A generalized polar-coordinate integration formula, oscillatory integral techniques, and applications to convolution powers of complex-valued functions on \mathbb{Z}^d

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Abstract

Acknowledgments

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Chapter 1

Introduction

Chapter 2

Positive Homogeneous Functions

In this chapter, we make precise and extend the definition of positive homogeneous polynomials as appeared in [13]. We also introduce the notions of **subhomogeneous functions** and **strongly subhomogeneous functions**, which appear alongside positive homogeneous polynomials/functions in the Taylor expansion of Γ I'll explain what Γ is, "formally," in the introduction.. These classes of functions are key ingredients for estimating the convolution powers of ϕ via the Van der Corput lemma, particularly in the cases where oscillatory integrals are involved.

To start, we shall review the notion of continuous one-parameter groups defined in Definition 2.0.1. Continuous one-parameter groups appear naturally when one characterizes the radial scaling in a general quasi-polar parameterization. For instance, in the usual polar coordinates we have **isotropic** scaling in the radial variable r and can choose our continuous one-parameter group to be $\{T_r\}_{r>0} = r^E$ where E is a linear transformation with standard matrix representation diag(1/2,1/2). It is clear that $\{T_r\}_{r>0}$ is a group, and that any point $v \in \mathbb{R}^d \setminus \{0\}$ can be written uniquely as $r^E \eta$ where $\eta \in \mathbb{S}$. In general, E needs not be a multiple of the identity transformation \mathbb{I} (whose standard representation is of course diag(1,1)). In such cases, we have **anisotropic** scaling. For example, let E' be diag(1/2,1/4) in the standard basis. Any point $v \in \mathbb{R}^2 \setminus \{0\}$ can also be written uniquely as $r^{E'} \eta'$ where $\eta' \in \{(\eta'_1, \eta'_2) \in \mathbb{R}^2 : \eta'_1^2 + \eta'_2^4 = 1\}$ and r > 0.

Definition 2.0.1. $\{T_r\}_{r>0}\subseteq\operatorname{Gl}(\mathbb{R}^d)$ is said to be a continuous one-parameter group if

- 1. $T_r: \mathbb{R}^d \to \mathbb{R}^d$ is a linear operator for each r > 0;
- 2. T_r is continuous in r (in the usual topology in $Gl \mathbb{R}^d$);
- 3. $T_0 = \mathbb{I}_d$, where \mathbb{I}_d denotes the identity on $Gl(\mathbb{R}^d)$;
- 4. $T_{uv} = T_u T_v$ for all u, v > 0.

It is well-known (c.f., [4, 5, 13]) that every continuous one-parameter group $\{T_r\}$ has the unique representation

$$T_r = r^E = \exp((\ln r)E) = \sum_{k=0}^{\infty} \frac{(\ln r)^k}{k!} E^k$$

for some $E \in \operatorname{End}(\mathbb{R}^d)$. E is called (infinitesimal) generator of $\{T_r\}$, and $\{T_r\}$ is said to be generated by E. For simplicity, we shall use this representation throughout the paper, i.e., any element of $\{T_r\}_{r>0}$ has the form r^E . We write $T_r = r^E$ to be explicit about E.

An important aspect of the continuous one-parameter groups $\{T_r\}_{r>0}$ is the **contracting** property. Intuitively, this means that the radial scaling becomes vanishingly small as the parameter r approaches 0. This property will become useful in our construction of the measure on S where it is used to contract sufficiently nice sets $F \subseteq S$ into the unit ball.

Definition 2.0.2. A continuous one-parameter group $\{T_r\}$ is said to be contracting if

$$\lim_{r\to 0}||T_r||=0.$$

The next two results follow immediately from this definition. The first result is by virtue of continuity and the Banach-Steinhaus theorem (or, the uniform boundedness principle), while the second by the continuity of the determinant and the operator-norm topology on $\operatorname{End} \mathbb{R}^d$. The proof of Proposition 2.0.3 is presented in Appendix 6.1.

Proposition 2.0.3. Let $\{T_r\}$ be a continuous one-parameter group. Then $\{T_r\}$ is contracting if and only if for all $x \in \mathbb{R}^d$,

$$\lim_{t \to 0} |T_r x| = 0. {(2.1)}$$

Proposition 2.0.4. Let $\{T_r\} \subseteq Gl(\mathbb{R}^d)$ be a continuous one-parameter group with generator E. If $\{T_r\}$ is contracting, then $\operatorname{tr} E > 0$.

Proof. The supposition that $\{T_r\}$ is a contracting group implies that $r^E \to \mathbf{0}$ as $r \to 0$ in the operator-norm topology on $\operatorname{End}(\mathbb{R}^d)$; here $\mathbf{0}$ is zero transformation. Because the determinant is a continuous function from $\operatorname{End}(\mathbb{R}^d)$, equipped with operator-norm topology, into \mathbb{R} , we have

$$0 = \det(\mathbf{0}) = \det\left(\lim_{r \to 0} r^E\right) = \lim_{r \to 0} \det\left(r^E\right) = \lim_{r \to 0} r^{\operatorname{tr} E}$$

in view of the preceding proposition. Therefore, $\operatorname{tr} E > 0$.

Now we are ready to define the notion of "positive homogeneous functions." Recall the previous example of anisotropic scaling in which any point $v \in \mathbb{R}^2$ can be written as $r^E \eta$ where E has standard matrix representation $\operatorname{diag}(1/2,1/4)$ and $\eta \in \{(\eta_1,\eta_2) \in \mathbb{R}^2 : P(\eta_1,\eta_2) = \eta_1^2 + \eta_2^4 = 1\}$. We observe that the function P is "homogeneous" with respect to the matrix E in the sense that

$$P(r^E x) = (r^{1/2}x_1)^2 + (r^{1/4}x_2)^4 = rP(x)$$

for all $x \in \mathbb{R}^d$ and r > 0. Because such functions (which are not necessarily polynomials, as we will see) often show up in estimating convolution powers and their attractors, we will generalize and make precise this notion of homogeneity. In what follows, for a function $P : \mathbb{R}^d \to \mathbb{R}$, we shall call

$$S = \{ \eta \in \mathbb{R}^d : P(\eta) = 1 \}$$

the unital level set of P. We say that P is **positive definite** if P is non-negative and P(x) = 0 only when x = 0. Given a continuous one-parameter group $\{T_r\}$, we say that P is **homogeneous with respect to** $\{T_r\}$ if

$$rP(x) = P(T_r x) = P(r^E x)$$
(2.2)

for all r > 0 and $x \in \mathbb{R}^d$. For a given map $P : \mathbb{R}^d \to \mathbb{R}$, the set of all E for which (2.2) is satisfied is called the **exponent set of** P and shall be denoted by $\operatorname{Exp}(P)$. We note that even though P often appears as a polynomial in the convolution powers setting, we only require that it be a continuous function.

