\mathcal{M}$ such that $\left.X\right|_{U}$ is a smooth embedding (see [Lee13, Theorem 4.25, p. 87]). As an open subset of $\mathcal{M}$ it is easy to show that $U$ is a smooth $m$-dimensional manifold (by restricting the coordinate charts to $U$, see [Lee13, Example 1.26, p. 19]). Then the map $\left.X\right|_{U}: U \rightarrow \mathcal{N}$ is a smooth embedding of a smooth $m$-dimensional manifold into a smooth $n$-dimensional manifold. Another result ([Lee13, Proposition 5.2, p. 99]) shows that the
image of a smooth embedding is again a smooth manifold of the same dimension as the manifold in the domain of the embedding. Thus, the image $S:=\left.X\right|_{U}(U) \subset \mathcal{N}$ is a smooth $m$-manifold in the ambient smooth $n$-manifold $\mathcal{N}$ wherefore it makes sense to consider the tangent space $T_{X(p)} S$ to $S$ at a point $X(p) \in \mathcal{N}$.

## Lemma 1.3.1

Let $X: \mathcal{M}^{m} \rightarrow \mathcal{N}^{n}$ be a smooth immersion, $p$ a point in $\mathcal{M}$ and $U \subset \mathcal{M}$ an open neighbourhood of $p$. Then the tangent space to $S:=X(U)$ at $X(p)$ is given by

$$
T_{X(p)} S=d\left(\left.X\right|_{U}\right)_{p}\left(T_{p} U\right)
$$

Proof: By assumption $d\left(\left.X\right|_{U}\right)_{p}: T_{p} U \rightarrow T_{X(p)} \mathcal{N}$ is an injective linear map and therefore bijective onto its image $d\left(\left.X\right|_{U}\right)_{p}\left(T_{p} U\right)$. Hence, both vector spaces $T_{X(p)} S$ and $d\left(\left.X\right|_{U}\right)_{p}\left(T_{p} U\right)$ are $m$-dimensional. For dimensional reasons we need only show that $d\left(\left.X\right|_{U}\right)_{p}\left(T_{p} U\right) \subset T_{X(p)} S$ to finish the proof. Thus, suppose $w \in d\left(\left.X\right|_{U}\right)_{p}\left(T_{p} U\right)$. Then there is $v \in T_{p} U$ such that $w=d\left(\left.X\right|_{U}\right)_{p}(v)$. As a composition of linear maps, $w$ is a linear map from $C^{\infty}(S)$ to $\mathbb{R}$. We need to show the product rule (1.3.22). Let $f, g \in C^{\infty}(S)$. Then $\left.f \circ X\right|_{U},\left.g \circ X\right|_{U} \in C^{\infty}(U)$. Unraveling the definitions and using (1.3.22) for $v \in T_{p} U$ we get

$$
\begin{aligned}
w(f g) & =d\left(\left.X\right|_{U}\right)_{p}(v)(f g) \\
& =v\left(\left.(f g) \circ X\right|_{U}\right) \\
& =v\left(\left(\left.f \circ X\right|_{U}\right)\left(\left.g \circ X\right|_{U}\right)\right) \\
& =\left(\left.f \circ X\right|_{U}\right)(p) v\left(\left.g \circ X\right|_{U}\right)+\left(\left.g \circ X\right|_{U}\right)(p) v\left(\left.f \circ X\right|_{U}\right) \\
& =f(X(p)) d\left(\left.X\right|_{U}\right)_{p}(v)(g)+g(X(p)) d\left(\left.X\right|_{U}\right)_{p}(v)(f) \\
& =f(X(p)) w(g)+g(X(p)) w(f)
\end{aligned}
$$

Therefore, $w$ is a derivation at $X(p)$ and so $w \in T_{X(p)} S$.

Furthermore, the injectivity of the differential of a smooth immersion immediately yields the next result.

## Corollary 1.3.2

Let $X: \mathcal{M}^{m} \rightarrow \mathcal{N}^{n}$ be a smooth immersion and $\left(U,\left(u^{\alpha}\right)\right)$ a smooth coordinate chart of $\mathcal{M}$ at some
point $p$. Then

$$
\left(d X_{p}\left(\left.\frac{\partial}{\partial u^{1}}\right|_{p}\right), d X_{p}\left(\left.\frac{\partial}{\partial u^{2}}\right|_{p}\right), \ldots, d X_{p}\left(\left.\frac{\partial}{\partial u^{m}}\right|_{p}\right)\right)
$$

is an ordered basis for the m-dimensional subspace $T_{X(p)} X(U)$ of the $n$-dimensional tangent space $T_{X(p)} \mathcal{N}$ to the ambient manifold $\mathcal{N}$ at $X(p)$.

### 1.3.2 Finsler manifolds

Let $\mathcal{N}$ be a smooth $n$-dimensional manifold and $T \mathcal{N}=\bigsqcup_{p \in \mathcal{N}} T_{p} \mathcal{N}$ be its tangent bundle. The subset $o:=\{(p, 0) \in T \mathcal{N}\}$ of $T \mathcal{N}$ is called the zero section of $T \mathcal{N}$.

A non-negative function $F: T \mathcal{N} \rightarrow[0, \infty)$ is called a Finsler metric if the following three conditions are satisfied:

- Regularity : $F \in C^{\infty}(T \mathcal{N} \backslash o) \cap C^{0}(T \mathcal{N})$.
- Positive homogeneity: $F(p, t v)=t F(p, v)$ for all $t>0$ and $(p, v) \in T \mathcal{N}$.
- Ellipticity: The coefficients $g_{i j}^{F}(p, v):=\left(\frac{1}{2} F^{2}\right)_{y^{i} y^{j}}(p, v)$ form a positive definite matrix for any $(p, v) \in T \mathcal{N} \backslash o$.

The third condition is to be understood in the following way. Recall that a smooth coordinate chart $\left(U,\left(x^{1}, x^{2}, \ldots, x^{n}\right)\right)$ containing $p \in \mathcal{N}$ induces natural coordinates $\widetilde{\varphi}=\left(x^{i}, y^{i}\right)$ on the tangent bundle by $v=\left.y^{i} \frac{\partial}{\partial x^{i}}\right|_{p} \in T_{p} \mathcal{N}$ (see (1.3.23)). Then in local coordinate representation $F$ is given by

$$
F(p, v)=\widehat{F}\left(x^{1}, x^{2}, \ldots, x^{n}, y^{1}, y^{2}, \ldots, y^{n}\right)
$$

It is with respect to the local expression $\widehat{F}=F \circ \widetilde{\varphi}^{-1}$ that we take the derivatives with respect to $y^{i}$ and $y^{j}$. The collection of coefficients $\left(g_{i j}^{F}(p, v)\right)_{i j}$ are called the fundamental tensor and locally define the fundamental form by

$$
\left.g^{F}\right|_{(p, v)}(u, w):=\left.\frac{1}{2} \frac{\partial^{2}}{\partial s \partial t}\right|_{s=t=0} F^{2}(p, v+s u+t w)=g_{i j}^{F}(p, v) u^{i} w^{j}
$$

where $u=\left.u^{i} \frac{\partial}{\partial x^{i}}\right|_{p}, w=\left.w^{i} \frac{\partial}{\partial x^{i}}\right|_{p} \in T_{p} \mathcal{N}$. A Finsler metric is called reversible if $F(p, v)=F(p,-v)$ for all $(p, v) \in T \mathcal{N}$.

A smooth $n$-manifold $\mathcal{N}$ together with a Finsler metric $F$ is called a Finsler manifold and will be denoted by $(\mathcal{N}, F)$ or $\left(\mathcal{N}^{n}, F\right)$.

Let $p \in \mathcal{N}$ be fixed. One can show that the three conditions for a Finsler metric imply that $F(p, \cdot)$ is an asymmetric (only positive homogeneous) norm on the tangent space $T_{p} \mathcal{N}$, that is, the following additional two conditions hold (see [BCS00, Theorem 1.2.2, pp. 7-9]).

- Positivity: $F(p, v)>0$ whenever $v \neq 0 \in T_{p} \mathcal{N}$.
- Triangle inequality: $F\left(p, v_{1}+v_{2}\right) \leq F\left(p, v_{1}\right)+F\left(p, v_{2}\right)$ for all $v_{1}, v_{2} \in T_{p} \mathcal{N}$ where equality holds if and only if $v_{2}=\alpha v_{1}$ for some $\alpha \geq 0$.

If, in addition, the Finsler metric is reversible then one has absolute homogeneity and the tangent space $T_{p} \mathcal{N}$ becomes a normed space with norm $F(p, \cdot)$ in the usual sense of functional analysis. The positive homogeneity and the triangle inequality imply that $F(\cdot, \cdot)$ is a convex function in its second component.

A Finsler metric $F$ is called a local Minkowski metric if for each $p \in \mathcal{N}$ there are smooth coordinate charts $\left(U_{p},\left(x^{1}, x^{2}, \ldots, x^{n}\right)\right)$ such that the coordinate representation of $F$ at $p$ is independent of $x^{i}$, that is,

$$
F(p, v)=\widehat{F}\left(x^{1}, x^{2}, \ldots, x^{n}, y^{1}, y^{2}, \ldots, y^{n}\right)=\widehat{F}\left(y^{1}, y^{2}, \ldots, y^{n}\right)
$$

where $v=\left.y^{i} \frac{\partial}{\partial x^{i}}\right|_{p} \in T_{p} \mathcal{N}$. A Finsler manifold $(\mathcal{N}, F)$ where $F$ is a local Minkowski metric is called locally Minkowskian.

## Chapter 2

## The Busemann-Hausdorff definition of area in Finsler geometry

In the first section of this chapter, we formalise the notion of Finsler volume introduced by Busemann in [Bus47]. Subsequently, the Finsler area of an immersion arises. Similarly to Overath [Ove13] we introduce the $m$-dimensional Busemann-Hausdorff area integrand which in turn leads to the $m$-dimensional Busemann-Hausdorff area density.

The second section mainly contains an analysis of Burago and Ivanov's paper 'Minimality of planes in normed spaces' [BI12]. In their work they proved the convexity of the two-dimensional Busemann-Hausdorff area density. Using multilinear algebra we introduce the concept of a calibrator for an area density. We show that the convexity of a density is equivalent to the existence of such calibrators. Subsequently, the defining inequality for a calibrator will be reformulated into the language of convex geometry and an explicit construction for calibrators will be given. This finally results in an inequality for polytopes on the two-dimensional plane.

Note that in contrast to Euclidean space or more generally Riemannian manifolds, there is no unique notion of volume when dealing merely with a normed space or more generally a Finsler manifold. The Busemann-Hausdorff definition is only one of many possible choices of volume in the setting of a Finsler manifold. An alternative is the Holmes-Thompson volume ([Tho96, Chapter 6]) which we do not address here.

### 2.1 The m-dimensional Busemann-Hausdorff area density

Suppose $\left(\mathcal{M}^{m}, F\right)$ is a Finsler manifold of dimension $m$. Let $p \in \mathcal{M}$ be an arbitrary but fixed point and $\left(U_{p}, \varphi\right)$ be a smooth coordinate chart containing $p$ where we denote the coordinate functions of $\varphi$
as $\left(u^{1}, u^{2}, \ldots, u^{m}\right)$. Then we define the (closed) Finslerian unit ball at $\boldsymbol{p}$ as

$$
B_{p}^{F}:=\left\{v \in T_{p} \mathcal{M} \mid F(p, v) \leq 1\right\} \subset T_{p} \mathcal{M}
$$

The Busemann-Hausdorff volume form is given by the volume ratio of the Euclidean and the Finslerian unit ball. To make sense of the latter notion we need to use the local coordinate chart $\left(U_{p}, \varphi\right)$ and the Hausdorff measure in $\mathbb{R}^{m}$. Let $\pi: T \mathcal{M} \rightarrow \mathcal{M},(q, v) \rightarrow q$ be the natural projection map of the tangent bundle and recall that $\varphi$ induces natural coordinates $\left(\pi^{-1}\left(U_{p}\right), \widetilde{\varphi}\right)$ on the tangent bundle $T \mathcal{M}$ by (1.3.23), that is,

$$
\widetilde{\varphi}: \pi^{-1}\left(U_{p}\right) \rightarrow \varphi\left(U_{p}\right) \times \mathbb{R}^{m},\left(q,\left.v^{\alpha} \frac{\partial}{\partial u^{\alpha}}\right|_{q}\right) \mapsto\left(u^{1}(q), u^{2}(q), \ldots, u^{m}(q), v^{1}, v^{2}, \ldots, v^{m}\right) .
$$

The Finslerian unit ball at $p$ in the tangent bundle - the set $\left\{(p, v) \mid v \in B_{p}^{F}\right\} \subset T \mathcal{M}$ - is contained in $\pi^{-1}\left(U_{p}\right)$. Let $\pi_{2}: \mathbb{R}^{m} \times \mathbb{R}^{m} \rightarrow \mathbb{R}^{m},(x, y) \mapsto y$ be the projection map onto the second component. Then the Finslerian unit ball at $p$ corresponds to the set

$$
\widetilde{B}_{p}^{F}:=\pi_{2} \circ \widetilde{\varphi}\left(\left\{(p, v) \mid v \in B_{p}^{F}\right\}\right) \subset \mathbb{R}^{m}
$$

Note that by using the natural coordinates one can rewrite this set as

$$
\widetilde{B}_{p}^{F}=\left\{v=\left(v^{1}, v^{2}, \ldots, v^{m}\right) \in \mathbb{R}^{m} \left\lvert\, F\left(p,\left.v^{\alpha} \frac{\partial}{\partial u^{\alpha}}\right|_{p}\right) \leq 1\right.\right\}
$$

The $\boldsymbol{m}$-dimensional Busemann-Hausdorff volume form is given by

$$
\begin{equation*}
d V_{F}(p):=\sigma_{F}(p) d u^{1} \wedge d u^{2} \wedge \cdots \wedge d u^{m} \tag{2.1.1}
\end{equation*}
$$

where

$$
\begin{equation*}
\sigma_{F}(p):=\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\widetilde{B}_{p}^{F}\right)} \tag{2.1.2}
\end{equation*}
$$

Here $\varepsilon_{m}:=\mathcal{H}^{m}\left(B_{1}^{m}(0)\right)$ denotes the volume of the $m$-dimensional Euclidean unit ball in $\mathbb{R}^{m}$ and $\mathcal{H}^{m}$ is the $m$-dimensional Hausdorff measure on $\mathbb{R}^{m}$.

The Busemann-Hausdorff volume of a subset $\Omega \subset \mathcal{M}$ is defined as

$$
\begin{equation*}
\operatorname{vol}_{F}(\Omega):=\int_{p \in \Omega} d V_{F}(p) \tag{2.1.3}
\end{equation*}
$$

If $F$ is a local Minkowski metric, that is, independent of the point $p$ and $\mathcal{M}=\mathbb{R}^{m}$ then the Busemann-Hausdorff volume of the Finslerian unit ball $B^{F}=\left\{v \in \mathbb{R}^{m} \mid F(v) \leq 1\right\}$ equals the volume of the Euclidean unit ball

$$
\begin{aligned}
\operatorname{vol}_{F}\left(B^{F}\right) & =\int_{p \in \Omega} d V_{F}(p) \\
& =\int_{p \in B^{F}} \frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\widetilde{B}_{p}^{F}\right)} d u^{1} \wedge d u^{2} \wedge \cdots \wedge d u^{m} \\
& =\int_{p \in B^{F}} \frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\widetilde{B}^{F}\right)} d u^{1} \wedge d u^{2} \wedge \cdots \wedge d u^{m} \\
& =\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(B^{F}\right)} \int_{p \in B^{F}} d u^{1} d u^{2} \cdots d u^{m} \\
& =\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(B^{F}\right)} \mathcal{H}^{m}\left(B^{F}\right) \\
& =\varepsilon_{m} .
\end{aligned}
$$

It is possible to show that for absolutely homogeneous Finsler metrics $F$, the Finsler manifold $(\mathcal{M}, F)$ inherits a metric space structure ([BCS00, $\S 6.2$, p. 145]). Busemann showed in [Bus47] that the Finsler volume $\operatorname{vol}_{F}$ of a set $\Omega \subset \mathcal{M}$ coincides with the Hausdorff measure induced by this metric space structure. Using the Busemann-Hausdorff definition of volume we can introduce a way of measuring surface area of a submanifold immersed into an ambient manifold.

Let $m \leq n$ and $X: \mathcal{M}^{m} \rightarrow \mathcal{N}^{n}$ be a smooth immersion of a smooth $m$-manifold $\mathcal{M}$ into a Finsler $n$-manifold $(\mathcal{N}, F)$. Then $\mathcal{M}$ inherits a pull-back Finsler metric through

$$
\begin{equation*}
X^{*} F(p, v):=F\left(X(p), d X_{p}(v)\right) \tag{2.1.4}
\end{equation*}
$$

for $(p, v) \in T_{p} \mathcal{M}$. The map $X^{*} F: T \mathcal{M} \rightarrow[0, \infty)$ is indeed a Finsler metric, because $X$ is smooth and its differential at each $p \in \mathcal{M}$ injective. The Busemann-Hausdorff area of the immersion $X$ for a measurable subset $\Omega \subset \mathcal{M}$ is given by

$$
\begin{equation*}
\operatorname{area}_{\Omega}^{F}(X):=\int_{p \in \Omega} d V_{X^{*} F}(p) \tag{2.1.5}
\end{equation*}
$$

Therein, $d V_{X{ }^{*} F}(p)$ is given by the formula (2.1.1).
Let $p \in \mathcal{M}$ be an arbitrary but fixed point and $\left(U_{p},\left(u^{1}, u^{2}, \ldots, u^{m}\right)\right)$ be a smooth coordinate chart containing $p$. The local Busemann-Hausdorff integrand $\sigma_{X^{*} F}(p)$ can be rewritten as

$$
\sigma_{X^{*} F}(p)=\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\left\{v=\left(v^{1}, v^{2}, \ldots, v^{m}\right) \in \mathbb{R}^{m} \left\lvert\, X^{*} F\left(p,\left.v^{\alpha} \frac{\partial}{\partial u^{\alpha}}\right|_{p}\right) \leq 1\right.\right\}\right)}
$$

$$
=\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\left\{v \in \mathbb{R}^{m} \left\lvert\, F\left(X(p), v^{\alpha} d X_{p}\left(\left.\frac{\partial}{\partial u^{\alpha}}\right|_{p}\right)\right) \leq 1\right.\right\}\right)}
$$

Therefore, we define the map

$$
\begin{align*}
a_{m}^{F}: & \bigsqcup_{q \in \mathcal{N}} G C_{m}\left(T_{q} \mathcal{N}\right) \rightarrow \mathbb{R}_{+},  \tag{2.1.6}\\
& \left(q, w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right) \mapsto \frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\left\{v=\left(v^{1}, v^{2}, \ldots, v^{m}\right) \in \mathbb{R}^{m} \mid F\left(q, v^{\alpha} w_{\alpha}\right) \leq 1\right\}\right)}
\end{align*}
$$

Recall that $G C_{m}\left(T_{q} \mathcal{N}\right) \subset \bigwedge^{m}\left(T_{q} \mathcal{N}\right)$ denotes the set of simple $m$-vectors on $T_{q} \mathcal{N}$. Note that this map is well-defined; for if $\widetilde{w}_{1} \wedge \widetilde{w}_{2} \wedge \cdots \wedge \widetilde{w}_{m}$ is another representation of $w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}$ then the injectivity of the Plücker embedding (Proposition 1.1.9) implies that the vectors $\left\{\widetilde{w}_{i}\right\}_{i=1}^{m}$ and $\left\{w_{i}\right\}_{i=1}^{m}$ span the same $m$-dimensional subspace. Therefore the denominator in (2.1.6) does not change value.