The forthcoming relates the unital level set S to the contracting group $\{r^E\}$ for a positive-definite P and lists the (equivalent) defining characteristics of a **positive-homogeneous function**. This characterization highlights an interplay between positive definite functions, contracting groups, and unital level sets and sets the stage for our definition of positive homogeneous functions

Just before we state and prove this proposition, however, it is helpful to fix some notation and introduce some basic topological objects connected to positive homogeneous functions. Given a positive homogeneous function P, we define

$$B_r = \{ \xi \in \mathbb{R}^d : P(\xi) < r \}$$
 and $A_s^r = \{ \xi \in \mathbb{R}^d : s \le P(\xi) < r \}$

for r>0 and $0 \le s < r$; these are P-adapted analogues of the Euclidean ball, \mathbb{B}_r , and the annulus of inner radius s and outer radius r, respectively. In view of Proposition 2.0.5 and the continuity of P, we see that, for each r>0, B_r is open and $\overline{B_r}$ is compact. Further, by setting $B=B_1$, it is a straightforward exercise to see that $\overline{B}=B\cup S$ where $\partial B=S$ is the unital level set associated to P.

Proposition 2.0.5. Let $P: \mathbb{R}^d \to \mathbb{R}$ be continuous, positive definite, and have $\text{Exp}(P) \neq \emptyset$. The following are equivalent:

- (a) S is compact.
- (b) There is a positive number M for which

for all $|x| \geq M$.

- (c) For each $E \in \text{Exp}(P)$, $T_r = r^E$ is contracting.
- (d) There exists $E \in \text{Exp}(P)$ for which $T_r = r^E$ is contracting.
- (e) We have

$$\lim_{x \to \infty} P(x) = \infty.$$

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Proof. In the case that d = 1, it is easy to see that every function satisfying the hypotheses is of the form

$$P(x) = \begin{cases} P(1)x^{\alpha} & \text{for } x \ge 0\\ P(-1)(-x)^{\alpha} & \text{for } x < 0 \end{cases}$$

for some $\alpha > 0$ where P(1), P(-1) > 0 and $\operatorname{Exp}(P)$ consists only of the linear function $x \mapsto x/\alpha$. In this setting, it is easy to see that Conditions (a)–(e) are satisfied (always and) simultaneously. We shall therefore assume that d > 1 for the remainder of the proof.

 $(a)\Rightarrow (b)$. Given that S is compact, it is bounded and so we have a positive number M for which $P(x)\neq 1$ for all $|x|\geq M$. Observe that, if for two points $x_1,x_2\in\mathbb{R}^d\setminus\mathbb{B}_M$, $P(x_1)<1< P(x_2)$ or $P(x_2)<1< P(x_1)$, then by virtue of the path connectedness of $\mathbb{R}^d\setminus\mathbb{B}_M$ and the intermediate value theorem, we would be able to find $x_0\in\mathbb{R}^d\setminus\mathbb{B}_M$ for which $P(x_0)=1$, an impossibility. Therefore, to show that Condition (b) holds, we must simply rule out the case in which P(x)<1 for all $|x|\geq M$. Let us therefore assume, to reach a contradiction, that this alternate condition holds. In this case, we take $E\in \operatorname{Exp}(P)$ and $y\in\mathbb{R}^d\setminus\{0\}$ and observe that

$$\lim_{r \to \infty} P(r^E y) = \lim_{r \to \infty} r P(y) = \infty.$$

By virtue of our supposition, we find that $|r^Ey| < M$ for all sufficiently large r. In particular, there exists a sequence $r_k \to \infty$ for which $|r_k^Ey| \le M$ for all k and

$$\lim_{k \to \infty} P(r_k^E y) = \infty.$$

Because $\overline{\mathbb{B}_M}$, the closure of \mathbb{B}_M , is compact, $\{r_k^Ey\}$ has a convergent subsequence which we also denote by $\{r_k^Ey\}$ by a slight abuse of notation. In view of the continuity of P at $\eta=\lim_{k\to\infty}r_k^Ey$, we have

$$P(\eta) = \lim_{k \to \infty} P(r_k^E y) = \lim_{k \to \infty} r_k P(y) = \infty,$$

which is impossible. Thus Condition (b) holds.

 $(b)\Rightarrow(c)$. We shall prove the contrapositive statement. Suppose that, for $E\in \operatorname{Exp}(P)$, $\{r^E\}$ is not contracting. In this case, by virtue of Proposition 2.0.3, there exists $x\in\mathbb{R}^d\setminus\{0\}$ and a sequence $r_k\to 0$ for which r_k^Ex does not converge to zero in \mathbb{R}^d . If our sequence $\{r_k^Ex\}$ is bounded, then it must have a convergent subsequence $\{r_{k_m}^Ex\}$ with non-zero subsequential limit $\eta=\lim_{m\to\infty}r_{k_m}^Ex$. By the continuity of P, we have

$$P(\eta) = \lim_{m \to \infty} P(r_{k_m}^E x) = \lim_{k \to \infty} r_{k_m} P(x) = 0$$

which cannot be true for it would violate the positive definiteness of P. We must therefore consider the other possibility: The sequence $\{r_k^Ex\}$ is unbounded. In particular, there must be some k_0 for which $r_{k_0} < 1/P(x)$ and $|r_{k_0}^Ex| > M$. Upon putting $y = r_{k_0}^Ex$, we have |y| > M and $P(y) = P(r_{k_0}^Ex) = r_{k_0}P(x) < 1$ which shows that Condition (b) cannot hold.

$$(c) \Rightarrow (d)$$
. This is immediate.