We call $a_{m}^{F}$ the $\boldsymbol{m}$-dimensional Busemann-Hausdorff area integrand or $\boldsymbol{m}$-dimensional Finsler area integrand. Notice the similarity to the Finsler area integrand defined by Overath in [Ove13, Equation (2.1.5), p. 68]. Therein, Overath uses the representation of the differential in local coordinates through its Jacobian matrix (see (1.3.24)) and defines the Finsler area integrand locally over $\mathcal{N} \times \mathbb{R}^{n \times m}$. In the present work we use the above definition to apply Burago and Ivanov's result. Note that our definition is not a local one - which contrasts [Ove13].

In local coordinates about $p$ in $\mathcal{M}$ we thereby get

$$
\sigma_{X^{*} F}(p)=a_{m}^{F}\left(X(p), d X_{p}\left(\left.\frac{\partial}{\partial u^{1}}\right|_{p}\right) \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{2}}\right|_{p}\right) \wedge \cdots \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{m}}\right|_{p}\right)\right)
$$

Recall from the previous chapter that locally the image of the immersion $X$ is a submanifold of dimension $m$ in the ambient $n$-manifold $\mathcal{N}$. By Corollary 1.3 .2 we know that $d X_{p}\left(\left.\frac{\partial}{\partial u^{1}}\right|_{p}\right) \wedge$ $d X_{p}\left(\left.\frac{\partial}{\partial u^{2}}\right|_{p}\right) \wedge \cdots \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{m}}\right|_{p}\right)$ corresponds to the tangent space $T_{X(p)} X\left(U_{p}\right)$ through the Plücker embedding (Proposition 1.1.9). Thus, we may rewrite the expression (2.1.5) for the BusemannHausdorff area of an immersion $X$ as

$$
\begin{equation*}
\operatorname{area}_{\Omega}^{F}(X)=\int_{p \in \Omega} a_{m}^{F}\left(X(p), d X_{p}\left(\left.\frac{\partial}{\partial u^{1}}\right|_{p}\right) \wedge \cdots \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{m}}\right|_{p}\right)\right) d u^{1} \wedge \cdots \wedge d u^{m} . \tag{2.1.7}
\end{equation*}
$$

Suppose now that $\mathcal{N}=\mathbb{R}^{n}$ and $F$ is a reversible Finsler metric. Fix $q \in \mathbb{R}^{n}$. Then the tangent space $T_{q} \mathcal{N}=T_{q} \mathbb{R}^{n}$ is canonically isomorphic to $\mathbb{R}^{n}$ itself ([Lee13, Proposition 3.13, p. 59]). Since $F$ is a reversible Finsler metric, the function $\|\cdot\|_{V}:=F(q, \cdot)$ defines a norm on the vector space $V:=\mathbb{R}^{n} \cong T_{q} \mathbb{R}^{n}$. Denote the closed unit ball by $B:=B_{q}^{F}=\{v \in V \mid F(q, v) \leq 1\}$. Then $\left(V,\|\cdot\|_{V}\right)$ is
an $n$-dimensional normed space. More generally, we can make the following definition.

## Definition 2.1.1

An $\boldsymbol{m}$-dimensional density on an $n$-dimensional normed space $\left(V,\|\cdot\|_{V}\right)$ is a continuous function $A: G C_{m}(V) \rightarrow \mathbb{R}_{+}$which is absolutely homogeneous, that is $A(\lambda \sigma)=|\lambda| A(\sigma)$ for all $\lambda \in \mathbb{R}$ and $\sigma \in G C_{m}(V)$.

## Definition 2.1.2

The $\boldsymbol{m}$-dimensional Busemann-Hausdorff area density $A^{b h}=A_{V, m}^{b h}: G C_{m}(V) \rightarrow \mathbb{R}_{+}$in a normed space $(V,\|\cdot\|)$ is defined by

$$
\begin{equation*}
A^{b h}(\sigma):=\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(L_{\sigma}^{-1}(B)\right)} \tag{2.1.8}
\end{equation*}
$$

wherein $\sigma=v_{1} \wedge v_{2} \wedge \ldots \wedge v_{m}$ and $L_{\sigma}: \mathbb{R}^{m} \rightarrow V$ is the linear map that takes the standard basis $\left(e_{1}, e_{2}, \ldots, e_{m}\right)$ of $\mathbb{R}^{m}$ to $\left(v_{1}, v_{2}, \ldots, v_{m}\right)$, that is, $L_{\sigma} e_{\alpha}=v_{\alpha}$ for $\alpha=1,2, \ldots, m$.

Notice that $v_{1}, v_{2}, \ldots, v_{m} \in V$ are linearly independent (because $\sigma \in G C_{m}(V)$ ) and the map $L_{\sigma}$ is always injective (and thus, bijective onto its image). Notice further that the map $L_{\sigma}$ is not well-defined as a function of $\sigma$. However, the map $A^{b h}$ is well-defined as we will see in Theorem 2.1.3.

The reason for the introduction of the function $A^{b h}$ becomes clear if we restrict our point of view again to the specific normed space $\left(\mathbb{R}^{n}, F(q, \cdot)\right)$ with unit ball $B=B_{q}^{F}$. Then for fixed $q \in \mathbb{R}^{n}=\mathcal{N}$ we calculate the value of $A_{q, m}^{b h}=A^{b h}$ to

$$
\begin{aligned}
A_{q, m}^{b h}\left(v_{1} \wedge v_{2} \wedge \ldots \wedge v_{m}\right) & =\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\left\{w \in \mathbb{R}^{m} \mid L w \in B\right\}\right)} \\
& =\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\left\{w=w^{\alpha} e_{\alpha} \in \mathbb{R}^{m} \mid w^{\alpha} v_{\alpha} \in B_{q}^{F}\right\}\right)} \\
& =\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\left\{w=w^{\alpha} e_{\alpha} \in \mathbb{R}^{m} \mid F\left(q, w^{\alpha} v_{\alpha}\right) \leq 1\right\}\right)}
\end{aligned}
$$

So, the $m$-dimensional Finsler area integrand given by (2.1.6) corresponds to the BusemannHausdorff area density through

$$
\begin{equation*}
a_{m}^{F}\left(q, v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right)=A_{q, m}^{b h}\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right) \tag{2.1.9}
\end{equation*}
$$

## Theorem 2.1.3

The Busemann-Hausdorff area density $A^{b h}$ is an m-dimensional density as defined in Definition 2.1.1. Proof: First, we show that the map $A^{b h}$ is both well-defined and absolutely homogeneous. Let
$\sigma=v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}$ and $\tau=w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}$ be two simple $m$-vectors and let there be a non-zero scalar $\lambda$ such that $\sigma=\lambda \tau$. This means $\sigma$ and $\tau$ span the same one-dimensional subspace of the $m$ th exterior power. We recall that the Plücker embedding maps an $m$-dimensional subspace of $V$ to the one-dimensional line spanned by the $m$-wedge of a basis of this subspace. In fact, because the Plücker embedding is injective (Proposition 1.1.9) we know that

$$
\operatorname{span}\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}=P=\operatorname{span}\left\{w_{1}, w_{2}, \ldots, w_{m}\right\}
$$

Then there is a bijective, linear map $T: P \rightarrow P$ which sends $w_{\alpha}$ to $v_{\alpha}$. By virtue of Lemma 1.1.7 (v) we know, that $\sigma=\operatorname{det} T \tau$ and thus $\operatorname{det} T=\lambda$. Consider the two linear maps from the definition of the Busemann-Hausdorff density

$$
\begin{aligned}
& L_{\sigma}: \mathbb{R}^{m} \rightarrow P, e_{\alpha} \mapsto v_{\alpha} \\
& L_{\tau}: \mathbb{R}^{m} \rightarrow P, e_{\alpha} \mapsto w_{\alpha}
\end{aligned}
$$

Each map is bijective onto its image, so we can define $\widetilde{v}_{\alpha}:=L_{\tau}^{-1}\left(v_{\alpha}\right) \in \mathbb{R}^{m}$. Then $\left\{\widetilde{v}_{1}, \widetilde{v}_{2}, \ldots, \widetilde{v}_{m}\right\}$ is a basis of $\mathbb{R}^{m}$. Now consider the linear map $\widetilde{T}: \mathbb{R}^{m} \rightarrow \mathbb{R}^{m}$ sending the standard basis vector $e_{\alpha}$ to $\widetilde{v}_{\alpha}$. We can calculate

$$
\begin{aligned}
& \widetilde{T} e_{\alpha}=\widetilde{v}_{\alpha}=L_{\tau}^{-1}\left(v_{\alpha}\right)=L_{\tau}^{-1} \circ L_{\sigma} e_{\alpha} \\
& \text { and } \quad L_{\sigma} e_{\alpha}=v_{\alpha}=T w_{\alpha}=T \circ L_{\tau} e_{\alpha} .
\end{aligned}
$$

Since this holds for $\alpha=1,2, \ldots, m$, we get

$$
\begin{equation*}
L_{\sigma}=L_{\tau} \circ \widetilde{T}=T \circ L_{\tau} \tag{2.1.10}
\end{equation*}
$$

From (2.1.10) we conclude that $\operatorname{det} L_{\tau} \operatorname{det} \widetilde{T}=\operatorname{det}\left(L_{\tau} \circ \widetilde{T}\right)=\operatorname{det}\left(T \circ L_{\tau}\right)=\operatorname{det} T \operatorname{det} L_{\tau}$ and thus because $L_{\tau}$ is bijective we know

$$
\operatorname{det} \widetilde{T}=\operatorname{det} T=\lambda
$$

Therefore, it holds that

$$
\begin{align*}
\mathcal{H}^{m}\left(L_{\sigma}^{-1}(B)\right) & =\mathcal{H}^{m}\left(\left(L_{\tau} \circ \widetilde{T}\right)^{-1}(B)\right)=\mathcal{H}^{m}\left(\widetilde{T}^{-1} \circ L_{\tau}^{-1}(B)\right)  \tag{2.1.11}\\
& =\left|\operatorname{det} \widetilde{T}^{-1}\right| \mathcal{H}^{m}\left(L_{\tau}^{-1}(B)\right)=\left|\lambda^{-1}\right| \mathcal{H}^{m}\left(L_{\tau}^{-1}(B)\right)
\end{align*}
$$

where we used the fact that linear transformations change the $m$-dimensional Lebesgue measure of $m$-dimensional sets by a constant factor - the absolute value of the determinant of the transformation (see e.g. [For09, $\S 5$, Satz 2, p. 48]). By definition of the Busemann-Hausdorff area density we finally conclude with (2.1.11) that

$$
\begin{equation*}
A^{b h}(\lambda \tau)=A^{b h}(\sigma)=\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(L_{\sigma}^{-1}(B)\right)}=|\lambda| \frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(L_{\tau}^{-1}(B)\right)}=|\lambda| A^{b h}(\tau) \tag{2.1.12}
\end{equation*}
$$

Thus, we have both shown absolute homogeneity and well-definedness. For if $\lambda=1$, that is, if $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}=w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}$ are two representations of the same simple $m$-vector $\sigma$, then $A^{b h}\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right)=A^{b h}\left(w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right)$ by (2.1.12).

For the continuity of the Busemann-Hausdorff area density we sketch a proof here and refer to [BI12] or [ÁPT04, Exercise 3.5, p. 11 and pp. 18-19]. We recall from Lemma 1.1.10 that the line spanned by a simple $m$-vector $\sigma \in G C_{m}(V)$ is in the image of the Plücker embedding. Then the injectivity of the latter mapping (Proposition 1.1.9) shows the existence of a unique $m$-dimensional plane $P_{\sigma} \in G_{m}(V)$ such that $\rho\left(P_{\sigma}\right)=[\sigma]_{\sim}$. Then $\operatorname{im}\left(L_{\sigma}\right)=P_{\sigma}$, so that

$$
\mathcal{H}^{m}\left(L_{\sigma}^{-1}(B)\right)=\mathcal{H}^{m}\left(L_{\sigma}^{-1}\left(B \cap P_{\sigma}\right)\right) .
$$

The isomorphism $L_{\sigma}: \mathbb{R}^{m} \rightarrow P_{\sigma}$ induces an inner product on $P_{\sigma}$ (see the introduction to Section 1.1.4 for details why this defines an inner product). Thus, we can choose an orthogonal basis of $P_{\sigma}$ and write $\sigma=v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}$ where

$$
\left\langle v_{i}, v_{j}\right\rangle_{P_{\sigma}}=0 \quad \text { for } i \neq j
$$

The collection of convex bodies in $\mathbb{R}^{m}$, denoted by $\mathcal{C}_{b}\left(\mathbb{R}^{m}\right)$, turns into a metric space by means of the Hausdorff metric $d_{\mathcal{H}}=d_{\mathcal{H}, \mathbb{R}^{m}}$ given by

$$
d_{\mathcal{H}}(A, B)=\max \left(\sup _{a \in A} \inf _{b \in B}\|a-b\|_{R^{m}}, \sup _{b \in B} \inf _{a \in A}\|a-b\|_{R^{m}}\right)
$$

where $\|\cdot\|_{\mathbb{R}^{m}}$ is the standard Euclidean norm on $\mathbb{R}^{m}$ (see [Tho96, Proposition 2.4.2, p. 61]). A result from convex geometry states that the $m$-dimensional Lebesgue measure $\mathcal{L}^{m}$ on $\mathbb{R}^{m}$ is continuous on the collection of compact convex sets in $\mathbb{R}^{m}$ with respect to the Hausdorff metric (see [Val68, Satz 12.7, p. 153] or [Bee74, p. 64]). Note that the $m$-dimensional Hausdorff measure equals the $m$-dimensional Lebesgue measure on $\mathbb{R}^{m}$ (see [EG92, Theorem 2, p. 70]).

Due to Lemma 1.2.16 (ii) and (iv), the set $L_{\sigma}^{-1}\left(B \cap P_{\sigma}\right)$ is a symmetric convex body in $\mathbb{R}^{m}$. Thus,


Figure 2.1: Planes corresponding to simple $m$-vectors where $B$ is an ellipsoid
it suffices to prove the continuity of the map

$$
\varphi: G C_{m}(V) \rightarrow \mathcal{C}_{b}\left(\mathbb{R}^{m}\right), \sigma=v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m} \mapsto L_{\sigma}^{-1}\left(B \cap P_{\sigma}\right)
$$

where we chose an orthogonal basis $\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ of $P_{\sigma}$ as stated above. Since $P_{\sigma}$ and so $\langle\cdot, \cdot\rangle_{P_{\sigma}}$ is unique, this map is well-defined.

To show the continuity of $\varphi$, let $\left(\sigma_{k}\right)_{k \in \mathbb{N}} \subset G C_{m}(V)$ be a sequence converging to $\sigma \in G C_{m}(V)$. Denote the corresponding $m$-dimensional planes by $P_{k}$ and $P$ and write $\sigma_{k}=v_{1}^{(k)} \wedge v_{2}^{(k)} \wedge \cdots \wedge v_{m}^{(k)}$ and $\sigma=v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}$. Remember that we chose the bases of $P_{k}$ and $P$ to be orthogonal (in $P_{k}$ and $P$ respectively). Further, we can rotate them within the planes $P_{k}$ and $P$, so that

$$
\begin{aligned}
& \quad\left\langle v_{i}, v_{j}\right\rangle_{P}=\left\langle v_{i}^{(k)}, v_{j}^{(k)}\right\rangle_{P_{k}}=0 \text { for } i \neq j \text { and all } k \in \mathbb{N} \\
& \text { and } \quad v_{i}^{(k)} \xrightarrow[k \rightarrow \infty]{\longrightarrow} v_{i} \text { for } i=1,2, \ldots, m
\end{aligned}
$$

Here the inner products $\langle\cdot, \cdot\rangle_{P}$ and $\langle\cdot, \cdot\rangle_{P_{k}}$ are induced by the isomorphisms $L_{\sigma}: \mathbb{R}^{m} \rightarrow P$ and $L_{\sigma_{k}}: \mathbb{R}^{m} \rightarrow P_{k}$ respectively. We recall that the Plücker embedding is a topological embedding, that is, a homeomorphism onto its image (see remark after Proposition 1.1.9). Therefore, if $\sigma_{k} \xrightarrow[k \rightarrow \infty]{ } \sigma$ then $P_{k} \xrightarrow[k \rightarrow \infty]{\longrightarrow} P$ in $G_{m}(V)$.

We need to estimate

$$
\begin{aligned}
d_{\mathcal{H}}\left(L_{\sigma_{k}}^{-1}(B),\right. & \left.L_{\sigma}^{-1}(B)\right)=\max \left(\sup _{a \in L_{\sigma_{k}}^{-1}(B)} \inf _{b \in L_{\sigma}^{-1}(B)}\|a-b\|_{R^{m}}, \sup _{b \in L_{\sigma}^{-1}(B)} \inf _{a \in L_{\sigma_{k}}^{-1}(B)}\|a-b\|_{R^{m}}\right) \\
& =\max \left(\sup _{x \in B \cap P_{k}} \inf _{y \in B \cap P}\left\|L_{\sigma_{k}}^{-1} x-L_{\sigma}^{-1} y\right\|_{R^{m}}, \sup _{y \in B \cap P} \inf _{x \in B \cap P_{k}}\left\|L_{\sigma_{k}}^{-1} x-L_{\sigma}^{-1} y\right\|_{R^{m}}\right) .
\end{aligned}
$$

Hence, let $x \in B \cap P_{k}$ und $y \in B \cap P$. Further, let $\Pi_{k}: V \rightarrow P_{k}$ be the projection onto $P_{k}$. Then

$$
\begin{aligned}
\left\|L_{\sigma_{k}}^{-1} x-L_{\sigma}^{-1} y\right\|_{\mathbb{R}^{m}} & \leq\left\|L_{\sigma_{k}}^{-1} x-L_{\sigma_{k}}^{-1} \Pi_{k} y\right\|_{\mathbb{R}^{m}}+\left\|L_{\sigma_{k}}^{-1} \Pi_{k} y-L_{\sigma}^{-1} y\right\|_{\mathbb{R}^{m}} \\
& =\left\|L_{\sigma_{k}}^{-1} \Pi_{k}(x-y)\right\|_{\mathbb{R}^{m}}+\left\|\left(L_{\sigma_{k}}^{-1} \Pi_{k}-L_{\sigma}^{-1}\right) y\right\|_{\mathbb{R}^{m}} \\
& \leq\left\|L_{\sigma_{k}}^{-1} \Pi_{k}\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)}\|x-y\|_{\mathbb{R}^{m}}+\left\|\left(L_{\sigma_{k}}^{-1} \Pi_{k}-L_{\sigma}^{-1}\right) y\right\|_{\mathbb{R}^{m}}
\end{aligned}
$$

and similarly,

$$
\left\|L_{\sigma_{k}}^{-1} x-L_{\sigma}^{-1} y\right\|_{\mathbb{R}^{m}} \leq\left\|L_{\sigma}^{-1} \Pi\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)}\|x-y\|_{\mathbb{R}^{m}}+\left\|\left(L_{\sigma_{k}}^{-1}-L_{\sigma}^{-1} \Pi\right) x\right\|_{\mathbb{R}^{m}}
$$

where $\Pi: V \rightarrow P$ is the projection onto $P$. Thus, we can estimate

$$
\begin{align*}
\inf _{y \in B \cap P}\left\|L_{\sigma_{k}}^{-1} x-L_{\sigma}^{-1} y\right\|_{\mathbb{R}^{m}} & \leq \inf _{y \in B \cap P}\left(\left\|L_{\sigma_{k}}^{-1} \Pi_{k}\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)}\|x-y\|_{\mathbb{R}^{m}}+\left\|\left(L_{\sigma_{k}}^{-1} \Pi_{k}-L_{\sigma}^{-1}\right) y\right\|_{\mathbb{R}^{m}}\right) \\
& \leq \inf _{y \in B \cap P}\left\|L_{\sigma_{k}}^{-1} \Pi_{k}\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)}\|x-y\|_{\mathbb{R}^{m}}+\sup _{y \in B \cap P}\left\|\left(L_{\sigma_{k}}^{-1} \Pi_{k}-L_{\sigma}^{-1} \Pi\right) y\right\|_{\mathbb{R}^{m}} \\
& \leq\left\|L_{\sigma_{k}}^{-1} \Pi_{k}\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)} \operatorname{dist}(x, B \cap P)+\left\|L_{\sigma_{k}}^{-1} \Pi_{k}-L_{\sigma}^{-1} \Pi\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)} \tag{2.1.13}
\end{align*}
$$

and similarly,

$$
\begin{equation*}
\inf _{x \in B \cap P_{k}}\left\|L_{\sigma_{k}}^{-1} x-L_{\sigma}^{-1} y\right\|_{\mathbb{R}^{m}} \leq\left\|L_{\sigma}^{-1} \Pi\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)} \operatorname{dist}\left(y, B \cap P_{k}\right)+\left\|L_{\sigma_{k}}^{-1} \Pi_{k}-L_{\sigma}^{-1} \Pi\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)} \tag{2.1.14}
\end{equation*}
$$

Therefore, (2.1.13) and (2.1.14) together yield

$$
\begin{align*}
& d_{\mathcal{H}, \mathbb{R}^{m}}\left(L_{\sigma_{k}}^{-1}(B), L_{\sigma}^{-1}(B)\right) \\
& \leq \max \left(\left\|L_{\sigma_{k}}^{-1} \Pi_{k}\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)},\left\|L_{\sigma}^{-1} \Pi\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)}\right) d_{\mathcal{H}, V}\left(B \cap P_{k}, B \cap P\right)  \tag{2.1.15}\\
& \quad+\left\|L_{\sigma_{k}}^{-1} \Pi_{k}-L_{\sigma}^{-1} \Pi\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)}
\end{align*}
$$

Since we chose the bases $\left\{v_{i}^{(k)}\right\}_{i=1}^{m} \subset P_{k}$ and $\left\{v_{i}\right\}_{i=1}^{m} \subset P$ so that $v_{i}^{(k)} \xrightarrow[k \rightarrow \infty]{ } v_{i}$, we know that

$$
\begin{aligned}
\left\|L_{\sigma_{k}}-L_{\sigma}\right\|_{\mathcal{L}\left(R^{m}, V\right)} & =\sup _{\sum_{i=1}^{m}\left(a^{i}\right)^{2}=1}\left\|\sum_{i=1}^{m} a^{i}\left(L_{\sigma_{k}}-L_{\sigma}\right) e_{i}\right\|_{V} \\
& \leq \sup _{\sum_{i=1}^{m}\left(a^{i}\right)^{2}=1} \sum_{i=1}^{m}\left|a^{i}\right|\left\|v_{i}^{(k)}-v_{i}\right\|_{V} \xrightarrow[k \rightarrow \infty]{ } 0,
\end{aligned}
$$

that is, $L_{\sigma_{k}} \xrightarrow[k \rightarrow \infty]{ } L_{\sigma}$ in $\mathcal{L}\left(\mathbb{R}^{m}, V\right)$. Note that the linear map $L_{\sigma_{k}}^{-1} \Pi_{k}: V \rightarrow \mathbb{R}^{m}$ is the MoorePenrose pseudoinverse (or generalised inverse) of the linear map $L_{\sigma_{k}}: \mathbb{R}^{m} \rightarrow P_{k} \subset V$. All of the mappings $L_{\sigma_{k}}$ are of rank $m$. Under these conditions, the map sending a matrix to its Moore-Penrose inverse is continuous (see [Rak91, Theorem 4.2, p. 166] and [Rak97]). Therefore, we have that $L_{\sigma_{k}}^{-1} \Pi_{k} \xrightarrow[k \rightarrow \infty]{ } L_{\sigma}^{-1} \Pi$ in $\mathcal{L}\left(V, \mathbb{R}^{m}\right)$. Thus, the second summand in (2.1.15) vanishes for $k \rightarrow \infty$ and consequently, the factor $\max \left(\left\|L_{\sigma_{k}}^{-1} \Pi_{k}\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)},\left\|L_{\sigma}^{-1} \Pi\right\|_{\mathcal{L}\left(V, \mathbb{R}^{m}\right)}\right)$ is bounded above. It remains to show that

$$
d_{\mathcal{H}, V}\left(B \cap P_{k}, B \cap P\right) \underset{k \rightarrow \infty}{ } 0
$$

This is true because $P_{k} \xrightarrow[k \rightarrow \infty]{\longrightarrow} P$ in $G_{m}(V)$ as stated earlier (see Figure 2.1). So, $\varphi$ and therewith, $A^{b h}$ is continuous.

### 2.2 Convexity of the two-dimensional Busemann-Hausdorff area density

## Definition 2.2.1

A density $A: G C_{m}(V) \rightarrow \mathbb{R}_{+}$is called convex if it can be extended to an absolutely homogeneous, continuous and convex function on the entire vector space $\bigwedge^{m} V$.

Busemann proved in [Bus49] that the Busemann-Hausdorff area density is convex in codimension one, that is, $\operatorname{dim}(V)=m+1$, and the general case was left as a conjecture (see [Tho96, Problem 7.1.1, p. 310]). In their work [BI12] Burago and Ivanov proved the conjecture for the two-dimensional Busemann-Hausdorff area density. The next theorem is the main result of the present chapter. The rest of this section will be devoted to extensively discuss Burago and Ivanov's proof.

Theorem 2.2.2 (Convexity of the two-dimensional area density [BI12, Theorem 1, p. 630]) In every finite-dimensional normed space $V$, the two-dimensional Busemann-Hausdorff area density defined by (2.1.8) is convex on $\bigwedge^{2} V$.

In what follows, we will present the geometric argument Burago and Ivanov used to prove their statement. Using the language of convex geometry the statement will be reformulated into an inequality for polytopes on the two-dimensional plane. In order to do so, let us first introduce the concept of a calibrator for an $m$-dimensional density.

Definition 2.2.3 ([BI12, Definition 2.1, p. 631])
Let $V$ be a finite-dimensional vector space, $A: G C_{m}(V) \rightarrow \mathbb{R}_{+}$an $m$-dimensional density and $P \subset V$ an $m$-dimensional subspace. A calibrator or calibrating form for $P$ with respect to $A$ is an exterior $m$-form $\omega \in \bigwedge^{m}\left(V^{*}\right) \cong\left(\bigwedge^{m}(V)\right)^{*}$ such that for every simple $m$-vector $\sigma \in G C_{m}(V)$ one has $|\omega(\sigma)| \leq A(\sigma)$ and this inequality turns into equality if $\sigma \in G C_{m}(V) \cap \bigwedge^{m}(P)$.

Notice that $\bigwedge^{m}(P)$ is one-dimensional and therefore every $m$-vector in $\bigwedge^{m}(P)$ is simple. So that intersection $G C_{m}(V) \cap \bigwedge^{m}(P)$ really is just $\bigwedge^{m}(P)$ itself. The preceding notion of a calibrator is a powerful concept because the following central result holds.

## Lemma 2.2.4

An m-dimensional density $A: G C_{m}(V) \rightarrow \mathbb{R}_{+}$is convex on $\bigwedge^{m}(V)$ if and only if every m-dimensional plane $P \in G_{m}(V)$ admits a calibrator $\omega_{P} \in\left(\bigwedge^{m}(V)\right)^{*}$ with respect to $A$.

Proof: Suppose the $m$-dimensional density $A$ admits an absolutely homogeneous, continuous and convex extension $\mathcal{A}: \bigwedge^{m}(V) \rightarrow \mathbb{R}_{+}$. The following is due to [BRS12, p. 7]. The absolute homogeneity and convexity of the extension imply that $\mathcal{A}$ is a sublinear functional because for $\sigma, \tau \in \bigwedge^{m}(V)$ we
have

$$
\begin{aligned}
\mathcal{A}(\sigma+\tau) & =\mathcal{A}\left(\frac{1}{2}(2 \sigma)+\frac{1}{2}(2 \tau)\right) \leq \frac{1}{2} \mathcal{A}(2 \sigma)+\frac{1}{2} \mathcal{A}(2 \tau) \\
& =\mathcal{A}(\sigma)+\mathcal{A}(\tau)
\end{aligned}
$$

Let $P \in G_{m}(V)$ be arbitrary and spanned be a basis $\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$. Recall that the top exterior power $\Lambda^{m}(P)$ is $\binom{m}{m}=1$-dimensional. Note that any element of the top exterior power $\Lambda^{m}(P)$ is of the form $\lambda v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}$, that is, any element is a simple $m$-vector or $\bigwedge^{m}(P)=G C_{m}(V) \cap \bigwedge^{m}(P)$. Let us define $\omega_{P, 0}: \bigwedge^{m}(P) \rightarrow \mathbb{R}$ on this subspace by setting

$$
\omega_{P, 0}\left(\lambda v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right):=\lambda A\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right)
$$

This is a linear form on $\bigwedge^{m}(P)$ and for $\nu \in \bigwedge^{m}(P)$ we have

$$
\begin{aligned}
\left|\omega_{P, 0}(\nu)\right| & =\left|\omega_{P, 0}\left(\lambda v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right)\right| \\
& =\left|\lambda A\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right)\right| \\
& =|\lambda| A\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right) \\
& =A\left(\lambda v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right) \\
& =A(\nu)=\mathcal{A}(\nu)
\end{aligned}
$$

where we used the fact that the density $A$ is non-negative and absolutely homogeneous. Therefore, $\omega_{P, 0}$ is bounded above by $\mathcal{A}$ on the subspace $\bigwedge^{m}(P)$.

Then the Hahn-Banach theorem ([Tho96, Theorem 1.3.2, p. 33]) applied to $\omega_{P, 0}$ and $\mathcal{A}$, asserts the existence of a linear form $w_{P}: \bigwedge^{m}(V) \rightarrow \mathbb{R}$ such that

$$
\begin{aligned}
& \omega_{P}(\sigma)=\omega_{P, 0}(\sigma) \\
& \text { and all } \sigma \in \bigwedge^{m}(P) \\
& \omega_{P}(\sigma) \leq \mathcal{A}(\sigma) \text { for all } \sigma \in \bigwedge^{m}(V)
\end{aligned}
$$

In particular, this implies

$$
\begin{array}{ll} 
& \left|\omega_{P}(\sigma)\right|=A(\sigma) \quad \text { for all } \sigma \in \bigwedge^{m}(P)=G C_{m}(V) \cap \bigwedge^{m}(P) \\
\text { and } \quad\left|\omega_{P}(\sigma)\right| \leq A(\sigma) \quad \text { for all } \sigma \in G C_{m}(V)
\end{array}
$$

So this $m$-form is a calibrator for $P$ with respect to $A$.
Conversely, suppose every $m$-dimensional plane $P$ admits a calibrator $\omega_{P} \in\left(\bigwedge^{m}(V)\right)^{*}$ with respect
to $A$. Then the mapping $\left|\omega_{P}\right|$ is convex because $\omega_{P}$ is linear. In addition, for every $P \in G_{m}(V)$ the density $A$ is an upper bound to $\left|\omega_{P}\right|$ on the domain $G C_{m}(V)$ of $A$, that is, it holds that

$$
\begin{align*}
& \left|\omega_{P}(\sigma)\right| \leq A(\sigma) \quad \text { for all } \sigma \in G C_{m}(V) \\
\text { and } \quad\left|\omega_{P}(\sigma)\right|=A(\sigma) & \text { for all } \sigma \in G C_{m}(V) \cap \bigwedge^{m}(P) . \tag{2.2.16}
\end{align*}
$$

Construct the set of absolutely homogeneous, proper convex functions bounded above by $A$ on $G C_{m}(V)$,

$$
\mathcal{E}:=\left\{l:\left.\bigwedge^{m}(V) \rightarrow \mathbb{R}_{+}|l|\right|_{G C_{m}(V)} \leq A, l \text { is absolutely homogeneous and proper convex }\right\}
$$

Note that by the preceding argumentation $\left|\omega_{P}\right| \in \mathcal{E}$ for every $P \in G_{m}(V)$, thus $\mathcal{E}$ is non-empty. Then let us define the extension of $A$ as its convex envelope, that is, set

$$
\mathcal{A}(\sigma):=\operatorname{conv}(A)(\sigma):=\sup _{l \in \mathcal{E}} l(\sigma)
$$

for arbitrary $\sigma \in \bigwedge^{m}(V)$. The function $\mathcal{A}: \bigwedge^{m}(V) \rightarrow \mathbb{R}_{+}$is convex because each $l \in \mathcal{E}$ is and

$$
\begin{aligned}
\mathcal{A}(t \sigma+(1-t) \nu) & =\sup _{l \in \mathcal{E}} l(t \sigma+(1-t) \nu) \leq \sup _{l \in \mathcal{E}} t l(\sigma)+(1-t) l(\nu) \\
& \leq t \sup _{l \in \mathcal{E}} l(\sigma)+(1-t) \sup _{l \in \mathcal{E}} l(\nu) \\
& =t \mathcal{A}(\sigma)+(1-t) \mathcal{A}(\nu)
\end{aligned}
$$

Further, $\mathcal{A}$ is absolutely homogeneous since for $\sigma \in \bigwedge^{m}(V), \lambda \in \mathbb{R}$

$$
\mathcal{A}(\lambda \sigma)=\sup _{l \in \mathcal{E}} l(\lambda \sigma)=\sup _{l \in \mathcal{E}}|\lambda| l(\sigma)=|\lambda| \mathcal{A}(\sigma)
$$

For an arbitrary simple $m$-vector $\sigma \in G C_{m}(V)$ there is a corresponding $m$-dimensional plane $Q \in G_{m}(V)$ such that $\sigma \in \bigwedge^{m}(Q)$, namely, if $\sigma=v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}$ then $Q:=\operatorname{span}\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$. (Recall the Plücker embedding - in particular, Lemma 1.1.10.)

Thus, for the calibrator $\omega_{Q}$ of the plane $Q$ the inequality

$$
A(\sigma)=\left|\omega_{Q}(\sigma)\right| \leq \sup _{l \in \mathcal{E}} l(\sigma)=\mathcal{A}(\sigma)
$$

holds. On the other hand, because each $l \in \mathcal{E}$ is bounded above by $A$ on $G C_{m}(V)$, we have

$$
\mathcal{A}(\sigma)=\sup _{l \in \mathcal{E}} l(\sigma) \leq \sup _{l \in \mathcal{E}} A(\sigma)=A(\sigma)
$$

by definition of $\mathcal{E}$. So, $\left.\mathcal{A}\right|_{G C_{m}(V)} \equiv A$ and $\mathcal{A}$ is an absolutely homogeneous and convex extension of $A$. It remains to show the continuity of $\mathcal{A}$. Since each $l \in \mathcal{E}$ is proper the pointwise supremum over this set nowhere takes the value $-\infty$. On the set $G C_{m}(V)$ the function $\mathcal{A}$ is bounded above by $A$ which in turn only takes finite values. Now let $\left\{e_{1}, e_{2}, \ldots, e_{n}\right\}$ be a basis for $V$. Let $\sigma=\sum_{I \in \mathcal{I}} a_{\sigma}^{I} e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}$ be an arbitrary $m$-vector in $\bigwedge^{m}(V)$. As usual, $\mathcal{I}$ denotes the set of strictly increasing multi-indices of length $m$. Since $\mathcal{A}$ is convex and absolutely homogeneous, it is also sublinear by the same argument as in the first part of the proof. Thus, it follows that

$$
\begin{aligned}
\mathcal{A}(\sigma) & \leq \sum_{I \in \mathcal{I}}\left|a_{\sigma}^{I}\right| \mathcal{A}\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right) \\
& =\sum_{I \in \mathcal{I}}\left|a_{\sigma}^{I}\right| A\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right) \leq \widehat{C} \sum_{I \in \mathcal{I}}\left|a_{\sigma}^{I}\right|
\end{aligned}
$$

where $\widehat{C}:=\max _{I \in \mathcal{I}} A\left(e_{i_{1}} \wedge e_{i_{2}} \wedge \cdots \wedge e_{i_{m}}\right)<\infty$. Note that the Hölder-inequality for $p$-norms implies that

$$
\sum_{I \in \mathcal{I}}\left|a_{\sigma}^{I}\right|=\sum_{I \in \mathcal{I}}\left|a_{\sigma}^{I}\right| \cdot 1 \leq\left(\sum_{I \in \mathcal{I}}\left|a_{\sigma}^{I}\right|^{2}\right)^{\frac{1}{2}} \cdot\left(\sum_{I \in \mathcal{I}} 1\right)^{\frac{1}{2}} \leq D\binom{n}{m}^{\frac{1}{2}}\|\sigma\|_{\Lambda^{m}(V)}
$$

Therein, we used the equivalence of norms on the finite-dimensional vector space $\Lambda^{m}(V)$, since $\left(\sum_{I \in \mathcal{I}}\left|a_{\sigma}^{I}\right|^{2}\right)^{\frac{1}{2}}$ is the Euclidean norm (see Corollary 1.1.12). Hence, we get for $C:=\widehat{C} D$ that

$$
\begin{equation*}
\mathcal{A}(\sigma) \leq C\binom{n}{m}^{\frac{1}{2}}\|\sigma\|_{\Lambda^{m}(V)} \tag{2.2.17}
\end{equation*}
$$

Let $\varepsilon>0$ be arbitrary and consider an $m$-vector $\tau \in \bigwedge^{m}(V)$ with $\|\sigma-\tau\|_{\Lambda^{m}(V)}<\varepsilon$. Then using the sublinearity of $\mathcal{A}$ again yields

$$
\begin{aligned}
\mathcal{A}(\tau) & =\mathcal{A}((\sigma-\tau)+\sigma) \leq \mathcal{A}(\sigma-\tau)+\mathcal{A}(\sigma) \\
& \leq C\binom{n}{m}^{\frac{1}{2}}\left(\|\sigma-\tau\|_{\Lambda^{m}(V)}+\|\sigma\|_{\Lambda^{m}(V)}\right) \\
& <C\binom{n}{m}^{\frac{1}{2}}\left(\varepsilon+\|\sigma\|_{\Lambda^{m}(V)}\right)=: \widetilde{C}(n, m, \varepsilon ; \sigma)<\infty
\end{aligned}
$$

where we used (2.2.17) twice. Therefore, on a ball of radius $\varepsilon>0$ centered at $\sigma$ the convex function $\mathcal{A}$ is bounded above by a finite constant $\widetilde{C}(n, m, \varepsilon ; \sigma)$. Then a result from convex analysis ([Cla90, Proposition 2.2.6, p. 34]) states that $\mathcal{A}$ is Lipschitz-continuous on this neighbourhood of $\sigma$. Since $\sigma \in \bigwedge^{m}(V)$ was arbitrary we find that, in particular, $\mathcal{A}$ is continuous on $\bigwedge^{m}(V)$.