 $(d) \Rightarrow (e)$. Let $E \in \operatorname{Exp}(P)$ be such that the one-parameter group $\{r^E\}$ is contracting and let $\{x_k\} \subseteq \mathbb{R}^d$ be such that $\lim_{k \to \infty} |x_k| = \infty$. By virtue of Proposition 6.1.7, there exist sequences

 $\{r_k\}\subseteq (0,\infty)$ and $\{\eta_k\}\in \mathbb{S}$ for which $r_k^E\eta_k=x_k$ for all k and $\lim_{k\to\infty}r_k=\infty$. Given that P is continuous and strictly positive on the compact set \mathbb{S} , we have $\inf_{\eta\in\mathbb{S}}P(\eta)>0$ and therefore

$$\liminf_{k} P(x_k) = \liminf_{k} r_k P(\eta_k) \ge \liminf_{k} r_k \left(\inf_{\eta \in \mathbb{S}} P(\eta) \right) = \infty$$

showing that $\lim_{k\to\infty} P(x_k) = \infty$, as desired.

 $(e)\Rightarrow (a).$ Because S is the preimage of the closed singleton $\{1\}$ under the continuous function P, it is closed. By virtue of Condition (e), S is also be bounded and thus compact in view of the Heine-Borel theorem.

The characterizing Proposition 2.0.5 gives us the following definition.

Definition 2.0.6. Let $P : \mathbb{R}^d \to \mathbb{R}$ be continuous, positive definite and have $\operatorname{Exp}(P) \neq \emptyset$. If any one (and hence all) of the equivalent conditions in Proposition 2.0.5 are fulfilled, we say that P is positive homogeneous.

Example 1. For any $\alpha > 0$, the α th-power of the Euclidean norm $x \mapsto |x|^{\alpha}$ is positive homogeneous. In this case, the unital level set S is the standard unit sphere $\mathbb{S} = \mathbb{S}_1$, $B = \mathbb{B} = \mathbb{B}_1$, and

$$\operatorname{Exp}(|\cdot|^{\alpha}) = \frac{1}{\alpha}I + \mathfrak{o}(d)$$

where I is the identity and $\mathfrak{o}(d)$ is the Lie algebra of the orthogonal group $O_d(\mathbb{R})$ and is characterized by the set of skew-symmetric matrices.

Example 2. In the language of L. Hörmander [8], consider semi-elliptic polynomial of the form

$$P(x) = \sum_{|\alpha:\mathbf{n}|=1} a_{\alpha} x^{\alpha}, \tag{2.3}$$

where $\mathbf{n}=(n_1,n_2,\ldots,n_d)$ is a d-tuple of positive even natural numbers¹, and, for each multi-index $\alpha=(\alpha_1,\alpha_2,\ldots,\alpha_d)\in\mathbb{N}^d$,

$$|\alpha:\mathbf{n}|:=\sum_{k=1}^d \frac{\alpha_k}{n_k},$$

and

$$x^{\alpha} = (x^1)^{\alpha_1} (x^2)^{\alpha_2} \cdots (x^d)^{\alpha_d}$$

for $x = (x^1, x^2, \dots, x^d) \in \mathbb{R}^d$. If we consider $E \in \text{End}(\mathbb{R}^d)$ whose standard matrix representation is $\text{diag}(1/n_1, 1/n_2, \dots, 1/n_d)$, we have

$$P(r^{E}x) = \sum_{|\alpha:\mathbf{n}|=1} a_{\alpha} (r^{1/n_{1}}x^{1})^{\alpha_{1}} (r^{1/n_{2}}x^{2})^{\alpha_{2}} \cdots (r^{1/n_{d}}x^{d})^{\alpha_{d}} = \sum_{|\alpha:\mathbf{n}|=1} a_{\alpha}r^{|\alpha:\mathbf{n}|}x^{\alpha} = rP(x)$$

for all $x \in \mathbb{R}^d$ and r > 0 and therefore $E \in \text{Exp}(P)$. It is easy to see that $T_r = r^E$ is a contracting group and so we have the following statement by virtue of Proposition 2.0.5:

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¹In Section ??, we will write this as $\mathbf{n} = 2\mathbf{m}$ for $\mathbf{m} \in \mathbb{N}^d_+$.

If a semi-elliptic polynomial P(x) of the form (2.3) is positive definite, then it is positive homogeneous.

For two concrete examples, consider the polynomials P_1 and P_2 on \mathbb{R}^2 defined by

$$P_1(x,y) = x^2 + y^4$$
 and $P_2(x,y) = x^2 + \frac{3}{2}xy^2 + y^4$

defined for $(x,y) \in \mathbb{R}^2$. It is straightforward to see that P_1 and P_2 are both positive definite and semi-elliptic of the form (2.3) with $\mathbf{n}=(2,4)$. Figure 2.1 illustrates P_1 and P_2 along with their associated unital level sets S_1 and S_2 and corresponding sets $B_1=\{(x,y)\in\mathbb{R}^2:P_1(x,y)<1\}$ and $B_2=\{(x,y)\in\mathbb{R}^2:P_2(x,y)<1\}$ written with a slight abuse of notation.

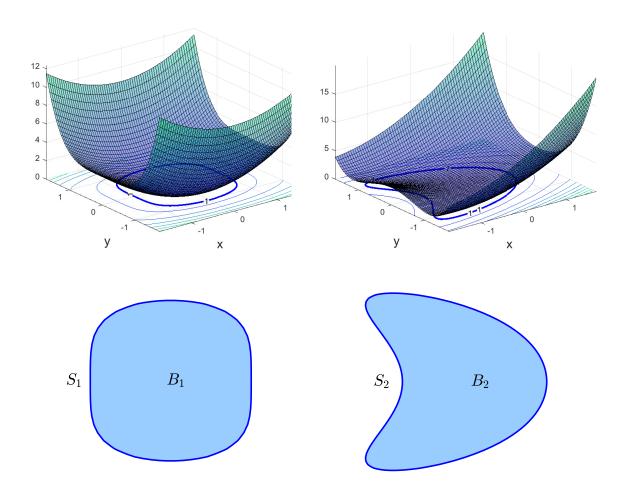


Figure 2.1: The left column illustrates the graph of P_1 with its associated S_1 and convex B_1 . The right column illustrates the graph of P_2 with its associated S_2 and non-convex B_2 .