In what follows, we justify that we can make the subsequent simplifying assumption.

## Assumption 2.2.5

The unit ball $B$ of $\left(V,\|\cdot\|_{V}\right)$ is a symmetric polytope.

Denote the closed unit ball of the $n$-dimensional normed space $\left(V,\|\cdot\|_{V}\right)$ by $B$. Then the triangle inequality and absolute homogeneity of the norm imply that the unit ball $B$ is convex and symmetric. It has non-empty interior because $0 \in B$ and $\|0\|_{V}=0<1$. Since $V$ is finite-dimensional $B$ is a compact set (see [Tho96, Theorem 1.2.6, p. 30]). Therefore, $B$ is a symmetric, convex, compact set with non-empty interior or in the terminology of Section 1.2.4 the unit ball $B$ is a symmetric convex body.

As stated earlier, the collection of convex bodies in $V$, denoted by $C_{b}=C_{b}(V)$ and the Hausdorff metric $d_{\mathcal{H}}=d_{\mathcal{H}, V}$ given by

$$
d_{\mathcal{H}}(A, B)=\max \left(\sup _{a \in A} \inf _{b \in B}\|a-b\|_{V}, \sup _{b \in B} \inf _{a \in A}\|a-b\|_{V}\right)
$$

turns into a metric space. A result from convex geometry (see [Tho96, Theorem 2.5.1, p. 64]) states that the set of polytopes is dense in this metric space. In particular, this means that for a convex body $C \in C_{b}$ and any $\varepsilon>0$ there is a polytope $P$ contained in $C$ such that $d_{\mathcal{H}}(P, C)<\varepsilon$. An analysis of the proof of this density statement gives the following variation which shows that symmetric convex bodies can be approximated by symmetric polytopes.

## Lemma 2.2.6

If $C$ is a symmetric convex body then for any $\varepsilon>0$ there is a symmetric polytope $P$ such that $d_{\mathcal{H}}(P, C)<\varepsilon$.

Proof: The compactness of $C$ implies that for any $\varepsilon>0$ there is a finite set of points $F^{\prime}:=$ $\left\{x_{1}, x_{2}, \ldots, x_{N}\right\} \subset C$ such that $C \subset \bigcup_{i=1}^{N} B_{\varepsilon}\left(x_{i}\right)=F^{\prime}+\varepsilon B$. Define $F:=F^{\prime} \cup\left(-F^{\prime}\right)$. (In the original proof Thompson only chooses $F:=F^{\prime}$.) Since $C$ is symmetric, the extended set $F$ is contained in $C$ as well. Of course, $C \subset F+\varepsilon B$. Define the centrally symmetric polytope $P:=\operatorname{conv}(F)$. Since $C$ is
convex and taking the convex hull preserves inclusions, the sequence of set inclusions

$$
P \subset \operatorname{conv}(C)=C \subset F+\varepsilon B \subset P+\varepsilon B
$$

holds. Hence, $\sup _{p \in P} \inf _{c \in C}\|p-c\|=0$ by the first inclusion and $\sup _{c \in C} \inf _{p \in P}\|p-c\|<\varepsilon$ by the last. Therefore, $P$ is a centrally symmetric polytope such that $d_{\mathcal{H}}(P, C)<\varepsilon$.

If we apply the previous result to the convex body $B$ for $\varepsilon_{l}:=\frac{1}{l}$, we may choose a sequence of symmetric polytopes $\left(B_{l}\right)_{l \in \mathbb{N}} \subset C_{b}$ such that $B_{l} \xrightarrow[l \rightarrow \infty]{ } B$ in $\left(C_{b}, d_{\mathcal{H}}\right)$. These polytopes are themselves convex sets by definition of a polytope. Without loss of generality, we assume that each $B_{l}$ is a convex set of full dimension $n$ (as defined in Section 1.2.1). Then each line through 0 meets $B_{l}$ in a non-trivial (because 0 is an interior point of $B_{l}$ ), closed and bounded segment (since $B_{l}$ is closed and bounded being a polytope, see Lemma 1.2.3). Then the Minkowskifunctional of each polytope $B_{l}$ defined by

$$
\|x\|_{B_{l}}:=\inf \left\{t \in \mathbb{R}_{+} \mid x \in t B_{l}\right\}, x \in V
$$

is a norm on $V$ (see [Tho96, Theorem 1.1.8, p. 17]). In addition, $B_{l}$ is the closed $\|\cdot\|_{B_{l}}$-unit ball. To see the latter, let $x \in V$ and $\|x\|_{B_{l}} \leq 1$. Then by definition of the Minkowski functional $x \in t B_{l} \subset B_{l}$ where $t \leq 1$ is the infimum in the definition. Conversely, if $x \in B_{l}$ then $\|x\|_{B_{l}} \leq 1$.

Now we define the sequence of densities $\left(A_{l}^{b h}\right)_{l \in \mathbb{N}}$ where each $A_{l}^{b h}$ is the Busemann-Hausdorff area density on the normed space $\left(V,\|\cdot\|_{B_{l}}\right)$ as defined by (2.1.8). Recall that the linear map $L=L_{\sigma}: \mathbb{R}^{m} \rightarrow V, L e_{\alpha}=v_{\alpha}$ for $\alpha=1,2, \ldots, m$ from the definition of the Busemann-Hausdorff area density is an injective linear map (or bijective onto its image). The intersection $B_{l} \cap \mathrm{im}(L)$ is a symmetric polytope of dimension $m$ because $\operatorname{im}(L)=\operatorname{span}\left\{v_{1}, v_{2}, \ldots, v_{m}\right\}$ is an $m$-dimensional plane. Due to Lemma 1.2.16 (ii) we therefore know that the preimage $L^{-1}\left(B_{l}\right)=L^{-1}\left(B_{l} \cap \mathrm{im}(L)\right)$ is a convex body. Notice that if $B_{l} \xrightarrow[l \rightarrow \infty]{\longrightarrow} B$ in $\left(C_{b}(V), d_{\mathcal{H}, V}\right)$ then $L^{-1}\left(B_{l}\right) \xrightarrow[l \rightarrow \infty]{\longrightarrow} L^{-1}(B)$ in $\left(C_{b}\left(\mathbb{R}^{m}\right), d_{\mathcal{H}, \mathbb{R}^{m}}\right)$ because

$$
\begin{aligned}
d_{\mathcal{H}, \mathbb{R}^{m}}\left(L^{-1}\left(B_{l}\right),\right. & \left.L^{-1}(B)\right) \\
& =\max \left(\sup _{x_{1} \in L^{-1}\left(B_{l}\right)} \inf _{x_{2} \in L^{-1}(B)}\left\|x_{1}-x_{2}\right\|_{\mathbb{R}^{m}}, \sup _{x_{2} \in L^{-1}(B)} \inf _{x_{1} \in L^{-1}\left(B_{l}\right)}\left\|x_{1}-x_{2}\right\|_{\mathbb{R}^{m}}\right) \\
& =\max \left(\sup _{y_{1} \in B_{l}} \inf _{y_{2} \in B}\left\|L^{-1}\left(y_{1}-y_{2}\right)\right\|_{\mathbb{R}^{m}}, \sup _{y_{2} \in B} \inf _{y_{1} \in B_{l}}\left\|L^{-1}\left(y_{1}-y_{2}\right)\right\|_{\mathbb{R}^{m}}\right) \\
& \leq\left\|L^{-1}\right\|_{\mathcal{L}\left(\operatorname{im}(L), \mathbb{R}^{m}\right)} \max \left(\sup _{y_{1} \in B_{l}} \inf _{y_{2} \in B}\left\|y_{1}-y_{2}\right\|_{V}, \sup _{y_{2} \in B} \inf _{y_{1} \in B_{l}}\left\|y_{1}-y_{2}\right\|_{V}\right)
\end{aligned}
$$

$$
=\left\|L^{-1}\right\|_{\mathcal{L}\left(\operatorname{im}(L), \mathbb{R}^{m}\right)} d_{\mathcal{H}, V}\left(B_{l}, B\right) \underset{l \rightarrow \infty}{\longrightarrow} 0
$$

As mentioned earlier, the Lebesgue measure $\mathcal{L}^{m}=\mathcal{H}^{m}$ is continuous on the collection of compact convex sets in $\mathbb{R}^{m}$ with respect to the Hausdorff metric (see [Val68, Satz 12.7, p. 153] or [Bee74, p. 64]). Thus, for any $\sigma \in G C_{m}(V)$ we obtain

$$
\begin{equation*}
A_{l}^{b h}(\sigma)=\frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(L_{\sigma}^{-1}\left(B_{l}\right)\right)} \xrightarrow[l \rightarrow \infty]{ } \frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(L_{\sigma}^{-1}(B)\right)}=A^{b h}(\sigma) \tag{2.2.18}
\end{equation*}
$$

as $l$ tends to infinity.
Let us assume that the convexity of the densities $A_{l}^{b h}$ has already been proved. With (2.2.18) at hand we will conclude the convexity of the original Busemann-Hausdorff area density $A^{b h}$ by using Lemma 2.2.4. Let $P \in G_{m}(V)$ be a fixed plane in $V$ and $\sigma \in G C_{m}(V)$ an arbitrary simple $m$-vector. Then the convexity of each $A_{l}^{b h}$ yields a sequence of calibrating $m$-forms $\left(w_{l}\right)_{l \in \mathbb{N}} \subset\left(\bigwedge^{m}(V)\right)^{*}$ for $P$ such that

$$
\begin{equation*}
\left|\omega_{l}(\sigma)\right| \leq A_{l}^{b h}(\sigma) \tag{2.2.19}
\end{equation*}
$$

with equality if $\sigma \in G C_{m}(V) \cap \bigwedge^{m}(P)$. The convergence of the sequence $\left(A_{l}^{b h}(\sigma)\right)_{l \in \mathbb{N}}$ with limit $A^{b h}(\sigma)$ has been established in (2.2.18). Hence, we find a number $N \in \mathbb{N}$ such that

$$
\left|\omega_{l}(\sigma)\right| \leq A_{l}^{b h}(\sigma)<A^{b h}(\sigma)+1
$$

for all $l>N$. Therefore, the sequence $\left(w_{l}(\sigma)\right)_{l \in \mathbb{N}}$ is bounded in $\mathbb{R}$ and by sequential compactness there is a convergent subsequence $\left(w_{l_{k}}(\sigma)\right)_{k \in \mathbb{N}}$. Passing to subsequences and renaming them, we find that for fixed $\sigma \in \Lambda^{m} V$ both $\left(A_{l}^{b h}(\sigma)\right)_{l \in \mathbb{N}}$ and $\left(\omega_{l}(\sigma)\right)_{l \in \mathbb{N}}$ converge and fulfil (2.2.19). Denote the limit of the latter sequence by $\omega(\sigma):=\lim _{l \rightarrow \infty} \omega_{l}(\sigma)$. The assignment $\sigma \rightarrow \omega(\sigma)$ is an $m$-form because it is the pointwise limit of such forms. Indeed, it is the sought-after calibrating form for $A^{b h}$ because it holds that

$$
\begin{aligned}
|\omega(\sigma)|=\lim _{l \rightarrow \infty}\left|\omega_{l}(\sigma)\right| & =\limsup _{l \rightarrow \infty}\left|\omega_{l}(\sigma)\right| \\
& \leq \limsup _{l \rightarrow \infty} A_{l}^{b h}(\sigma)=\lim _{l \rightarrow \infty} A_{l}^{b h}(\sigma)=A^{b h}(\sigma)
\end{aligned}
$$

with equality if $\sigma \in G C_{m}(V) \cap \bigwedge^{m}(P)$. This shows the convexity of the density $A^{b h}$ corresponding to the initial normed space $\left(V,\|\cdot\|_{V}\right)$ under the assumption that the densities $A_{l}^{b h}$ corresponding to the polyhedral unit balls $B_{l}$ are convex already.

Henceforth, Assumption 2.2 .5 shall hold. We leave the general argumentation behind and focus on the case $m=2$. As Lemma 2.2.4 suggests, the convexity of the corresponding two-dimensional Busemann-Hausdorff area density will be proved by explicitly constructing appropriate calibrators for $A^{b h}$ and $m=2$. We will reformulate the inequality Lemma 2.2.4 for a calibrator and prove a result of convex geometry on the two-dimensional plane.

Let us fix a two-dimensional linear subspace $P \in G_{2}(V)$ with basis $\left\{w_{1}, w_{2}\right\}$ and consider its intersection $B \cap P$ with the unit ball. Since $B$ is a symmetric polytope this intersection is a symmetric polygon (a two-dimensional polytope) with vertices $\hat{a}_{1}, \hat{a}_{2}, \ldots, \hat{a}_{2 N} \in V$. Recall from Theorem 1.2.7 that there is a half-space representation so that $B \cap P$ is the intersection of finitely many half-spaces. Let $F^{i} \in V^{*}$ be the linear form that supports $B \cap P$ on the segment $\left[\hat{a}_{i}, \hat{a}_{i+1}\right]=\left\{t \hat{a}_{i_{1}}+(1-t) \hat{a}_{i} \mid t \in[0,1]\right\}$. This means after appropriate scaling that $\left.F^{i}\right|_{B \cap P} \leq 1$ and $\left.F^{i}\right|_{\left[\hat{a}_{i}, \hat{a}_{i+1}\right]}=1$ (for this, we need that $0 \in \operatorname{int}(B \cap P)$, see (1.2.15)). Further, let $L_{P}: \mathbb{R}^{2} \rightarrow V$ be the linear map that takes the standard basis $\left\{e_{1}, e_{2}\right\}$ of $\mathbb{R}^{2}$ to $\left\{w_{1}, w_{2}\right\}$, that is, $L_{P} e_{i}=w_{i}$. Endowed with the standard Euclidean inner product $\mathbb{R}^{2}$ is an inner product space and the standard basis is orthonormal with respect to this inner product. Henceforth, whenever we use the words volume form and orientation we implicitly mean the standard volume form $\omega_{e}=e_{1}^{*} \wedge e_{2}^{*}$ corresponding to the standard basis and its induced standard orientation as introduced in Section 1.1.4.

Consider the preimage $K_{P}:=L_{P}^{-1}(B \cap P)$ of the unit ball $B$ in $\mathbb{R}^{2}$. This set corresponds to the intersection $B \cap P$. Then $K_{P}$ is a convex body due to Lemma 1.2.16 (ii). Furthermore, $K_{P}$ is the intersection of the half-spaces $\left\{x \in \mathbb{R}^{2} \mid f_{P}^{i}(x) \leq 1\right\}$ where $f_{P}^{i}:=F^{i} \circ L_{P} \in\left(\mathbb{R}^{2}\right)^{*}$ and therefore a polytope in $\mathbb{R}^{2}$ (Lemma 1.2.16 (iii)). In addition, according to Lemma 1.2 .16 (iv) $K_{P}$ is symmetric because $B$ is. Thus, $K_{P}=\left[a_{1} a_{2} \ldots a_{2 N}\right] \subset \mathbb{R}^{2}$ is a symmetric polygon with vertices $a_{i}=a_{i}(P) \in \mathbb{R}^{2}$. Suppose that the vertices are enumerated counterclockwise so that the 2 -vectors $a_{i} \wedge a_{j}$ are positively oriented with respect to the standard volume form $\omega_{e}=e_{1}^{*} \wedge e_{2}^{*}$ for $1 \leq i<j \leq N$.

Define a 2-form $\omega \in \bigwedge^{2}\left(V^{*}\right)$ by

$$
\begin{equation*}
\omega=\omega(P):=\pi \sum_{1 \leq i<j \leq N} p_{i} p_{j} F^{i} \wedge F^{j} . \tag{2.2.20}
\end{equation*}
$$

and set the coefficients as the area ratio of the triangle $\Delta 0 a_{i} a_{i+1}=\Delta 0 a_{i}(P) a_{i+1}(P):=\left[0 a_{i} a_{i+1}\right]$ and the entire polygon $K_{P}$, namely

$$
p_{i}=p_{i}(P):=2 \frac{\mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right)}{\mathcal{H}^{2}\left(K_{P}\right)} \in[0,1]
$$

for $i=1,2, \ldots, N$. Note that by the symmetry of $K_{P}$

$$
\mathcal{H}^{2}\left(K_{P}\right)=\sum_{i=1}^{2 N} \mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right)=2 \sum_{i=1}^{N} \mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right)
$$

and thus $\sum_{i=1}^{N} p_{i}=1$.
We are going to prove that $\omega$ is a calibrator for $P$, that is, for all $v_{1} \wedge v_{2} \in G C_{2}(V)$ we need to show that

$$
\begin{equation*}
\left|\omega\left(v_{1} \wedge v_{2}\right)\right| \leq A^{b h}\left(v_{1} \wedge v_{2}\right) \tag{2.2.21}
\end{equation*}
$$

with equality if $v_{1} \wedge v_{2} \in G C_{2}(V) \cap \bigwedge^{2}(P)$.
Consider a simple 2-vector $\sigma=v_{1} \wedge v_{2} \in G C_{2}(V)$ where $v_{1}, v_{2} \in V$ are linearly independent. Denote the two-dimensional plane spanned by these two vectors by $Q:=\operatorname{span}\left\{v_{1}, v_{2}\right\}$. As $\sigma$ runs through $G C_{2}(V)$ the corresponding plane runs through $G_{2}(V)$ (once again, this uses the injectivity of the Plücker embedding). The top exterior power $\bigwedge^{2}(P)$ of the two-dimensional subspace $P$ is one-dimensional and therefore spanned by $w_{1} \wedge w_{2}$. Thus, $v_{1} \wedge v_{2} \in G C_{2}(V) \cap \bigwedge^{2}(P)$ if and only if $v_{1} \wedge v_{2}=c w_{1} \wedge w_{2}$. Recall from Section 1.1.3 on the Plücker embedding $\rho: G C_{2}(V) \rightarrow \mathbb{P}\left(\bigwedge^{2}(V)\right)$ that the latter means

$$
\rho(Q)=\rho\left(\operatorname{span}\left\{v_{1}, v_{2}\right\}\right)=\left[v_{1} \wedge v_{2}\right]_{\sim}=\left[w_{1} \wedge w_{2}\right]_{\sim}=\rho\left(\operatorname{span}\left\{w_{1}, w_{2}\right\}\right)=\rho(P) .
$$

Since the Plücker embedding is injective by Proposition 1.1.9 we have the equivalence $v_{1} \wedge v_{2} \in$ $G C_{2}(V) \cap \bigwedge^{2}(P)$ if and only if $Q=P$.

Corresponding to the plane $Q$ let $L_{Q}: \mathbb{R}^{2} \rightarrow V$ be the linear map such that the standard basis $\left\{e_{1}, e_{2}\right\}$ of $\mathbb{R}^{2}$ maps to $\left\{v_{1}, v_{2}\right\}$, that is, $L_{Q} e_{i}=v_{i}$. The left-hand side of (2.2.21) then reads

$$
\begin{aligned}
\left|\omega\left(v_{1} \wedge v_{2}\right)\right| & =\left|\omega\left(L_{Q}\left(e_{1}\right) \wedge L_{Q}\left(e_{2}\right)\right)\right| \\
& =\pi\left|\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left(F^{i} \wedge F^{j}\right)\left(L_{Q}\left(e_{1}\right) \wedge L_{Q}\left(e_{2}\right)\right)\right| \\
& =\pi\left|\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left(\left(F^{i} \circ L_{Q}\right) \wedge\left(F^{j} \circ L_{Q}\right)\right)\left(e_{1} \wedge e_{2}\right)\right| .
\end{aligned}
$$

Define $f_{Q}^{i}:=F^{i} \circ L_{Q} \in\left(\mathbb{R}^{2}\right)^{*}$. Any unit 2-vector $\sigma$ in the top exterior power $\bigwedge^{2}\left(\mathbb{R}^{2}\right)$ is simple and either has a representation $\sigma=e_{1} \wedge e_{2}$ or $\sigma=-e_{1} \wedge e_{2}$ (see Corollary 1.1.12). Because of the absolute
value and the linearity of $f_{Q}^{i} \wedge f_{Q}^{j} \in \Lambda^{2}\left(\left(\mathbb{R}^{2}\right)^{*}\right) \cong\left(\bigwedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}$ we can rewrite the last sum as

$$
\begin{aligned}
\left|\omega\left(v_{1} \wedge v_{2}\right)\right| & =\pi \sup \left\{\left|\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left(f_{Q}^{i} \wedge f_{Q}^{j}\right)(\sigma)\right| \mid \sigma \in\left\{e_{1} \wedge e_{2}, e_{2} \wedge e_{1}\right\}\right\} \\
& =\pi \sup \left\{\left|\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left(f_{Q}^{i} \wedge f_{Q}^{j}\right)(\sigma)\right| \mid\|\sigma\|_{\wedge^{2}\left(\mathbb{R}^{2}\right)}=1\right\} \\
& =\pi\left\|\sum_{1 \leq i<j \leq N} p_{i} p_{j} f_{Q}^{i} \wedge f_{Q}^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}
\end{aligned}
$$

In the same way as earlier we can show that $K_{Q}:=L_{Q}^{-1}(B \cap P)$ is a symmetric polygon in $\mathbb{R}^{2}$. By definition of the two-dimensional Busemann-Hausdorff area density in (2.1.8) the right-hand side of (2.2.21) reads

$$
A^{b h}\left(v_{1} \wedge v_{2}\right)=\frac{\varepsilon_{2}}{\mathcal{H}^{2}\left(L_{Q}^{-1}(B \cap P)\right)}=\frac{\pi}{\mathcal{H}^{2}\left(K_{Q}\right)}
$$

Thus, to prove that $\omega(P)$ is a calibrator for the fixed plane $P$, we need to show for each plane $Q \in G_{2}(V)$ that

$$
\begin{equation*}
\left\|\sum_{1 \leq i<j \leq n} p_{i}(P) p_{j}(P) f_{Q}^{i} \wedge f_{Q}^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \leq \frac{1}{\mathcal{H}^{2}\left(K_{Q}\right)} \tag{2.2.22}
\end{equation*}
$$

with equality if the plane $Q$ coincides with the fixed plane $P$. Therein all dependencies on $P$ and $Q$ have been explicitly stated and the coefficients are given by

$$
p_{i}(P)=2 \frac{\mathcal{H}^{2}\left(\Delta 0 a_{i}(P) a_{i+1}(P)\right)}{\mathcal{H}^{2}\left(K_{P}\right)}
$$

Note that the preceding discussion can similarly be made for arbitrary $m$. In fact, Burago and Ivanov give a characterisation for the convexity of the higher-dimensional Busemann-Hausdorff area density which is an $m$-dimensional analogue of (2.2.22) with coefficients $\mu_{i_{1} i_{2} \ldots i_{m}}$ on the left-hand side. The crux in the two-dimensional case is that one can choose $\mu_{i j}=p_{i} p_{j}$ so that the subsequent results hold. Burago and Ivanov mention that the most straight-forward generalisation of the two-dimensional construction (subdividing $K_{P}$ into tetrahedrons instead of triangles) does not work, [BI12, Remark 4.2, p. 637].

To show (2.2.22) we prove the following more general statement from convex geometry on the plane.

Theorem 2.2.7 ([BI12, Proposition 2.2, p. 632])
Let $K \subset \mathbb{R}^{2}$ be a symmetric convex polygon and suppose $f^{1}, f^{2}, \ldots, f^{N} \in\left(\mathbb{R}^{2}\right)^{*}$ are linear forms such that $\left.f^{i}\right|_{K} \leq 1$ for all $i=1,2, \ldots, N$ and $p_{1}, p_{2}, \ldots, p_{N} \in[0,1]$ such that $\sum_{i=1}^{N} p_{i}=1$. Then

$$
\left\|\sum_{1 \leq i<j \leq N} p_{i} p_{j} f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \leq \sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \leq \frac{1}{\mathcal{H}^{2}(K)}
$$

In addition, if $K$ is a convex $2 N$-gon with vertices $a_{1}, a_{2}, \ldots, a_{2 N} \in \mathbb{R}^{2}$ and for $i=1,2, \ldots, N$ the functions $f^{i}$ are supporting linear forms of $K$ corresponding to its sides (such that $\left.f^{i}\right|_{\left[a_{i}, a_{i+1}\right]}=1$ ) and $p_{i}=2 \mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right) / \mathcal{H}^{2}(K)$ where $\Delta 0 a_{i} a_{i+1}:=\left[0 a_{i} a_{i+1}\right]$ denotes the triangle with vertices $0, a_{i}, a_{i+1}$, then the above inequalities turn into equalities.

The convexity of the two-dimensional Busemann-Hausdorff area density - Theorem 2.2.2 - is an immediate consequence of the previous result.

Proof of Theorem 2.2.2: Let $P \in G_{2}(V)$ be arbitrary but fixed. Due to the preceding discussion we need only show the inequality $(2.2 .22)$ for $Q \in G_{2}(V)$ where equality holds if $Q=P$. As previously shown $K_{Q}$ is a symmetric $2 N$-gon with vertices $a_{i}=a_{i}(Q)$ and $f_{Q}^{i}$ are supporting linear forms corresponding to its sides. Further, $\sum_{i=1}^{N} p_{i}(P)=1$. Theorem 2.2.7 for $K=K_{Q}$ yields the desired inequality. If, in addition, $Q=P$ note that the coefficient

$$
p_{i}(P)=2 \frac{\mathcal{H}^{2}\left(\Delta 0 a_{i}(P) a_{i+1}(P)\right)}{\mathcal{H}^{2}\left(K_{P}\right)}=2 \frac{\mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right)}{\mathcal{H}^{2}\left(K_{Q}\right)}=p_{i}(Q)
$$

which is the equality case in Theorem 2.2.7. Thus, $\omega=\omega(P)$ as defined in (2.2.20) is a indeed a calibrator for $P$. Since $P \in G_{2}(V)$ was arbitrary, Lemma 2.2.4 finally shows that the two-dimensional Busemann-Hausdorff area density $A^{b h}$ admits a convex extension.

In the rest of this chapter we will prove Theorem 2.2.7 through several elementary lemmata. The next result provides us with a technical identity for the area of a polygon.

Lemma 2.2.8 ([BI12, Lemma 2.3, p. 632])
Let $K=\left[a_{1} a_{2} \ldots a_{2 N}\right]$ be a symmetric $2 N$-gon in $\mathbb{R}^{2}$. For $i=1,2, \ldots, N$ define $v_{i}:=a_{i+1}-a_{i}$. Then

$$
\mathcal{H}^{2}(K)=\sum_{1 \leq i<j \leq N}\left\|v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}=\left\|\sum_{1 \leq i<j \leq N} v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}
$$

Proof: Note that the vertices of $K$ are positively oriented. Then the convexity of the polygon $K$
implies that the vectors $v_{i}$ - which correspond to the edges of $K$ - are also positively oriented, that is, $\omega_{e}\left(v_{i} \wedge v_{j}\right)>0$ for any $1 \leq i<j \leq N$.

Since $K$ is symmetric we observe that $\mathcal{H}^{2}(K)=2 \mathcal{H}^{2}\left(\left[a_{1} a_{2} \ldots a_{N}\right]\right)$. Let us denote by $\Delta a_{1} a_{j} a_{j+1}:=$ $\left[a_{1} a_{j} a_{j+1}\right]$ the triangle with vertices $a_{1}, a_{j}, a_{j+1}$. Note that the intersection of any two of such triangles is a set of two-dimensional Lebesgue measure zero. The polygon $\left[a_{1} a_{2} \ldots a_{N}\right.$ ] is the union of all triangles $\Delta a_{1} a_{j} a_{j+1}$ where $j=2,3, \ldots, N$. Then by the countable additivity of Hausdorff measure

$$
\begin{equation*}
\mathcal{H}^{2}(K)=2 \mathcal{H}^{2}\left(\left[a_{1} a_{2} \ldots a_{N}\right]\right)=2 \sum_{j=2}^{N} \mathcal{H}^{2}\left(\Delta a_{1} a_{j} a_{j+1}\right) . \tag{2.2.23}
\end{equation*}
$$

Observe further that the area of the triangle $\Delta a_{1} a_{j} a_{j+1}$ can be expressed via the norm of the wedge product of its edges due to Proposition 1.1.15. Then

$$
\mathcal{H}^{2}\left(\Delta a_{1} a_{j} a_{j+1}\right)=\frac{1}{2}\left\|\left(a_{j}-a_{1}\right) \wedge\left(a_{j+1}-a_{j}\right)\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}=\frac{1}{2}\left\|\left(a_{j}-a_{1}\right) \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}
$$

By telescoping $a_{j}-a_{1}=\sum_{i=1}^{j-1} a_{i+1}-a_{i}=\sum_{i=1}^{j-1} v_{i}$ and

$$
\mathcal{H}^{2}\left(\Delta a_{1} a_{j} a_{j+1}\right)=\frac{1}{2}\left\|\left(a_{j}-a_{1}\right) \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}=\frac{1}{2}\left\|\sum_{i=1}^{j-1} v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}
$$

We apply Corollary 1.1.20 to the preceding identity and substitute the result into (2.2.23) which finally yields

$$
\begin{aligned}
\mathcal{H}^{2}(K) & =2 \sum_{j=2}^{N} \mathcal{H}^{2}\left(\Delta a_{1} a_{j} a_{j+1}\right)=2 \sum_{j=2}^{N} \frac{1}{2}\left\|\sum_{i=1}^{j-1} v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)} \\
& =\sum_{j=2}^{N} \sum_{i=1}^{j-1}\left\|v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)} \\
& =\sum_{1 \leq i<j \leq N}\left\|v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)} .
\end{aligned}
$$

The following lemma takes care of the equality case in Theorem 2.2.7.

Lemma 2.2.9 ([BI12, Lemma 2.4, p. 633])
Let $K=\left[a_{1} a_{2} \ldots a_{2 N}\right]$ be a symmetric $2 N$-gon in $\mathbb{R}^{2}$. For $i=1,2, \ldots, N$ define $v_{i}:=a_{i+1}-a_{i}$ and $p_{i}:=2 \mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right) / \mathcal{H}^{2}(K)$ and let $f^{i} \in\left(\mathbb{R}^{2}\right)^{*}$ be such that $\left.f^{i}\right|_{K} \leq 1$ and $\left.f^{i}\right|_{\left[a_{i}, a_{i+1}\right]}=1$. Then

$$
p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\Lambda^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}}\left\|v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}
$$

for all $i, j=1,2, \ldots, N$ and therefore

$$
\left\|\sum_{1 \leq i<j \leq N} p_{i} p_{j} f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\frac{1}{\mathcal{H}^{2}(K)}
$$

Proof: Define $S_{i}:=2 \mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right)>0$. Then by Corollary 1.1.20 the area of the parallelogram spanned by $a_{i}$ and $a_{i+1}$ can be expressed via $S_{i}=\left\|a_{i} \wedge a_{i+1}\right\|_{\wedge^{2}\left(\mathbb{R}^{2}\right)}$. Of course, $p_{i}=S_{i} / \mathcal{H}^{2}(K)$. We recall from Section 1.1.4 that the standard area form $\omega=\omega_{e}$ induces a Riesz-type isomorphism $\iota_{\omega}: \mathbb{R}^{2} \rightarrow\left(\mathbb{R}^{2}\right)^{*}$ given by (1.1.8). We will now show that $\iota_{\omega}\left(v_{i}\right)=S_{i} f^{i}$. Note that there is no implied summation on $i$ in this formula. This is an equation in the dual space $\left(\mathbb{R}^{2}\right)^{*}$, so we need to show $\iota_{\omega}\left(v_{i}\right)(u)=\omega\left(u \wedge v_{i}\right)=S_{i} f^{i}(u)$ for all $u \in \mathbb{R}^{2}$. Moreover, it suffices to prove this for a basis of $\mathbb{R}^{2}$ because both sides of the equation are linear functions. Observe that since $K$ is convex, $a_{i}$ and $v_{i}=a_{i+1}-a_{i}$ are linearly independent vectors and hence form a basis of $\mathbb{R}^{2}$. Indeed, (again there is no implied summation over $i$ here)

$$
S_{i} f^{i}\left(v_{i}\right)=S_{i} f^{i}\left(a_{i+1}-a_{i}\right)=S_{i}\left(f^{i}\left(a_{i+1}\right)-f\left(a_{i}\right)\right)=S_{i}(1-1)=0=\omega(0)=\omega\left(v_{i} \wedge v_{i}\right)
$$

because $f^{i}$ is linear and $\left.f^{i}\right|_{\left[a_{i}, a_{i+1}\right]}=1$. Further, by (1.1.7) calculate

$$
S_{i} f^{i}\left(a_{i}\right)=S_{i}=\left\|a_{i} \wedge a_{i+1}\right\|_{\wedge^{2}\left(\mathbb{R}^{2}\right)}=\omega\left(a_{i} \wedge a_{i+1}\right)=\omega\left(a_{i} \wedge\left(a_{i+1}-a_{i}\right)\right)=\omega\left(a_{i} \wedge v_{i}\right)
$$

where the third equality uses that the 2 -vectors $a_{i} \wedge a_{i+1}$ are positively oriented. The isometric property of $\iota_{\omega}$ that was established in Proposition 1.1.18 yields

$$
\begin{aligned}
p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} & =\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}} S_{i} S_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \\
& =\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}}\left\|\left(S_{i} f^{i}\right) \wedge\left(S_{j} f^{j}\right)\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \\
& =\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}}\left\|\iota_{\omega}\left(v_{i}\right) \wedge \iota_{\omega}\left(v_{j}\right)\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \\
& =\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}}\left\|v_{i} \wedge v_{i}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}
\end{aligned}
$$

The second identity follows from Lemma 2.2.8 because

$$
\begin{aligned}
\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} & =\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}} \sum_{1 \leq i<j \leq N}\left\|v_{i} \wedge v_{i}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)} \\
& =\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}} \mathcal{H}^{2}(K)
\end{aligned}
$$

$$
=\frac{1}{\mathcal{H}^{2}(K)}
$$

Recall the Riesz-type isomorphism $\iota_{\omega}$ induces a dual volume form $\omega^{*}$ given by (1.1.11). The edges $v_{i}$ of the polygon $K$ are consistently (positively) oriented with respect to $\omega$ and therefore the corresponding supporting linear forms $f^{i}$ are consistently (positively) oriented with respect to $\omega^{*}$ because

$$
\omega^{*}\left(f^{i} \wedge f^{j}\right)=\frac{1}{S_{i} S_{j}} \omega\left(\left(\iota_{\omega}\right)^{-1}\left(S_{i} f^{i}\right) \wedge\left(\iota_{\omega}\right)^{-1}\left(S_{j} f^{j}\right)\right)=\frac{1}{S_{i} S_{j}} \omega\left(v_{i} \wedge v_{j}\right)>0
$$

for $1 \leq i<j \leq N$, wherein $S_{i}>0$. Finally, by Corollary 1.1.20

$$
\left\|\sum_{1 \leq i<j \leq N} p_{i} p_{j} f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\frac{1}{\mathcal{H}^{2}(K)}
$$

It remains to show the inequality part of Theorem 2.2.7. The next lemma takes care of the principal case where the linear forms $f^{i}$ support the polygon $K$ at its edges.

Lemma 2.2.10 ([BI12, Lemma 2.5, p. 633])
Let $K=\left[a_{1} a_{2} \ldots a_{2 N}\right]$ be a symmetric $2 N$-gon in $\mathbb{R}^{2}$. For $i=1,2, \ldots, N$ let $f^{i} \in\left(\mathbb{R}^{2}\right)^{*}$ be such that $\left.f^{i}\right|_{K} \leq 1$ and $\left.f^{i}\right|_{\left[a_{i}, a_{i+1}\right]}=1$. Let $p_{1}, p_{2}, \ldots, p_{N} \in[0,1]$ be nonnegative real numbers such that $\sum_{i=1}^{N} p_{i}=1$. Then

$$
\left\|\sum_{1 \leq i<j \leq N} p_{i} p_{j} f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \leq \frac{1}{\mathcal{H}^{2}(K)}
$$

Proof: The first identity follows because all the 2-forms $f^{i} \wedge f^{j}$ are of the same orientation for $i<j$ by the same argumentation as used in the proof of Lemma 2.2.9. Again for $i=1,2, \ldots, N$ define the vectors $v_{i}:=a_{i+1}-a_{i}$ corresponding to the edges of $K$. Further, let $q_{i}:=2 \mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right) / \mathcal{H}^{2}(K)$ and $\lambda_{i}:=p_{i} / q_{i}$ for $i=1,2, \ldots, N$. Due to the first assertion of Lemma 2.2.9

$$
q_{i} q_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}}\left\|v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)}
$$

Define the vectors $v_{i}^{\prime}:=\lambda_{i} v_{i}$. By Lemma 1.2.15 there is a symmetric $2 N$-gon $K^{\prime}=\left[a_{1}^{\prime} a_{2}^{\prime} \ldots a_{2 N}^{\prime}\right]$ such that $a_{i+1}^{\prime}-a_{i}^{\prime}=v_{i}^{\prime}$ for $i=1,2, \ldots, N$. Then

$$
\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\sum_{1 \leq i<j \leq N} \lambda_{i} \lambda_{j} q_{i} q_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}
$$

$$
\begin{aligned}
& =\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}} \sum_{1 \leq i<j \leq N} \lambda_{i} \lambda_{j}\left\|v_{i} \wedge v_{j}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)} \\
& =\frac{1}{\left(\mathcal{H}^{2}(K)\right)^{2}} \sum_{1 \leq i<j \leq N}\left\|v_{i}^{\prime} \wedge v_{j}^{\prime}\right\|_{\Lambda^{2}\left(\mathbb{R}^{2}\right)} .
\end{aligned}
$$

Applying Lemma 2.2.8 to $K^{\prime}$ and substituting the result into the last equation yields

$$
\begin{equation*}
\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\frac{\mathcal{H}^{2}\left(K^{\prime}\right)}{\left(\mathcal{H}^{2}(K)\right)^{2}} \tag{2.2.24}
\end{equation*}
$$

Note that the 1-dimensional Lebesgue measure of the edges satisfy

$$
\mathcal{L}^{1}\left(\left[a_{i}^{\prime}, a_{i+1}^{\prime}\right]\right)=\left\|v_{i}^{\prime}\right\|_{\mathbb{R}^{2}}=\lambda_{i}\left\|v_{i}\right\|_{\mathbb{R}^{2}}=\lambda_{i} \mathcal{L}^{1}\left(\left[a_{i}, a_{i+1}\right]\right)
$$

by construction of $K^{\prime}$. Consider the support function $h_{K}$ of the convex set $K$. By virtue of Proposition 1.2.6 the value $h_{i}:=h_{K}\left(\frac{f^{i}}{\left\|f^{i}\right\|_{\left(\mathbb{R}^{2}\right)^{*}}}\right)$ is the distance of $0 \in K$ to the edge $\left[a_{i}, a_{i+1}\right]$ corresponding to $f^{i}$. Then the two-dimensional Hausdorff measure of the triangle $\Delta 0 a_{i} a_{i+1}$ is given by

$$
\mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right)=\frac{1}{2} h_{i} \mathcal{L}^{1}\left(\left[a_{i}, a_{i+1}\right]\right) .
$$

Now we can calculate

$$
\begin{align*}
1=\sum_{i=1}^{N} p_{i} & =\sum_{i=1}^{N} \lambda_{i} q_{i}=\frac{2}{\mathcal{H}^{2}(K)} \sum_{i=1}^{N} \lambda_{i} \mathcal{H}^{2}\left(\Delta 0 a_{i} a_{i+1}\right) \\
& =\frac{2}{\mathcal{H}^{2}(K)} \frac{1}{2} \sum_{i=1}^{N} \lambda_{i} h_{i} \mathcal{L}^{1}\left(\left[a_{i}, a_{i+1}\right]\right)  \tag{2.2.25}\\
& =\frac{1}{\mathcal{H}^{2}(K)} \frac{1}{2} \sum_{i=1}^{2 N} h_{i} \mathcal{L}^{1}\left(\left[a_{i}^{\prime}, a_{i+1}^{\prime}\right]\right)
\end{align*}
$$

where we used the symmetry of $K$ to extend the sum in the last equality. The explicit formula for mixed volume in Lemma 1.2.14 (use $A=K^{\prime}$ and $K=K$ ) shows that

$$
\frac{1}{2} \sum_{i=1}^{2 N} h_{i} \mathcal{L}^{1}\left(\left[a_{i}^{\prime}, a_{i+1}^{\prime}\right]\right)=\mathrm{V}\left(K^{\prime}, K\right)
$$

The preceding identity (2.2.25) thus transforms to

$$
\mathcal{H}^{2}(K)=\mathrm{V}\left(K^{\prime}, K\right)
$$

Now we use the Minkowski inequality for convex bodies (Proposition 1.2.12) to see that

$$
\left(\mathcal{H}^{2}(K)\right)^{2}=\left(\mathrm{V}\left(K^{\prime}, K\right)\right)^{2} \geq \mathcal{H}^{2}\left(K^{\prime}\right) \mathcal{H}^{2}(K)
$$

or, equivalently, $\mathcal{H}^{2}(K) \geq \mathcal{H}^{2}\left(K^{\prime}\right)$. Using this estimate on (2.2.24) proves the assertion of the lemma.