Remark 1. In [13], a positive homogeneous polynomial P is, by definition, a complex-valued multivariate polynomial on \mathbb{R}^d for which $\operatorname{Exp}(P)$ contains an element of $\operatorname{End}(\mathbb{R}^d)$ whose spectrum is purely real and for which $R = \operatorname{Re} P$ is positive definite (see Proposition 2.0.7 below). By virtue of Proposition 2.2 of [13], for each such polynomial P and $E \in \operatorname{Exp}(P)$ with real spectrum, there exists $A \in \operatorname{Gl}(\mathbb{R}^d)$ (representing a change of basis of \mathbb{R}^d) and a d-tuple of even positive natural numbers $\mathbf{n} = (n_1, n_2, \ldots, n_d) \in \mathbb{N}^d_+$ for

which $A^{-1}EA$ has standard matrix representation $\operatorname{diag}(1/n_1,1/n_2,\ldots,1/n_d)$ and $(P\circ A)(x)$ is semi-elliptic of the form (2.3) with, in this case, complex coefficients. It follows that every real-valued positive homogeneous polynomial (in the sense of [13]) is a positive homogeneous function in the sense of the present article. Of course, the semi-elliptic polynomials discussed above are positive homogeneous polynomials in the sense of [13] where A=I. We refer the reader to Section 7.3 of [13] which presents a real-valued positive homogeneous polynomial which is not semi-elliptic (and so $A \neq I$).

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Example 3. Let Q be a positive homogeneous function with exponent set Exp(Q) and compact level set $S_Q = \{\eta : Q(\eta) = 1\}$. Given any $f \in C^0(S_Q)$ for which $f(\eta) > 0$ for all $\eta \in S_Q$ and $E \in \text{Exp}(Q)$, define $P = P_{f,E,Q} : \mathbb{R}^d \to \mathbb{R}$ by

$$P(x) = \begin{cases} Q(x)f\left((Q(x))^{-E}x\right) & x \neq 0\\ 0 & x = 0 \end{cases}$$

for $x \in \mathbb{R}^d$. We claim that P is positive homogeneous and $E \in \text{Exp}(P)$.

Proof. To see this, we first observe that, for any $x \in \mathbb{R}^d \setminus \{0\}$, $Q((Q(x))^{-E}x) = Q(x)/Q(x) = 1$ and hence $Q(x)^{-E}x \in S_Q$ and so the above formula makes sense and ensures that P is continuous on $\mathbb{R}^d \setminus \{0\}$. Furthermore, because f is continuous and positive on the compact set S_Q , we have $0 < \min f \le \max f < \infty$. From this it follows that P is positive definite and, by virtue of the squeeze theorem, continuous at x = 0. For any x > 0 and $x \in \mathbb{R}^d$, we have

$$P(r^Ex) = Q(r^Ex)f(Q(r^Ex)^{-E}r^Ex) = rQ(x)f(Q(x)^{-E}x) = rP(x)$$

and therefore $E \in \text{Exp}(P)$. Upon noting that $\{r^E\}$ is contracting by virtue of Proposition 2.0.5, we conclude that P is positive homogeneous.

The utility of this construction allows us to see that "most" positive homogeneous functions are not smooth. To see this, we fix a positive homogeneous function $Q \in C^{\infty}(\mathbb{R}^d)$ and remark that S_Q is necessarily a compact smooth embedded hypersurface of \mathbb{R}^d (see Proposition 3.2.1). If $P = P_{f,Q}$ is $C^{\infty}(\mathbb{R}^d)$, $P|_{S_Q} = f$ is necessarily $C^{\infty}(S_Q)$. It follows that $P \notin C^{\infty}(\mathbb{R}^d)$ whenever f is chosen from $C^0(S_Q) \setminus C^{\infty}(S_Q)$. By precisely the same argument, we see that $P \in C^0(\mathbb{R}^d) \setminus C^k(\mathbb{R}^d)$ whenever $f \in C^0(S_Q) \setminus C^k(S_Q)$ for each $k \in \mathbb{N}_+$.

As a straightforward example, consider $Q(x,y)=|(x,y)|=\sqrt{x^2+y^2}$ on \mathbb{R}^2 with $S_Q=\mathbb{S}$ and define

$$f(x,y) = w(\operatorname{Arg}(x,y)) + 3$$

where $w: \mathbb{R} \to \mathbb{R}$ is defined by

$$w(t) = \sum_{n=0}^{\infty} 2^{-n} \cos(3^n t)$$

for $t \in \mathbb{R}$; w is a continuous 2π -periodic version of the Weierstrass function. The resulting positive homogeneous function P is continuous but nowhere differentiable. Figure 2.2 illustrates this function P alongside Q and together with their associated unital level sets. We note that $S_P \neq S_Q$ and this is generally the case unless $f \equiv 1$.

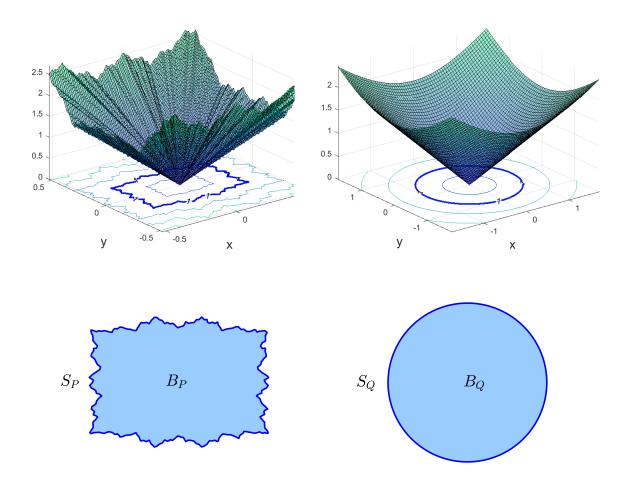


Figure 2.2: The left column illustrates P's graph and associated level set S_P containing B_P . The right column illustrates Q's graph and associated level set S_Q containing B_Q .