Finally, we can complete the proof of Theorem 2.2.7 which in turn implies the convexity of the two-dimensional Busemann-Hausdorff area density as already shown. Thus, what follows concludes this chapter.

Proof of Theorem 2.2.7: By means of Lemma 2.2.9 and Lemma 2.2.10 it remains to show the inequality for the case where $K$ is a symmetric polygon (not necessarily with $2 N$ vertices) and the linear forms $f^{1}, f^{2}, \ldots, f^{N} \in\left(\mathbb{R}^{2}\right)^{*}$ are such that $\left.f^{i}\right|_{K} \leq 1$. The triangle inequality for $\|\cdot\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}$ yields

$$
\left\|\sum_{1 \leq i<j \leq N} p_{i} p_{j} f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \leq \sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}
$$

Thus, we need to show that

$$
\begin{equation*}
\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \leq \frac{1}{\mathcal{H}^{2}(K)} \tag{2.2.26}
\end{equation*}
$$

Consider the left hand side of $(2.2 .26)$ as a function in just one variable $f^{k}$ with the remaining $f^{i}$ for $i \neq k$ staying fixed and let us call this function $G_{k}$. Then $G_{k}$ is a convex function because it is the sum of the norms of 2-forms. Recall the definition of the polar set $K^{\circ}$ from (1.2.16). Then $f^{i} \in K^{\circ}$. Furthermore, the polar set of a polygon is itself a polygon by Proposition 1.2.9. This means that $G_{k}$ is a convex function mapping the polar polygon $K^{\circ}$ to $\mathbb{R}$. In addition, $G_{k}$ is bounded above because $K^{\circ}$ is bounded as a polytope and for each of the summands

$$
\begin{aligned}
\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}^{2} & =\operatorname{det}\left(\begin{array}{cc}
\left\langle f_{i}, f_{i}\right\rangle_{\left(\mathbb{R}^{2}\right)^{*}} & \left\langle f_{i}, f_{j}\right\rangle_{\left(\mathbb{R}^{2}\right)^{*}} \\
\left\langle f_{j}, f_{i}\right\rangle_{\left(\mathbb{R}^{2}\right)^{*}} & \left\langle f_{j}, f_{j}\right\rangle_{\left(\mathbb{R}^{2}\right)^{*}}
\end{array}\right) \\
& =\left\|f_{i}\right\|_{\left(\mathbb{R}^{2}\right)^{*}}^{2}\left\|f_{j}\right\|_{\left(\mathbb{R}^{2}\right)^{*}}^{2}-2\left\langle f_{i}, f_{j}\right\rangle_{\left(\mathbb{R}^{2}\right)^{*}} \\
& \leq\left\|f_{i}\right\|_{\left(\mathbb{R}^{2}\right)^{*}}^{2}\left\|f_{j}\right\|_{\left(\mathbb{R}^{2}\right)^{*}}^{2}+2\left\|f_{i}\right\|_{\left(\mathbb{R}^{2}\right)^{*}}\left\|f_{j}\right\|_{\left(\mathbb{R}^{2}\right)^{*}} \\
& =\left(\left\|f_{i}\right\|_{\left(\mathbb{R}^{2}\right)^{*}}+\left\|f_{j}\right\|_{\left(\mathbb{R}^{2}\right)^{*}}\right)^{2}<(2 R)^{2}
\end{aligned}
$$

where $R>0$ such that $K^{\circ}$ is contained in the ball $B_{R}(0) \subset\left(\mathbb{R}^{2}\right)^{*}$ of radius $R$. Proposition 1.2.18
shows that $G_{k}$ attains its maximum on a vertex of $K^{\circ}$. But by Lemma 1.2.10 the vertices of $K^{\circ}$ are the linear forms that support $K$ at its edges. Thus, it suffices to show (2.2.26) for the case where the $f^{i}$ are (possibly duplicate) linear forms supporting $K$ at its edges. Now if any two of the functions $f^{i}$ and $f^{j}$ coincide (without loss of generality $f_{1}=f_{N}$ ) the problem can be reduced to a smaller number of functions by proceeding as follows. Calculate on the left hand side of (2.2.26)

$$
\begin{aligned}
\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} & =\sum_{i=1}^{N-1} \sum_{j=i+1}^{N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \\
& =\sum_{j=2}^{N} p_{1} p_{j}\left\|f^{1} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}+\sum_{i=2}^{N-1} \sum_{j=i+1}^{N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}
\end{aligned}
$$

By expanding the second sum and using $f^{1} \wedge f^{N}=0$ on the first sum this reduces to

$$
\begin{aligned}
\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}= & \sum_{j=2}^{N-1} p_{1} p_{j}\left\|f^{1} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \\
& +\sum_{i=2}^{N-1}\left(\sum_{j=i+1}^{N-1} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}+p_{i} p_{N}\left\|f^{i} \wedge f^{N}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}\right) \\
= & \sum_{j=2}^{N-1} p_{1} p_{j}\left\|f^{1} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}+\sum_{i=2}^{N-1} p_{i} p_{N}\left\|f^{i} \wedge f^{N}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \\
& +\sum_{i=2}^{N-2} \sum_{j=i+1}^{N-1} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}
\end{aligned}
$$

Again we can use that $f^{1}=f^{N}$ and thus,

$$
\left\|f^{1} \wedge f^{i}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\left\|-f^{i} \wedge f^{N}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}=\left\|f^{i} \wedge f^{N}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}
$$

Renaming the indices in the first two sums then yields

$$
\begin{aligned}
\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} & =\sum_{j=2}^{N-1}\left(p_{1}+p_{N}\right) p_{j}\left\|f^{1} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}+\sum_{i=2}^{N-2} \sum_{j=i+1}^{N-1} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \\
& =\sum_{i=1}^{N-2} \sum_{j=i+1}^{N-1} p_{i}^{\prime} p_{j}^{\prime}\left\|f^{i} \wedge f^{j}\right\|_{\left(\wedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}}
\end{aligned}
$$

where $p_{1}^{\prime}:=p_{1}+p_{N}$ and $p_{i}^{\prime}:=p_{i}$ for $i>1$. That is, the problem can be reduced to a smaller number of functions by dropping $f_{N}$ from the list of functions and replacing the set of coefficients by $p_{1}+p_{N}, p_{2}, \ldots, p_{N-1}$. Additionally, note that changing the sign of one of the functions $f^{i}$ does not
change the left hand side of (2.2.26). Therefore, it further suffices to assume that all the functions $\pm f_{1}, \pm f_{2}, \ldots, \pm f_{N}$ are distinct.

If $N=1$ then the left hand side of (2.2.26) equals zero and the inequality holds trivially. Therefore, suppose $N>1$ and consider the polyhedral set given by the half-space representation

$$
K^{\prime}:=\bigcap_{i=1}^{N}\left\{x \in \mathbb{R}^{2}| | f^{i}(x) \mid \leq 1\right\} .
$$

Since $K$ is symmetric and $\left.f^{i}\right|_{K} \leq 1$ by hypothesis, we know that $f^{i}(x) \leq 1$ and $-f^{i}(x)=f^{i}(-x) \leq 1$ for $i=1,2, \ldots, N$ and any $x \in K$. Therefore, $\left|f^{i}(x)\right| \leq 1$ for $i=1,2, \ldots, N$ and any $x \in K$. In other words, $K \subset K^{\prime}$. Therefore,

$$
\begin{equation*}
\mathcal{H}^{2}(K) \leq \mathcal{H}^{2}\left(K^{\prime}\right) \tag{2.2.27}
\end{equation*}
$$

Observe that $K^{\prime}$ is a symmetric polyhedral set and that $\pm f_{1}, \pm f_{2}, \ldots, \pm f_{N}$ are those linear forms which give the half-space representation. By our assumption they correspond to $N$ distinct pairs of opposing supporting hyperplanes. Since $N \geq 2$ there are at least two distinct pairs of such opposing supporting hyperplanes. Every polyhedral set in the two-dimensional space $\mathbb{R}^{2}$ that is supported by at least two distinct pairs of opposing hyperplanes is consequently bounded. Thus, $K^{\prime}$ is a bounded, polyhedral set which means by Theorem 1.2 .7 that $K^{\prime}$ is a polygon. In two dimensions the number of vertices and facets (which are its edges) coincide (see [Brø83, §16 Euler's Relation, Theorem 16.1, p. 98]). Thus, $K^{\prime}$ is a $2 N$-gon. By definition of $K^{\prime}$ through its half-space representation, we see that $f^{i}$ are linear forms that support $K^{\prime}$ at its edges. Then all the prerequisites to apply Lemma 2.2.10 to $K^{\prime}$ are fulfilled. Using this fact and (2.2.27) finally proves (2.2.26) because

$$
\sum_{1 \leq i<j \leq N} p_{i} p_{j}\left\|f^{i} \wedge f^{j}\right\|_{\left(\bigwedge^{2}\left(\mathbb{R}^{2}\right)\right)^{*}} \leq \frac{1}{\mathcal{H}^{2}\left(K^{\prime}\right)} \leq \frac{1}{\mathcal{H}^{2}(K)}
$$

## Chapter 3

## The Plateau problem in arbitrary codimension in the Finsler setting

In this final chapter we formulate and solve the Plateau problem in $n$-dimensional Finsler space $\left(\mathbb{R}^{n}, F\right)$ for a reversible Finsler metric $F$. To achieve this goal, we first reference the work of Hildebrandt and von der Mosel in which they developed the theory of Cartan functionals. We state their result on the Plateau problem for Cartan integrands (Theorem 3.2.1). Finally, we solve the Plateau problem in the Finsler setting by identifying the Busemann-Hausdorff area integrand $a_{m}^{F}$ as a Cartan integrand for $m=2$. For that, we use the convexity of the area integrand which we established in Chapter 2 . Further, we use a representation of the area integrand, found by Overath in [Ove13].

### 3.1 Formulation of the Plateau problem in Finsler space

In this section we formalise the Plateau problem which was mentioned in the introduction. Let us recall the question it poses:

Given a closed rectifiable Jordan curve $\Gamma$, is there a minimal surface spanned by $\Gamma$ ?

For the following, denote the two-dimensional Euclidean unit ball in $\mathbb{R}^{2}$ by $B$. Its boundary is the set of vectors of Euclidean unit length and we denote it by $\mathbb{S}^{1}$.

A curve in $\mathbb{R}^{n}$ is a continuous mapping $c:[a, b] \rightarrow \mathbb{R}^{n}$ of an interval $[a, b]$ into $\mathbb{R}^{n}$. If a curve is injective then we call it a simple curve. A curve $c$ is said to be a closed curve if $c(a)=c(b)$. We can reparametrise a curve to be a continuous mapping $c: \mathbb{S}^{1} \rightarrow \mathbb{R}^{n}$ (Consider the map $c \circ \varphi^{-1}$ where $\left.\varphi:[a, b] \rightarrow \mathbb{S}^{1}, t \mapsto\left(\cos \left(2 \pi \frac{t-a}{b-a}\right), \sin \left(2 \pi \frac{t-a}{b-a}\right)\right)\right)$. The image $\Gamma=\operatorname{im}(c)$ of a simple closed curve $c: \mathbb{S}^{1} \rightarrow \mathbb{R}^{n}$ is called a Jordan curve. A rectifiable curve is a curve of finite Euclidean length,
that is,

$$
\mathcal{L}^{|\cdot|}(c)=\int_{\mathbb{S}^{1}}|\dot{c}(t)| d t<\infty
$$

Suppose $\Gamma$ is a closed Jordan curve in $\mathbb{R}^{n}$ and let $\varphi: \mathbb{S}^{1} \rightarrow \Gamma$ be a homeomorphism from $\mathbb{S}^{1}$ onto $\Gamma$. Then a continuous mapping $\psi: \mathbb{S}^{1} \rightarrow \Gamma$ from $\mathbb{S}^{1}$ onto $\Gamma$ is said to be weakly monotonic if there is a non-decreasing continuous function $\tau:[0,2 \pi] \rightarrow \mathbb{R}$ with $\tau(0)=0, \tau(2 \pi)=2 \pi$ such that

$$
\psi(\cos (\theta), \sin (\theta))=\varphi(\cos (\tau(\theta)), \sin (\tau(\theta)))
$$

for $0 \leq \theta \leq 2 \pi$ (see [DHKW92, Definition 2, p. 231]).
We now define the class of admissible surfaces for the Plateau problem. Recall that every function $X \in W^{1,2}\left(B, \mathbb{R}^{n}\right)$ has a trace $\left.X\right|_{\mathbb{S}^{1}}$ on the boundary $\mathbb{S}^{1}$ which is of class $L^{2}\left(\mathbb{S}^{1}, \mathbb{R}^{n}\right)$ (see [Eva10, Theorem 1, p. 272])

Given a closed Jordan curve $\Gamma$ in $\mathbb{R}^{n}$, a mapping $X: B \rightarrow \mathbb{R}^{n}$ is said to be of class $\mathcal{C}(\Gamma)$ if and only if $X \in W^{1,2}\left(B, \mathbb{R}^{n}\right)$ and its trace $\left.X\right|_{\mathbb{S}^{1}}$ can be represented by a weakly monotonic, continuous mapping $\varphi: \mathbb{S}^{1} \rightarrow \Gamma$ of $\mathbb{S}^{1}$ onto $\Gamma$ (which means, every $L^{2}\left(\mathbb{S}^{1}, \mathbb{R}^{n}\right)$-representative of $\left.X\right|_{\mathbb{S}^{1}}$ coincides with $\varphi$ except for a subset of zero one-dimensional Hausdorff measure). One can show that $\mathcal{C}(\Gamma)$ is non-empty if the Jordan curve $\Gamma$ is rectifiable (see [DHKW92, pp. 233-234]).

Finally, we can formulate the Plateau problem for Finsler area mentioned above. Notice that for $n=3$ this coincides with the Plateau problem in [OvdM14, Theorem 1.2, p. 278].

Theorem 3.1.1 (Plateau problem for Finsler area in arbitrary codimension )
Let $F$ be a reversible Finsler metric on $\mathbb{R}^{n}$ and assume in addition that

$$
0<m_{F}:=\inf _{\mathbb{R}^{n} \times \mathbb{S}^{n-1}} F(\cdot, \cdot) \leq \sup _{\mathbb{R}^{n} \times \mathbb{S}^{n-1}} F(\cdot, \cdot)=\|F\|_{L^{\infty}\left(\mathbb{R}^{n} \times \mathbb{S}^{n-1}\right)}=: M_{F}<\infty .
$$

Then for any given rectifiable Jordan curve $\Gamma \subset \mathbb{R}^{n}$ there exists a surface $X \in \mathcal{C}(\Gamma)$, such that

$$
\operatorname{area}_{B}^{F}(X)=\inf _{\mathcal{C}(\Gamma)} \operatorname{area}_{B}^{F}(\cdot)
$$

In addition, the minimiser $X$ is Euclidean conformally parametrised almost everywhere on $B$, that is,

$$
\left|\frac{\partial X}{\partial u^{1}}\right|^{2}=\left|\frac{\partial X}{\partial u^{2}}\right|^{2} \quad \text { and } \quad\left\langle\frac{\partial X}{\partial u^{1}}, \frac{\partial X}{\partial u^{2}}\right\rangle_{\mathbb{R}^{2}}=0 \quad \mathcal{H}^{2}-\text { a.e. on } B .
$$

Furthermore, the minimiser $X$ has the following regularity property,

$$
X \in C^{0}\left(\bar{B}, \mathbb{R}^{n}\right) \cap C^{0, \sigma}\left(B, \mathbb{R}^{n}\right) \cap W^{1, q}\left(B, \mathbb{R}^{n}\right)
$$

for $\sigma:=\left(\frac{m_{F}}{M_{F}}\right)^{2} \in(0,1]$ and some $q>2$.

### 3.2 Cartan functional theory

In the work of Hildebrandt and von der Mosel [HvdM99, HvdM02, HvdM03a, HvdM03b, HvdM03c, HvdM05, HvdM09] the theory of Cartan (or parametric) functionals has been developed. In particular, the Plateau problem for Cartan integrands has been solved and higher regularity of the minimizers has been established. We want to apply their theory to the Busemann-Hausdorff area integrand and prove Theorem 3.1.1 from the preceding section.

First, let us define the basics of Cartan functional theory. For two integers $n, m$ with $n \geq m$, set $N:=$ $\binom{n}{m}$. Recall from Corollary 1.1 .13 that we can isometrically identify the spaces $\left(\Lambda^{m}\left(\mathbb{R}^{n}\right),\|\cdot\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)}\right)$ and $\left(\mathbb{R}^{N},|\cdot|\right)$ where $|\cdot|$ is the standard Euclidean norm on $\mathbb{R}^{N}$. We denote the image of the set $G C_{m}\left(\mathbb{R}^{n}\right) \subset \bigwedge^{m}\left(\mathbb{R}^{n}\right)$ under this identification as $\mathbb{G} \mathbb{C}_{m}\left(\mathbb{R}^{n}\right) \subset \mathbb{R}^{N}$.

A function $I: \mathbb{R}^{n} \times \mathbb{R}^{N} \rightarrow \mathbb{R}$ is called a Cartan integrand if it is continuous, that is,

$$
\begin{equation*}
I \in C^{0}\left(\mathbb{R}^{n} \times \mathbb{R}^{N}\right) \tag{R}
\end{equation*}
$$

and if it is homogeneous of degree one in its second variable, that is,

$$
\begin{equation*}
I(x, t z)=t I(x, z) \quad \text { for all } t>0,(x, z) \in \mathbb{R}^{n} \times \mathbb{R}^{N} \tag{H}
\end{equation*}
$$

A Cartan integrand $I$ is said to be positive definite if there are two constants $M_{1}$ and $M_{2}$ with $0<M_{1} \leq M_{2}$ such that

$$
\begin{equation*}
M_{1}|z| \leq I(x, z) \leq M_{2}|z| \quad \text { for all }(x, z) \in \mathbb{R}^{n} \times \mathbb{G C}_{m}\left(\mathbb{R}^{n}\right) \tag{D}
\end{equation*}
$$

A Cartan integrand $I$ is said to be semi-elliptic on $\Omega \times \mathbb{R}^{N}$ for $\Omega \subset \mathbb{R}^{n}$ if it is convex in its second variable, that is, if

$$
\begin{equation*}
I\left(x, t z_{1}+(1-t) z_{2}\right) \leq t I\left(x, z_{1}\right)+(1-t) I\left(x, z_{2}\right) \quad \text { for all } x \in \Omega, z_{1}, z_{2} \in \mathbb{R}^{N} \text { and } t \in[0,1] . \tag{C}
\end{equation*}
$$

Given a Cartan integrand $I$ we can define the Cartan functional $\mathcal{I}$ in the following way. Suppose $\mathcal{M}$ is a smooth $m$-manifold and $X: \mathcal{M} \rightarrow \mathbb{R}^{n}$ a smooth immersion. Then we define

$$
\mathcal{I}(X):=\int_{p \in \mathcal{M}} i(p)
$$

where the differential $m$-form $i$ on $\mathcal{M}$ is given in local coordinates $\left(u^{1}, u^{2}, \ldots, u^{m}\right): U \subset \mathcal{M} \rightarrow \Omega \subset \mathbb{R}^{n}$
by

$$
i(p):=I\left(X(p), d X_{p}\left(\left.\frac{\partial}{\partial u^{1}}\right|_{p}\right) \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{2}}\right|_{p}\right) \wedge \cdots \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{m}}\right|_{p}\right)\right) .
$$

This form is well-defined globally due to its invariance under a coordinate change ([Ove13, p. 47]). If we choose $\mathcal{M}=\Omega$ for an open subset $\Omega \subset \mathbb{R}^{m}$ and $\left(u^{\alpha}\right)$ to be the standard coordinates on $\mathbb{R}^{m}$, then $\mathcal{I}(X)$ computes to

$$
\mathcal{I}(X)=\int_{u \in \Omega} I\left(X(u), \frac{\partial X}{\partial u^{1}}(u) \wedge \frac{\partial X}{\partial u^{2}}(u) \wedge \cdots \wedge \frac{\partial X}{\partial u^{m}}(u)\right) d u^{1} d u^{2} \ldots d u^{m} .
$$

The value $\mathcal{I}(X)$ can be guaranteed to be finite by assuming $X \in W^{1, m}\left(\Omega, \mathbb{R}^{n}\right)$ and $I$ to satisfy the positive definiteness relation (D). We can see this by combining these assumptions, Lemma 1.1.14 and the identification $\mathbb{R}^{N} \cong \bigwedge^{m}\left(\mathbb{R}^{n}\right)$ to get the estimate

$$
\begin{aligned}
I\left(X(u), \frac{\partial X}{\partial u^{1}}(u) \wedge \frac{\partial X}{\partial u^{2}}(u) \wedge \cdots \wedge \frac{\partial X}{\partial u^{m}}(u)\right) & \leq M_{2}\left|\frac{\partial X}{\partial u^{1}}(u) \wedge \frac{\partial X}{\partial u^{2}}(u) \wedge \cdots \wedge \frac{\partial X}{\partial u^{m}}(u)\right| \\
& =M_{2}\left\|\frac{\partial X}{\partial u^{1}}(u) \wedge \frac{\partial X}{\partial u^{2}}(u) \wedge \cdots \wedge \frac{\partial X}{\partial u^{m}}(u)\right\|_{\wedge^{m}\left(\mathbb{R}^{n}\right)} \\
& \leq M_{2} \prod_{i=1}^{m}\left\|\frac{\partial X}{\partial u^{i}}(u)\right\|_{\mathbb{R}^{n}} \\
& =M_{2} \prod_{i=1}^{m}\left|\frac{\partial X}{\partial u^{i}}(u)\right|
\end{aligned}
$$

On the last expression we can apply the inequality of arithmetic and geometric means and the Hölder inequality and get

$$
M_{2} \prod_{i=1}^{m}\left|\frac{\partial X}{\partial u^{i}}(u)\right| \leq M_{2} \frac{1}{m^{m}}\left(\sum_{i=1}^{m}\left|\frac{\partial X}{\partial u^{i}}(u)\right|\right)^{m} \leq M_{2} \frac{1}{m^{m}} m^{m-1} \sum_{i=1}^{m}\left|\frac{\partial X}{\partial u^{i}}(u)\right|^{m} .
$$

Thus, it follows that

$$
\begin{aligned}
\mathcal{I}(X) & =\int_{u \in \Omega} I\left(X(u), \frac{\partial X}{\partial u^{1}}(u) \wedge \frac{\partial X}{\partial u^{2}}(u) \wedge \cdots \wedge \frac{\partial X}{\partial u^{m}}(u)\right) d u^{1} d u^{2} \ldots d u^{m} \\
& \leq M_{2} \frac{1}{m} \int_{u \in \Omega} \sum_{i=1}^{m}\left|\frac{\partial X}{\partial u^{i}}(u)\right|^{m} d u^{1} d u^{2} \ldots d u^{m} \\
& \leq M_{2} \frac{1}{m}\|X\|_{W^{1, m}\left(B, \mathbb{R}^{n}\right)}^{m}<\infty .
\end{aligned}
$$

Hildebrandt and von der Mosel proved the following result on the minimisation of the Cartan
functional in the class of admissible surface $\mathcal{C}(\Gamma)$ introduced in Section 3.1. Note that therein we chose $\mathcal{M}=B$ as the parameter domain of the competing surfaces $X$. Further, $m=2$ and $X \in W^{1,2}\left(B, \mathbb{R}^{n}\right)$ for $X \in \mathcal{C}(\Gamma)$ such that the value of the Cartan functional is finite by the preceding calculation.

Theorem 3.2.1 (Plateau problem for Cartan integrands, [HvdM03b, Theorems 1.4 and 1.5, p. 928]) Suppose $I: \mathbb{R}^{n} \times \mathbb{R}^{N} \rightarrow \mathbb{R}$ satisfies $(\mathrm{R})$, ( H ), (D) and (C). Then for any given rectifiable Jordan curve $\Gamma \subset \mathbb{R}^{n}$ there exists a surface $X \in \mathcal{C}(\Gamma)$, such that

$$
\mathcal{I}(X)=\inf _{\mathcal{C}(\Gamma)} \mathcal{I}(\cdot)
$$

In addition, the minimiser $X$ is Euclidean conformally parametrised almost everywhere on $B$, that is,

$$
\left|\frac{\partial X}{\partial u^{1}}\right|^{2}=\left|\frac{\partial X}{\partial u^{2}}\right|^{2} \quad \text { and } \quad\left\langle\frac{\partial X}{\partial u^{1}}, \frac{\partial X}{\partial u^{2}}\right\rangle_{\mathbb{R}^{2}}=0 \quad \mathcal{H}^{2}-\text { a.e. on } B .
$$

Furthermore, the minimiser $X$ has the following regularity property,

$$
X \in C^{0}\left(\bar{B}, \mathbb{R}^{n}\right) \cap C^{0, \sigma}\left(B, \mathbb{R}^{n}\right) \cap W^{1, q}\left(B, \mathbb{R}^{n}\right)
$$

for $\sigma:=\left(\frac{m_{F}}{M_{F}}\right)^{2} \in(0,1]$ and some $q>2$.
Note that in [HvdM03b] a Cartan integrand is called positive definite if the inequalities in (D) hold for all of $\mathbb{R}^{n} \times \mathbb{R}^{N}$. However, an analysis of their proof shows that condition (D) as stated above suffices to prove Theorem 3.2.1.

### 3.3 Representing the Busemann-Hausdorff area as a Cartan functional

To prove Theorem 3.1.1, we wish to identify the two-dimensional Busemann-Hausdorff area integrand $a_{2}^{F}$ for a reversible Finsler metric on $\mathcal{N}=\mathbb{R}^{n}$ as a positive definite, semi-elliptic Cartan integrand. Recall from (2.1.6) that $a_{m}^{F}$ is the function

$$
\begin{aligned}
a_{m}^{F}: & \bigsqcup_{q \in \mathcal{N}} G C_{m}\left(T_{q} \mathbb{R}^{n}\right) \rightarrow \mathbb{R}_{+} \\
& \left(q, w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right) \mapsto \frac{\varepsilon_{m}}{\mathcal{H}^{m}\left(\left\{v \in \mathbb{R}^{m} \mid F\left(q, v^{\alpha} w_{\alpha}\right) \leq 1\right\}\right)}
\end{aligned}
$$

where we sum over Greek indices $\alpha=1,2, \ldots, m$.
As discussed in the previous chapter, the tangent spaces $T_{q} \mathbb{R}^{n}$ are canonically isomorphic to $\mathbb{R}^{n}$ itself and we mentioned above that $\bigwedge^{m}\left(T_{q} \mathbb{R}^{n}\right) \cong \bigwedge^{m}\left(\mathbb{R}^{n}\right) \cong \mathbb{R}^{N}$. So we can identify the disjoint union as

$$
\begin{equation*}
\bigsqcup_{q \in \mathbb{R}^{n}} \bigwedge^{m}\left(T_{q} \mathbb{R}^{n}\right) \cong \bigsqcup_{q \in \mathbb{R}^{n}} \mathbb{R}^{N}=\bigcup_{q \in \mathbb{R}^{n}}\{q\} \times \mathbb{R}^{N}=\mathbb{R}^{n} \times \mathbb{R}^{N} \tag{3.3.1}
\end{equation*}
$$

One can show that $\bigsqcup_{q \in \mathcal{N}} \bigwedge^{m}\left(T_{q} \mathcal{N}\right)$ is a vector bundle of rank $\boldsymbol{N}$ (see [Lee13, Exercise 14.14, p. 359]). In fact, if $\mathcal{N}=\mathbb{R}^{n}$ then it is a trivial bundle which means that the identification (3.3.1) is a homeomorphism. Under this identification, the Busemann-Hausdorff area integrand $a_{m}^{F}$ is a function mapping the subset $\mathbb{R}^{n} \times \mathbb{G} \mathbb{C}_{m}\left(\mathbb{R}^{n}\right) \subset \mathbb{R}^{n} \times \mathbb{R}^{N}$ to $\mathbb{R}_{+}$。

Recall as well that by (2.1.9)

$$
a_{m}^{F}\left(q, v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right)=A_{q, m}^{b h}\left(v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m}\right)
$$

for $q \in \mathbb{R}^{n}$ and $v_{1} \wedge v_{2} \wedge \cdots \wedge v_{m} \in G C_{m}\left(\mathbb{R}^{n}\right)$. Therein, $A^{b h}=A_{q, m}^{b h}$ was the $m$-dimensional Busemann-Hausdorff area density for the normed space $\left(\mathbb{R}^{n}, F(q, \cdot)\right) \cong\left(T_{q} \mathbb{R}^{n}, F(q, \cdot)\right)$.

### 3.3.1 SEMI-ELLIPTICITY AND HOMOGENEITY OF $\mathfrak{a}_{2}^{F}$ ON $\mathbb{R}^{n} \times \mathbb{R}^{N}$

For $m=2$ (and thus $N=\binom{n}{2}$ ) we showed in Theorem 2.2.2 that the two-dimensional BusemannHausdorff area density $A_{q, 2}^{b h}: G C_{2}\left(\mathbb{R}^{n}\right) \cong \mathbb{G C}_{2}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}_{+}$is convex, that is, it admits an absolutely homogeneous, continuous and convex extension $\mathcal{A}=\mathcal{A}_{q, 2}: \bigwedge^{2}\left(\mathbb{R}^{n}\right) \cong \mathbb{R}^{N} \rightarrow \mathbb{R}_{+}($see Definition 2.2.1)

Therefore, the two-dimensional Busemann-Hausdorff area integrand extends to a function

$$
\mathfrak{a}_{2}^{F}: \mathbb{R}^{n} \times \mathbb{R}^{N} \rightarrow \mathbb{R}_{+},(q, \sigma) \mapsto \mathcal{A}_{q, 2}(\sigma)
$$

which we call the extended two-dimensional Busemann-Hausdorff area integrand.
Since $\mathcal{A}$ is convex for every normed space $\left(\mathbb{R}^{n}, F(q, \cdot)\right)$, the function $\mathfrak{a}_{2}^{F}$ is convex in its second variable and therefore the extended two-dimensional Busemann-Hausdorff area integrand is semielliptic on $\mathbb{R}^{n} \times \mathbb{R}^{N}$, that is, (C) holds. Further, $\mathfrak{a}_{2}^{F}$ is homogeneous of degree one in its second variable because $\mathcal{A}$ is and so, (H) holds.

### 3.3.2 Continuity of $\mathfrak{a}_{2}^{F}$ ON $\mathbb{R}^{n} \times \mathbb{R}^{N}$

For the regularity of $\mathfrak{a}_{2}^{F}$, we consider the extended Busemann-Hausdorff area integrand as a mapping from $\bigsqcup_{q \in \mathbb{R}^{n}} \bigwedge^{2}\left(T_{q} \mathbb{R}^{n}\right)$ to $\mathbb{R}$.

Theorem 3.3.1 (Continuity of $\mathfrak{a}_{2}^{\boldsymbol{F}}$ )
The extended two-dimensional Busemann-Hausdorff area integrand $\mathfrak{a}_{2}^{F}: \bigsqcup_{q \in \mathbb{R}^{n}} \bigwedge^{2}\left(T_{q} \mathbb{R}^{n}\right) \rightarrow \mathbb{R}_{+}$is a continuous function.Therefore, by using the identification (3.3.1)

$$
\mathfrak{a}_{2}^{F} \in C^{0}\left(\mathbb{R}^{n} \times \mathbb{R}^{N}\right)
$$

that is, the condition ( R ) holds true.
Proof: We recall that $\mathfrak{a}_{2}^{F}(q, \sigma)=\mathcal{A}_{q, 2}(\sigma)$ for $(q, \sigma) \in \bigsqcup_{q \in \mathbb{R}^{n}} \Lambda^{2}\left(T_{q} \mathbb{R}^{n}\right)$. Each of the extensions $\mathcal{A}_{q, 2}$ is a continuous function by Theorem 2.2.2 and Definition 2.2.1. Since $\mathbb{R}^{n}$ is a smooth manifold, we know that its tangent spaces $T_{q} \mathbb{R}^{n}$ vary continuously in $q$. Thus, $\mathfrak{a}_{2}^{F}$ is also continuous in its first variable. Using local trivialisations of the vector bundle $\bigsqcup_{q \in \mathbb{R}^{n}} \bigwedge^{2}\left(T_{q} \mathbb{R}^{n}\right)$, it then follows that $\mathfrak{a}_{2}^{F}$ is continuous on all of $\bigsqcup_{q \in \mathbb{R}^{n}} \bigwedge^{2}\left(T_{q} \mathbb{R}^{n}\right)$.

### 3.3.3 Positive definiteness of $\mathfrak{a}_{2}^{F}$ on $\mathbb{R}^{n} \times \mathbb{R}^{N}$

Condition (D) is only a condition on $\mathbb{R}^{n} \times \mathbb{G C}_{m}\left(\mathbb{R}^{n}\right)$. Therefore, we need only show it for the Busemann-Hausdorff area integrand and not for its extension. In fact, (D) holds true for $a_{m}^{F}$ for arbitrary $m$ as we will see in the following. In the work of Overath [Ove13] it was found that the Busemann-Hausdorff area integrand can be represented as a spherical integral. This result will provide us with a way to quantify the positive definiteness of $a_{m}^{F}$.

Theorem 3.3.2 (Spherical integral representation of $\boldsymbol{a}_{\boldsymbol{m}}^{\boldsymbol{F}}$, [Ove13, Theorem 2.1.6, p. 70])
Let $n \geq m>0,\left(\mathcal{N}^{n}, F\right)$ be a Finsler manifold where $F$ is strictly positive on $T \mathcal{N} \backslash o$. Let $q \in \mathcal{N}$ and
$w_{1} \wedge w_{2} \wedge \ldots \wedge w_{m} \in G C_{m}\left(T_{q} \mathcal{N}\right)$. Then we have

$$
\begin{equation*}
a_{m}^{F}\left(q, w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right)=\left(\frac{1}{\mathcal{H}^{m-1}\left(\mathbb{S}^{m-1}\right)} \int_{\mathbb{S}^{m-1}} \frac{1}{\left(F\left(q, \theta^{\alpha} w_{\alpha}\right)\right)^{m}} d S(\theta)\right)^{-1} \tag{3.3.2}
\end{equation*}
$$

Therein $d S(\theta)$ is the spherical measure and $\mathbb{S}^{m-1}=\left\{x \in \mathbb{R}^{m}| | x \mid=1\right\}$ is the $m$-dimensional Euclidean unit sphere.