In most situations in estimating convolution powers, we are provided with a continuous, positive-definite P with some $E \in \operatorname{Exp}(P)$. To see if P is a positive-homogeneous function, we can check if the spectrum of E is real. The following proposition captures this result.

Proposition 2.0.7. *If* P *is continuous, positive definite and* Exp(P) *contains an* $E \in End(\mathbb{R}^d)$ *with real spectrum, then* $\{r^E\}$ *is contracting and hence all of the above conditions are (simultaneously) met.*

Proof. Since $\operatorname{Spec}(E)$ is real, the characteristic polynomial of E factors completely over $\mathbb R$ and so we may apply the Jordan-Chevalley decomposition to write E=D+N where D is diagonalizable, N is nilpotent, and DN=ND. Let $v_1,v_2,\ldots,v_d\in\mathbb R^d$ be an eigenbasis of D whose corresponding eigenvalues $\lambda_1,\lambda_2,\ldots,\lambda_d$ satisfy $\lambda_k\leq \lambda_{k+1}$ for all $k=1,2\ldots,d-1$.

Let us assume, to reach a contradiction, that $\{r^E\}$ is not contracting. Repeating the same argument given in $(b) \Rightarrow (c)$ in the proof of Proposition 2.0.5, leaves us with only one possibility:

 \triangle

There is a non-zero $x = \sum_{i=1}^d \alpha_i v_i \in \mathbb{R}^d$, and a sequence $r_k \to 0$ for which $|r_k^E x| \to \infty$. Let n+1 denote the index of N, then we have

$$r_{k}^{E}x = r_{k}^{N+D}x = r_{k}^{N}r_{k}^{D}x = \sum_{j=0}^{n}\sum_{i=1}^{d}\frac{r_{k}^{\lambda_{i}}(\log r_{k})^{j}}{j!}\alpha_{i}N^{j}v_{i}$$

for all k. Since $|r_k^E x| \to \infty$ and $r_k \to 0$, at least one eigenvalue of D must be non-positive. To see this, suppose $\lambda_i > 0$ for all i = 1, 2, ..., d, then in view of L'Hôpital's rule we have

$$\lim_{r_k \to 0} (\log r_k)^j r_k^{\lambda_i} = 0$$

for any $j=0,1,2,\ldots,n$ and $i=1,2,\ldots,d$, which implies that $\left|r_k^Ex\right| \not\to \infty$ as $r_k\to 0$, contradicting our assumption. Thus, $\lambda_1=\min\{\operatorname{Spec}(D)\}\le 0$. Let k be such that $N^kv_1\ne 0$ but $N^{k+1}v_1=0$, then

$$r^{E}N^{k}v_{1} = r^{D}r^{N}N^{k}v_{1} = r^{D}\sum_{j=0}^{\infty} \frac{(\log r)^{j}}{j!}N^{j}N^{k}v_{1} = r^{D}N^{k}v_{1} = N^{k}r^{D}v_{1} = r^{\lambda_{1}}N^{k}v_{1}$$

where we have used the fact that DN = ND. If $\lambda_1 = 0$, then

$$\infty = \lim_{r \to \infty} r P(N^k v_1) = \lim_{r \to \infty} P(r^D r^N N^k v_1) = \lim_{r \to \infty} P(r^0 N^k v_1) = \lim_{r \to \infty} P(N^k v_1) = P(N^k v_1)$$

which is impossible since *P* is continuous at $N^k v_1$. On the other hand, if $\lambda_1 < 0$, then

$$\infty = \lim_{r \to \infty} r P(N^k v_1) = \lim_{r \to \infty} P(r^D r^N N^k v_1) = \lim_{r \to \infty} P(r^{\lambda_1} N^k v_1) = P(0) = 0$$

which is also impossible.

Example 4. In Example 1, $\operatorname{Sym}(|\cdot|^{\alpha})$ is precisely the orthogonal group $\operatorname{O}_d(\mathbb{R})$ and $\mu_{|\cdot|^{\alpha}}=d/\alpha$. In Example 2, the symmetric set of a semi-elliptic polynomial P of the form (2.3) depends on the specific nature of the polynomial in question. Concerning the polynomials P_1 and P_2 in that example, it is easily shown that $\operatorname{Sym}(P_1)$ is the four-element dihedral group D_2 and $\operatorname{Sym}(P_2)$ is the two-element group consisting of the identity and the transformation $(x,y)\mapsto (x,-y)$. For a semi-elliptic polynomial P of the form (2.3),

$$\mu_P = |\mathbf{1} : \mathbf{n}| = \frac{1}{n_1} + \frac{1}{n_2} + \dots + \frac{1}{n_d}$$

and, in particular, $\mu_{P_1} = \mu_{P_2} = 1/2 + 1/4 = 3/4$.

Given a positive-homogeneous function P, let $\mathrm{Sym}(P)$ be the set of $O \in \mathrm{End}(\mathbb{R}^d)$ for which

$$P(Ox) = P(x)$$

for all $x \in \mathbb{R}^d$. By virtue of the positive-definiteness of P, it is easy to see that $\mathrm{Sym}(P)$ is a subgroup of $\mathrm{Gl}(\mathbb{R}^d)$. For this reason, $\mathrm{Sym}(P)$ is said to be the **symmetry group associated to** P. It turns out that $\mathrm{Sym}(P)$ is a subgroup of the orthogonal group. Later, we will show later that measures of sufficiently nice sets on S are invariant under transformation by any element of $\mathrm{Sym}(P)$, as expected.

Proposition 2.0.8. For each positive-homogeneous function P, $\operatorname{Sym}(P)$ is a compact subgroup of $\operatorname{Gl}(\mathbb{R}^d)$. In particular, it is a subgroup of the orthogonal group $O(\mathbb{R}^d)$.

Proof. By virtue of the Heine-Borel theorem (and the fact that $Gl(\mathbb{R}^d)$ is finite dimensional), we prove that Sym(P) is closed and bounded. To this end, let $\{O_n\} \subseteq Sym(P)$ be a sequence converging to $O \in Gl(\mathbb{R}^d)$. For each $x \in \mathbb{R}^d$, the continuity of P guarantees that

$$P(Ox) = P\left(\lim_{n \to \infty} O_n x\right) = \lim_{n \to \infty} P(O_n x) = \lim_{n \to \infty} P(x) = P(x).$$

Hence, $O \in \text{Sym}(P)$ and so Sym(P) is closed.

We assume, to reach a contradiction, that $\operatorname{Sym}(P)$ is not bounded. In this case, there is a sequence $\{\eta_n\}\subseteq \mathbb{S}$ for which $\lim_{n\to\infty}|O_n\eta_n|=\infty$. Given that \mathbb{S} is compact, by passing to a subsequence if needed, we may assume without loss in generality that $\lim_{n\to\infty}\eta_n=\eta\in\mathbb{S}$. By virtue of Proposition 2.0.5 and the continuity of P,

$$P(\eta) = \lim_{n \to \infty} P(\eta_n) = \lim_{n \to \infty} P(O_n \eta_n) = \infty$$

which is impossible. Hence Sym(P) is bounded.

Corollary 2.0.9. Let P be a positive-homogeneous function, then for all $E, E' \in \text{Exp}(P)$,

$$\operatorname{tr} E = \operatorname{tr} E' > 0.$$

Proof. By virtue of Propositions 2.0.4 and 2.0.5, $\operatorname{tr} E > 0$ for all $E \in \operatorname{Exp}(P)$. It remains to show that the trace map is constant on $\operatorname{Exp}(P)$. To this end, let $E, E' \in \operatorname{Exp}(P)$. Then, for all r > 0 and $x \in \mathbb{R}^d$,

$$P(x) = r(1/r)P(x) = rP((1/r)^{E'}x) = P(r^{E}(1/r)^{E'}x) = P(r^{E}r^{-E'}x).$$

Thus $O_r = r^E r^{-E'} \in \text{Sym}(P)$ for each r > 0. In view of the Propositions 6.1.1 and 2.0.8 and the homomorphism property of the determinant,

$$1 = \det(O_r) = \det(r^E r^{E'}) = \det(r^E) \det(r^{-E'}) = r^{\operatorname{tr} E} r^{\operatorname{tr} E} = r^{\operatorname{tr} E - \operatorname{tr} E'}$$

for all r > 0 and therefore $\operatorname{tr} E = \operatorname{tr} E'$.

In view of the preceding corollary, to each positive-homogeneous function P, we define the *homogeneous order of* P to be the unique positive number μ_P for which

$$\mu_P = \operatorname{tr} E$$

for all $E \in \text{Exp}(P)$.

We end this section by addressing a useful proposition which connects Exp(P) and Sym(P). Given $O \in \text{Sym}(P)$, we write

$$OF = \{O\eta : \eta \in F\}.$$

While a small result, the following proposition will be useful when we show that the measure on S (which we will construct from some $E \in \text{Exp}(P)$) is in fact independent of the choice of $E \in \text{Exp}(P)$. In other words, we will show that the measure is intrinsic to P.

Proposition 2.0.10. *For any* $O \in \text{Sym}(P)$

$$\operatorname{Exp}(P) = O^* \operatorname{Exp}(P)O.$$

In other words, the set Exp(P) *is invariant under conjugation by* Sym(P)*.*

Proof. Since $\operatorname{Sym}(P)$ is a subgroup of the orthogonal group $O(\mathbb{R}^d)$, $O^* = O^{-1} \in \operatorname{Sym} P$ whenever $O \in \operatorname{Sym}(P)$. For a fixed $O \in \operatorname{Sym}(P)$, it is easy to see that the map $\operatorname{Exp}(P) \ni E \mapsto O^*EO \in \operatorname{Exp}(P)$ is a bijection and hence $\operatorname{Exp}(P) = O^*\operatorname{Exp}(P)O$.

2.1 Subhomogeneous Functions

In this section, we introduce the notions of subhomogeneous functions and strongly subhomogeneous functions with respect to a given endomorphism $E \in \operatorname{End}(\mathbb{R}^d)$. As briefly discussed in the beginning of this chapter, these functions show up in the Taylor expansion of Γ and play an important role in estimating $|\phi^{(n)}|$ using the Van der Corput lemma. Roughly speaking, for some $E \in \operatorname{End}(\mathbb{R}^d)$, a sufficiently nice function $Q: \mathbb{R}^d \to \mathbb{C}$ is said to be subhomogeneous with respect to E if $Q(r^E\xi) = o(r)$ as $r \to 0$. If $|\partial_r Q(r^E\xi)| = o(1)$ as $r \to 0$, then Q is said to be strongly subhomogeneous with respect to E. It is then straightforward to use the mean value theorem to show that strong subhomogeneity implies subhomogeneity. The following definition and proposition make precise these ideas.

Definition 2.1.1. Let Q be a continuous and complex-valued function defined on an open neighborhood \mathcal{O} of 0 in \mathbb{R}^d and let $E \in \operatorname{End}(\mathbb{R}^d)$ be such that $\{r^E\}$ is a contracting group.

1. We say that Q is **subhomogeneous** with respect to E if, for each $\epsilon > 0$ and compact set $K \subseteq \mathbb{R}^d$, there is a $\delta > 0$ for which

$$\left|Q(r^E\xi)\right| \le \epsilon r$$

for all $0 < r < \delta$ and $\xi \in K$.

2. Given $k \ge 1$, we say that Q is **strongly subhomogeneous with respect to** E **of order** k if $Q \in C^k(\mathcal{O})$ and, for each $\epsilon > 0$ and compact set $K \subseteq \mathbb{R}^d$, there is a $\delta > 0$ for which

$$|r^j \partial_r^j Q(r^E \xi)| \le \epsilon r$$

for all
$$j = 1, 2, ..., k$$
, $0 < r < \delta$ and $\xi \in K$.

When the endomorphism E is understood (and fixed), we will say that Q is subhomogeneous if it is subhomogeneous with respect to E. Also, we will say that Q is k-strongly subhomogeneous if it is strongly subhomogeneous with respect to E of order k.

The following proposition, in particular, justifies our choice of vocabulary and give credence to the interpretation that 0-strongly subhomogeneous is synonymous with subhomogeneous.