Proof: We need only rewrite the Hausdorff measure in expression (2.1.6) by a change of variables (see [For09, Satz 8, p. 144]). Using polar coordinates we can write any $v=\left(v^{\alpha}\right)_{\alpha=1}^{m} \in \mathbb{R}^{m}$ as $v=s \theta$ where $\theta=v /|v| \in \mathbb{S}^{m-1}$ and $s=|v| \in[0, \infty)$. Then

$$
\begin{aligned}
& \mathcal{H}^{m}\left(\left\{v \in \mathbb{R}^{m} \mid F\left(q, v^{\alpha} w_{\alpha}\right) \leq 1\right\}\right) \\
&= \int_{\mathbb{R}^{m}} \chi_{\left\{v \in \mathbb{R}^{m} \mid F\left(q, v^{\alpha} w_{\alpha}\right) \leq 1\right\}}(x) d \mathcal{L}^{m}(x) \\
&=\left.\int_{\mathbb{S}^{m-1}} \int_{0}^{\infty} \chi_{\left\{\left(r \theta^{\prime}\right) \in \mathbb{R}^{m}\right.} \mid F\left(q, r \theta^{\prime \alpha} w_{\alpha}\right) \leq 1\right\} \\
&(s \theta) s^{m-1} d s d S(\theta)
\end{aligned}
$$

wherein $\chi_{M}$ is the characteristic function of $\boldsymbol{M}$, that is, $\chi_{M}(z):=1$ if $z \in M$ and 0 otherwise. The Finsler metric $F$ is positive homogeneous in its second component by definition. Further, $F$ is strictly positive on $T \mathcal{N} \backslash o$. Therefore,

$$
\begin{aligned}
&\left.\int_{\mathbb{S}^{m-1}} \int_{0}^{\infty} \chi_{\left\{\left(r \theta^{\prime}\right) \in \mathbb{R}^{m}\right.} \mid F\left(q, r \theta^{\prime \alpha} w_{\alpha}\right) \leq 1\right\} \\
&=\left.\int_{\mathbb{S}^{m-1}} \int_{0}^{\infty} \chi_{\left\{\left(r \theta^{\prime}\right) \in \mathbb{R}^{m}\right.} \mid r \leq\left(F\left(q, \theta^{\prime \alpha} w_{\alpha}\right)\right)^{-1}\right\} \\
&= \int_{\mathbb{S}^{m-1}} \int_{0}(s \theta) s^{m-1} d s d S(\theta) \\
&= \int_{\mathbb{S}^{m-1}} \frac{1}{\left.m\left(q, \theta^{\alpha} w_{\alpha}\right)\right)^{-1}} s^{m-1} d s d S(\theta) \\
&\left.m\left(q, \theta^{\alpha} w_{\alpha}\right)\right)^{m} d S(\theta) .
\end{aligned}
$$

Substituting into (2.1.6) and using $\mathcal{H}^{m-1}\left(\mathbb{S}^{m-1}\right)=m \mathcal{H}^{m}\left(B_{1}^{m}(0)\right)=m \varepsilon_{m}$, we find

$$
\begin{aligned}
a_{m}^{F}\left(q, w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right) & =\frac{\varepsilon_{m}}{\int_{\mathbb{S}^{m-1}} \frac{1}{m\left(F\left(q, \theta^{\alpha} w_{\alpha}\right)\right)^{m}} d S(\theta)} \\
& =\left(\frac{1}{\mathcal{H}^{m-1}\left(\mathbb{S}^{m-1}\right)} \int_{\mathbb{S}^{m-1}} \frac{1}{\left(F\left(q, \theta^{\alpha} w_{\alpha}\right)\right)^{m}} d S(\theta)\right)^{-1}
\end{aligned}
$$

The next result shows that the Busemann-Hausdorff area integrand coincides with the classical area integrand in the case that the Finsler metric is the Euclidean norm.

Lemma 3.3.3 (Euclidean area integrand)
Define $F: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow R_{+}$by $F(q, v):=|v|$ where $|\cdot|$ is the Euclidean norm on $\mathbb{R}^{n}$. Then $F$ is a Finsler metric on $\mathbb{R}^{n}$ which is strictly positive on $\mathbb{R}^{n} \times\left(\mathbb{R}^{n} \backslash\{0\}\right)$ and

$$
\begin{equation*}
a_{m}^{|\cdot|}(q, \sigma)=|\sigma| \quad \text { for all }(q, \sigma) \in \mathbb{R}^{n} \times \mathbb{G} \mathbb{C}_{m}\left(\mathbb{R}^{n}\right) \tag{3.3.3}
\end{equation*}
$$

where $|\cdot|$ denotes the $n$-dimensional Euclidean norm on the left-hand side and the $\binom{n}{m}$-dimensional one on the right-hand side.

Proof: Clearly, $F$ is non-negative on $\mathbb{R}^{n} \times \mathbb{R}^{n}$, strictly positive on $\mathbb{R}^{n} \times\left(\mathbb{R}^{n} \backslash\{0\}\right)$ and positively homogeneous in its second argument. In addition, $F \in C^{\infty}\left(\mathbb{R}^{n} \times\left(\mathbb{R}^{n} \backslash\{0\}\right)\right) \cap C^{0}\left(\mathbb{R}^{n} \times \mathbb{R}^{n}\right)$. Note that

$$
F(q, v)=\left(\delta_{i j} v^{i} v^{j}\right)^{\frac{1}{2}}
$$

where $\delta_{i j}$ is the Kronecker delta. Then the numbers

$$
g_{i j}(q, v)=\frac{1}{2} \frac{\partial^{2}}{\partial v^{i} \partial v^{j}} F^{2}(q, v)=\frac{1}{2} \delta_{i j}
$$

form a positive definite matrix, so that the ellipticity condition for $F$ is fulfilled.
To prove the second assertion, we follow the calculation in [Ove13, pp. 72-73]. The goal is to apply the area formula [EG92, Section 3.3.2, Theorem 1] to the denominator in (2.1.6). Fix $q \in \mathbb{R}^{n}$ and $\sigma \in G C_{m}\left(\mathbb{R}^{n}\right)$ and suppose $\sigma=w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}$ where $w_{1}, w_{2}, \cdots, w_{m} \in \mathbb{R}^{n}$ are linearly independent.

We will prove

$$
\begin{equation*}
\varepsilon_{m}=\|\sigma\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)} \mathcal{H}^{m}\left(\left\{v \in \mathbb{R}^{m}| | v^{\alpha} w_{\alpha} \mid \leq 1\right\}\right) \tag{3.3.4}
\end{equation*}
$$

wherewith it follows that

$$
a_{m}^{|\cdot|}(q, \sigma)=\|\sigma\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)}
$$

As mentioned in the beginning of the present section, we recall from Corollary 1.1.13 that the spaces $\left(\bigwedge^{m}\left(\mathbb{R}^{n}\right),\|\cdot\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)}\right)$ and $\left(\mathbb{R}^{N},|\cdot|\right)$ are isometrically isomorph. So

$$
\|\sigma\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)}=|\sigma|
$$

where we use the same symbol $\sigma$ for an element of $G C_{m}\left(\mathbb{R}^{n}\right) \subset \bigwedge^{m}\left(\mathbb{R}^{n}\right)$ and of $\mathbb{G} \mathbb{C}_{m}\left(\mathbb{R}^{n}\right) \subset \mathbb{R}^{N}$ on the left- and right-hand side respectively. Combining this with the above, the assertion holds true.

To prove (3.3.4) define the map $f: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}, v \mapsto v^{\alpha} w_{\alpha}$. This is an injective linear mapping with $\operatorname{im}(f)=\operatorname{span}\left\{w_{1}, w_{2}, \ldots, w_{m}\right\} \subset \mathbb{R}^{n}$. Note that the Jacobian matrix of $f$ is

$$
\begin{aligned}
D f(v) & :=\left(\begin{array}{ccc}
\frac{\partial f^{1}}{\partial v^{1}} & \cdots & \frac{\partial f^{1}}{\partial v^{m}} \\
\vdots & \ddots & \vdots \\
\frac{\partial f^{n}}{\partial v^{1}} & \cdots & \frac{\partial f^{n}}{\partial v^{m}}
\end{array}\right) \\
& =\left(w_{1}\left|w_{2}\right| \cdots \mid w_{m}\right) \in \mathbb{R}^{n \times m}
\end{aligned}
$$

and the Jacobian determinant of $f$ is

$$
\begin{aligned}
J f(v) & :=\sqrt{\operatorname{det}\left(D f(v)^{T} D f(v)\right)} \\
& =\sqrt{\operatorname{det}\left(\left\langle w_{\alpha}, w_{\beta}\right\rangle_{\mathbb{R}^{n}}\right)_{\alpha, \beta=1, \ldots, m}} \\
& =\left\|w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)}=\|\sigma\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)}
\end{aligned}
$$

Define the two sets

$$
\begin{aligned}
& A:=\left\{v \in \mathbb{R}^{m}| | v^{\alpha} w_{\alpha} \mid \leq 1\right\} \subset \mathbb{R}^{m} \\
& \Omega:=\overline{B_{1}^{n}(0)} \cap \operatorname{span}\left\{w_{1}, w_{2}, \ldots, w_{m}\right\} \subset \mathbb{R}^{n}
\end{aligned}
$$

wherein $B_{1}^{n}(0) \subset \mathbb{R}^{n}$ denotes the $n$-dimensional Euclidean unit ball. As a first step we will show $f(A)=\Omega$. Let $y \in f(A)$, then $y=v^{\alpha} w_{\alpha} \in \operatorname{span}\left\{w_{1}, w_{2}, \ldots, w_{m}\right\}$ and $\left|v^{\alpha} w_{\alpha}\right|=|y| \leq 1$. Therefore, $y \in \Omega$. Conversely, suppose $y \in \Omega$. Then $y=v^{\alpha} w_{\alpha}$ for some $v \in \mathbb{R}^{m}$ and $|y|=\left|v^{\alpha} w_{\alpha}\right| \leq 1$. Hence, $y \in f(A)$. In a next step, we show

$$
\mathcal{H}^{0}\left(A \cap f^{-1}(\{\cdot\})\right)=\chi_{\Omega}(\cdot)
$$

If $y \notin \operatorname{im}(f)=\operatorname{span}\left\{w_{1}, w_{2}, \ldots, w_{m}\right\}$ then $f^{-1}(\{y\})=\varnothing$ and $A \cap f^{-1}(\{y\})=\varnothing$. Therefore, $\mathcal{H}^{0}\left(A \cap f^{-1}(\{y\})\right)=0$ and $\chi_{\Omega}(y)=0$ because $\Omega \subset \operatorname{im}(f)$.

If $y \in \operatorname{im}(f) \cap f(A)=\operatorname{span}\left\{w_{1}, w_{2}, \ldots, w_{m}\right\} \cap \Omega$ then the injectivity of $f$ implies that the cardinality of $f^{-1}(\{y\})=A \cap f^{-1}(\{y\})$ is 1 . Therefore, $\mathcal{H}^{0}\left(A \cap f^{-1}(\{y\})\right)=1=\chi_{\Omega}(y)$. If $y \notin \Omega$ but still $y \in \operatorname{span}\left\{w_{1}, w_{2}, \ldots, w_{m}\right\}$ then $A \cap f^{-1}(\{y\})=\varnothing$ and consequently, $\mathcal{H}^{0}\left(A \cap f^{-1}(\{y\})\right)=0=\chi_{\Omega}(y)$.

Using the area formula, we get

$$
\begin{aligned}
\|\sigma\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)} & \mathcal{H}^{m}\left(\left\{v \in \mathbb{R}^{n} \mid F\left(q, v^{\alpha} w_{\alpha}\right) \leq 1\right\}\right) \\
& =\|\sigma\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)} \mathcal{H}^{m}(A)=\int_{v \in A}\|\sigma\|_{\Lambda^{m}\left(\mathbb{R}^{n}\right)} d \mathcal{L}^{m}(v) \\
& =\int_{v \in A} J f(v) d \mathcal{L}^{m}(v)=\int_{y \in \mathbb{R}^{n}} \mathcal{H}^{0}\left(A \cap f^{-1}(\{y\})\right) d \mathcal{H}^{m}(y) \\
& =\int_{y \in \mathbb{R}^{n}} \chi_{\Omega}(y) d \mathcal{H}^{m}(y)=\mathcal{H}^{m}\left(\overline{B_{1}^{n}(0)} \cap \operatorname{span}\left\{w_{1}, w_{2}, \ldots, w_{m}\right\}\right) \\
& =\mathcal{H}^{m}\left(\overline{B_{1}^{m}(0)}\right)=\varepsilon_{m}
\end{aligned}
$$

In the last equation we used that the $n$-dimensional Euclidean unit ball is centrally symmetric. This proves (3.3.4).

Now we can compare the value of the Busemann-Hausdorff area integrand for two different Finsler metrics.

Lemma 3.3.4 ([OvdM14, Lemma 2.4, p. 286])
Suppose $F_{1}, F_{2}$ are two Finsler metrics on $\mathbb{R}^{n}$ which are both strictly positive on $\mathbb{R}^{n} \times\left(\mathbb{R}^{n} \backslash\{0\}\right)$. If for $q \in \mathbb{R}^{n}$ there exist numbers $0<c_{1}(q) \leq c_{2}(q)$ with

$$
c_{1}(q) F_{1}(q, v) \leq F_{2}(q, v) \leq c_{2}(q) F_{1}(q, v) \quad \text { for all } v \in \mathbb{R}^{n}
$$

then

$$
\begin{equation*}
m_{1}(q) a_{m}^{F_{1}}(q, \sigma) \leq a_{m}^{F_{2}}(q, \sigma) \leq c_{2}(q) a_{m}^{F_{1}}(q, \sigma) \quad \text { for all } \sigma \in \mathbb{G C}_{m}\left(\mathbb{R}^{n}\right) \tag{3.3.5}
\end{equation*}
$$

where $m_{i}(q):=c_{i}^{m}(q)$ for $i=1,2$.
Proof: We use Theorem 3.3.2 to get for $\sigma=w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m} \in G C_{m}\left(\mathbb{R}^{n}\right)$

$$
\begin{aligned}
m_{1}(q) a_{m}^{F_{1}}\left(q, w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right) & =c_{1}^{m}(q)\left(\frac{1}{\mathcal{H}^{m-1}\left(\mathbb{S}^{m-1}\right)} \int_{\mathbb{S}^{m-1}} \frac{1}{\left(F_{1}\left(q, \theta^{\alpha} w_{\alpha}\right)\right)^{m}} d S(\theta)\right)^{-1} \\
& =\left(\frac{1}{\mathcal{H}^{m-1}\left(\mathbb{S}^{m-1}\right)} \int_{\mathbb{S}^{m-1}} \frac{1}{\left(c_{1}(q) F_{1}\left(q, \theta^{\alpha} w_{\alpha}\right)\right)^{m}} d S(\theta)\right)^{-1} \\
& \leq\left(\frac{1}{\mathcal{H}^{m-1}\left(\mathbb{S}^{m-1}\right)} \int_{\mathbb{S}^{m-1}} \frac{1}{\left(F_{2}\left(q, \theta^{\alpha} w_{\alpha}\right)\right)^{m}} d S(\theta)\right)^{-1} \\
& =a_{m}^{F_{2}}\left(q, w_{1} \wedge w_{2} \wedge \cdots \wedge w_{m}\right)
\end{aligned}
$$

The second inequality follows similarly.

Finally, we combine the preceding results to find the positive definiteness of $a_{m}^{F}$.

Theorem 3.3.5 (Positive definiteness of $\boldsymbol{a}_{\boldsymbol{m}}^{\boldsymbol{F}}$, [OvdM14, Corollary 2.5, p. 287])
Let $F$ be a Finsler metric on $\mathbb{R}^{n}$ with

$$
\begin{equation*}
0<c_{1}:=\inf _{\mathbb{R}^{n} \times \mathbb{S}^{n-1}} F(\cdot, \cdot) \leq \sup _{\mathbb{R}^{n} \times \mathbb{S}^{n-1}} F(\cdot, \cdot)=\|F\|_{L^{\infty}\left(\mathbb{R}^{n} \times \mathbb{S}^{n-1}\right)}=: c_{2}<\infty \tag{3.3.6}
\end{equation*}
$$

Then

$$
\begin{equation*}
m_{1}|\sigma| \leq a_{m}^{F}(q, \sigma) \leq m_{2}|\sigma| \quad \text { for all }(q, \sigma) \in \mathbb{R}^{n} \times \mathbb{G} \mathbb{C}_{m}\left(\mathbb{R}^{n}\right) \tag{3.3.7}
\end{equation*}
$$

where $m_{i}:=c_{i}^{m}$ for $i=1,2$, that is, (D) holds true.
Proof: Through the positive homogeneity of $F$ in its second argument, we see from (3.3.6) that

$$
c_{1}|v| \leq|v| F\left(q, \frac{v}{|v|}\right)=F(q, v) \leq c_{2}|v|
$$

for all $(q, v) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$. For $q \in \mathbb{R}^{n}$ set $c_{i}(q):=c_{i}$ for $i=1,2$. The functions $F_{1}(q, v):=|v|$ and $F_{2}(q, v):=F(q, v)$ both are Finsler metrics on $\mathbb{R}^{n}$ which are strictly positive on $\mathbb{R}^{n} \times\left(\mathbb{R}^{n} \backslash\{0\}\right)$. Due to Lemma 3.3.3 we find in addition that

$$
a_{m}^{|\cdot|}(q, \sigma)=|\sigma| \quad \text { for all }(q, \sigma) \in \mathbb{R}^{n} \times \mathbb{G}_{m}\left(\mathbb{R}^{n}\right)
$$

where $|\cdot|$ on the left and right-hand side means the Euclidean norm on $\mathbb{R}^{n}$ and $\mathbb{R}^{N}$ respectively. Therefore, we can apply Lemma 3.3.4 to $F_{1}$ and $F_{2}$ and obtain (3.3.7) from (3.3.5).

We conclude this chapter and the thesis with the proof of the Plateau problem in reversible Finsler space.

Proof of Theorem 3.1.1: In Sections 3.3.1, 3.3.2 and 3.3.3 we saw that the extended two-dimensional Busemann-Hausdorff area integrand $\mathfrak{a}_{2}^{F}: \mathbb{R}^{n} \times \mathbb{R}^{N} \rightarrow \mathbb{R}$ is a positive definite, semi-elliptic Cartan integrand. Recall that the Busemann-Hausdorff area functional area ${ }_{B}^{F}(\cdot)$ in (2.1.7) is given by

$$
\operatorname{area}_{\Omega}^{F}(X)=\int_{p \in \Omega} a_{m}^{F}\left(X(p), d X_{p}\left(\left.\frac{\partial}{\partial u^{1}}\right|_{p}\right) \wedge \cdots \wedge d X_{p}\left(\left.\frac{\partial}{\partial u^{m}}\right|_{p}\right)\right) d u^{1} \wedge \cdots \wedge d u^{m}
$$

Thus, $\operatorname{area}_{B}^{F}(\cdot)$ is the Cartan functional corresponding to $I=\mathfrak{a}_{2}^{F}$. Therefore, we can apply The-
orem 3.2.1 to the present situation which yields the stated result.

## Conclusion and prospects

In this thesis we solved the Plateau problem in Finsler space in arbitrary codimension for a reversible Finsler metric. We observed in Chapter 3 that the proof technique used in [OvdM14] for the codimension one case can be applied to solve the area minimisation problem in higher codimension. Essentially, only the convexity of the two-dimensional Busemann-Hausdorff area density, established in [BI12], was needed. It was especially important that the underlying Finsler metric is reversible, since otherwise all the arguments involving symmetric polygons in Chapter 2 do not hold.

Overath and von der Mosel, however, did not restrict themselves to reversible Finsler metrics. They considered the so-called $m$-harmonic symmetrisation $F_{\text {sym }}(x, y)=2^{\frac{1}{m}}\left(F(x, y)^{-m}+F(x,-y)^{-m}\right)^{-\frac{1}{m}}$ of a Finsler metric $F$. Note that $F_{\text {sym }}$ is always reversible and coincides with $F$, if the original Finsler metric is reversible. In [OvdM14] as a general assumption only such Finsler metrics are considered whose $m$-harmonic symmetrisation is also a Finsler metric. Their idea is to apply the convexity result to the reversible $m$-harmonic symmetrisation and use a result comparing the corresponding Finsler areas. In fact, Overath and von der Mosel showed that the Busemann-Hausdorff area functionals corresponding to $F$ and $F_{\text {sym }}$ coincide (see [OvdM14, Theorem $1.1+$ Lemma 2.3, p. 276+285]) - but their proof is restricted to codimension one.

Further analysis is needed to see if a similar result holds true in higher codimension to fully generalise Overath and von der Mosel's work. It remains to be investigated if their results on higher regularity of solutions (see [OvdM14, Theorem 1.4, p. 281]) generalise straightforwardly to higher codimension.

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Convex extension of the Busemann-Hausdorff area integrand
and the Plateau problem in arbitrary co-dimension


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