Proposition 2.1.2. Let Q be differentiable on an open neighborhood of 0 in \mathbb{R}^d and let $E \in \operatorname{End}(\mathbb{R}^d)$ be such that $\{r^E\}$ is a contracting group. If, for some $k \geq 1$, Q is strongly subhomogeneous with respect to E of order k and Q(0) = 0, then Q is subhomogeneous with respect to E.

Proof. Let $\epsilon > 0$ and K be a compact set. In view of our supposition that Q is strongly subhomogeneous with respect to E or order k, let $\delta > 0$ be given so that $\left| r \partial_r Q(r^E \xi) \right| \leq \epsilon r$ for all $\xi \in K$ and $0 < r < \delta$. Given that r^E is a contracting group and Q(0) = 0, it follows that, for each $\xi \in K$, $f_{\xi} : [0, \delta) \to \mathbb{C}$ defined by

$$f_{\xi}(r) = \begin{cases} Q(r^{E}\xi) & 0 < r < \delta \\ 0 & r = 0 \end{cases}$$

is differentiable on $(0, \delta)$ and continuous on $[0, \delta)$, for each $\xi \in K$. Consequently, for every $0 < r < \delta$ and $\xi \in K$, the mean value theorem guarantees a $c = c_{\xi,r} \in (0,r)$ for which

$$|f_{\xi}(r) - f_{\xi}(0)| \le r |f'_{\xi}(c)| = r |\partial_r Q(r^E \xi)|_{r=c} \le r\epsilon.$$

Consequently, for all $0 < r < \delta$ and $\xi \in K$,

$$|Q(r^E \xi)| = |f_{\xi}(r) - f_{\xi}(0)| \le r\epsilon.$$

When estimating convolution powers, we are often presented with a Q that is (strongly) subhomogeneous with respect to some $G=E/k\in \operatorname{End}\left(\mathbb{R}^d\right)$ where k>0 and E belongs to the exponent of some positive-homogeneous or imaginary-homogeneous function P. It is then useful to characterize how the notions of subhomogeneity and strong homogeneity with respect to E/k depend on k.

Proposition 2.1.3. Let P be positive homogeneous and \widetilde{P} be complex-valued and continuous on a neighborhood of 0 in \mathbb{R}^d . The following are equivalent:

- (a) $\widetilde{P}(\xi) = o(P(\xi))$ as $\xi \to 0$.
- (b) For every $E \in \text{Exp}(P)$, \widetilde{P} is subhomogeneous with respect to E.
- (c) There exists $E \in \text{Exp}(P)$ for which \widetilde{P} is subhomogeneous with respect to E.

Proof. (a) \Rightarrow (b). Let $\epsilon > 0$, K be a compact set and choose $E \in \text{Exp}(P)$. Given our supposition that $\widetilde{P}(\xi) = o(P(\xi))$ as $\xi \to 0$, we can find an open neighborhood \mathcal{O} of 0 for which

$$\left|\widetilde{P}(\xi)\right| \le \frac{\epsilon}{1 + \sup_{\eta \in K} P(\eta)} P(\xi)$$

for all $\xi \in \mathcal{O}$. Now, because r^E is contracting in view of Proposition 2.0.5, we can find a $\delta > 0$ for which $r^E \xi \in \mathcal{O}$ for all $0 < r < \delta$ and $\xi \in K$ by virtue of Proposition 6.1.8. Consequently, for all $0 < r < \delta$ and $\xi \in K$,

$$\left|\widetilde{P}(r^E\xi)\right| \leq \frac{\epsilon}{1+\sup_{\eta \in K} P(\eta)} P(r^E\xi) = \epsilon r \frac{P(\xi)}{1+\sup_{\eta \in K} P(\eta)} \leq r\epsilon.$$

//

 $(b) \Rightarrow (c)$. This implication is trivial.

//

 $(c) \Rightarrow (a)$. Let $\epsilon > 0$. Choose $E \in \text{Exp}(P)$ and let $S = \{ \eta \in \mathbb{R}^d : P(\eta) = 1 \}$. Using the supposition that \widetilde{P} is subhomogeneous with respect to E, we may choose $\delta > 0$ for which

$$\left| \widetilde{P}(r^E \eta) \right| \le \epsilon r$$

for all $0 < r < \delta$ and $\eta \in S$. We remark that, in view of the continuity of \widetilde{P} and the fact that r^E is contracting, this inequality ensures that $\widetilde{P}(0) = 0$. We fix \mathcal{O} to be the open set $B_{\delta} = \{\xi \in \mathbb{R}^d : P(\xi) < \delta\}$. For each non-zero $\xi \in \mathcal{O}$, we observe that $\xi = r^E \eta$ where $0 < r = P(\xi) < \delta$ and $\eta = P(\xi)^{-E} \xi \in S$ and therefore

$$\left|\widetilde{P}(\xi)\right| = \left|\widetilde{P}(r^E\eta)\right| \le r\epsilon = \epsilon P(\xi)$$

If $\xi=0$, obviously, $\left|\widetilde{P}(\xi)\right|=0=\epsilon P(0)=\epsilon P(\xi)$. Thus, for all $\xi\in\mathcal{O}$,

$$\left| \widetilde{P}(\xi) \right| \le \epsilon P(\xi),$$

as desired.

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Proposition 2.1.4. Let $E \in \operatorname{End}(\mathbb{R}^d)$ be for which $\{r^E\}$ is contracting and suppose that Q is strongly subhomogeneous with respect to E of order 2. Given $\alpha > 0$, set $F = \alpha E$. Then, for any $\epsilon > 0$ and compact set K,

$$|\theta \partial_{\theta} Q(\theta^F \eta)| \le \epsilon \theta^{\alpha}$$

and

$$\left|\theta^2 \partial_{\theta}^2 Q(\theta^F \eta)\right| \le \epsilon \theta^{\alpha}$$

for all $0 < \theta \le \delta^{1/\alpha}$ and $\eta \in K$.

Proof. Let $\epsilon>0$ and $K\subseteq\mathbb{R}^d$ be a compact set. By virtue of the strong subhomogeneity of Q, let $\delta>0$ be such that

$$|r\partial_r Q(r^E\eta)| \le \epsilon' r$$
 and $|r^2\partial_r^2 Q(r^E\eta)| \le \epsilon' r$

for $0 < r < \delta$ and $\eta \in K$ where $\epsilon' = \epsilon/(2\alpha^2 + \alpha)$. We set $r = \theta^{\alpha}$ so that $r^E = \theta^F$ and observe that

$$\left|\theta \partial_{\theta} Q(\theta^{F} \eta)\right| = \left|\theta \partial_{r} Q(r^{E} \eta) \frac{\partial r}{\partial \theta}\right| = \left|\theta \partial_{r} Q(r^{E} \eta) \alpha \theta^{\alpha - 1}\right| = \alpha r \left|\partial_{r} Q(r^{E} \eta)\right| \leq \alpha \epsilon' r < \epsilon \theta^{\alpha}$$

for all $0 < \theta \le \delta^{1/\alpha}$ and $\eta \in K$. Further, we have

$$\begin{split} \left| \theta^2 \partial_{\theta}^2 Q(\theta^F \eta) \right| &= \theta^2 \middle| \partial_r^2 Q(r^E \eta) \left(\frac{\partial r}{\partial \theta} \right)^2 + \partial_r Q(r^E \eta) \frac{\partial^2 r}{\partial \theta^2} \middle| \\ &= \theta^2 \middle| \partial_r^2 Q(r^E \eta) \alpha^2 \theta^{2\alpha - 2} + \partial_r Q(r^E \eta) \alpha(\alpha - 1) \theta^{\alpha - 2} \middle| \\ &\leq \alpha^2 \middle| r^2 Q(r^E \eta) \middle| + |\alpha(\alpha - 1)| \middle| r \partial_r Q(r^E \eta) \middle| \\ &< \alpha^2 \epsilon' r + \middle| \alpha^2 - \alpha \middle| \epsilon' r \\ &< \epsilon \theta^{\alpha} \end{split}$$

 $\text{ for all } 0<\theta<\delta^{1/\alpha} \text{ and } \eta\in K.$

Chapter 3

A generalized polar integration formula

In this chapter, we obtain our generalized polar integration formula from two perspectives: measure theory and smooth-manifold theory. In Subsection 3.1.1, we construct a measure on the unital level set of a positive homogeneous function. Our construction only uses results from point-set topology and measure theory and this measure will appear later as the surface measure which appears in our generalization of the polar coordinate integration formula. In a later section (Section 3.2), we will consider the case in which the positive homogeneous function is smooth and, in this case, we will find that its unital level set is a smooth hypersurface of \mathbb{R}^d and the measure constructed in Subsection 3.1.1 is characterized by a smooth d-1 form which related to the canonical Riemann volume form.

To start, we let a positive-homogeneous function P on \mathbb{R}^d be fixed with homogeneous order μ_P , exponent set $\operatorname{Exp}(P)$, and symmetry group $\operatorname{Sym}(P)$. Let \mathcal{M}_d denote the Lebesgue σ -algebra on $\mathbb{R}^d \setminus \{0\}$ and m the Lebesgue measure. We take the unital level set $S = \{\eta \in \mathbb{R}^d : P(\eta) = 1\}$ to be equipped with the relative topology inherited from \mathbb{R}^d and $(0,\infty) \times S$ with the product topology, where $(0,\infty)$ carries the usual topology on \mathbb{R} . Let $\mathcal{B}(S)$ denote the σ -algebra on S. We denote by $\mathcal{L} = \mathcal{L}(0,1)$ the σ -algebra of Lebesgue measurable sets on $(0,\infty)$ and let λ_P be the σ -finite measure on $((0,\infty),\mathcal{L})$ with $\lambda_P(dr) = r^{\mu_P-1} dr$. For basic notions of measure theory, the reader is encouraged to refer to [3], [15]. For selected definitions and theorems from measure theory, see Appendix 6.2.

3.1 A surface measure on S

Our main result of this section is as follows.

Theorem 3.1.1. There exists a σ -algebra Σ on S containing $\mathcal{B}(S)$ and a finite Radon measure σ_P on (S, Σ) which satisfies the following properties:

- 1. (S, Σ, σ_P) is the completion of $(S, \mathcal{B}(S), \sigma_P)$. In particular, (S, Σ, σ_P) is a complete measure space.
- 2. For any $F \in \Sigma_P$ and $O \in \operatorname{Sym}(P)$, $OF \in \Sigma_P$ and $\sigma_P(OF) = \sigma_P(F)$.

3. For any $F \in \Sigma$ and $E \in \text{Exp}(P)$,

$$\widetilde{F_E} := \bigcup_{0 < r < 1} \left(r^E F \right) = \left\{ r^E \eta \in \mathbb{R}^d \setminus \{0\} : 0 < r < 1, \eta \in F \right\}$$

is a Lebesgue measurable subset of $\mathbb{R}^d \setminus \{0\}$, i.e., $\widetilde{F_E} \in \mathcal{M}_d$, and

$$\sigma_P(F) = \mu_P \cdot m\left(\widetilde{F_E}\right)$$

where m denotes the Lebesgue measure on \mathbb{R}^d .

Further, denote by $((0, \infty) \times S, (\mathcal{L} \times \Sigma)', \lambda_P \times \sigma_P)$ the completion of the product measure space $((0, \infty) \times S, \mathcal{L} \times \Sigma, \lambda_P \times \sigma_P)$. We have

1. Given any $E \in \text{Exp}(P)$, the map $\psi_E : (0, \infty) \times S \to \mathbb{R}^d \setminus \{0\}$, defined by $\psi_E(r, \eta) = r^E \eta$ for r > 0 and $\eta \in S$, is a point isomorphism of the measure spaces $((0, \infty) \times S, (\mathcal{L} \times \Sigma)', \lambda_P \times \sigma_P)$ and $(\mathbb{R}^d \setminus \{0\}, \mathcal{M}_d, m)$. That is

$$\mathcal{M}_d = \left\{ A \subseteq \mathbb{R}^d \setminus \{0\} : \psi_E^{-1}(A) \in (\mathcal{L} \times \Sigma)' \right\}$$

and, for each $A \in \mathcal{M}_d$,

$$m(A) = (\lambda_P \times \sigma_P)(\psi_E^{-1}(A)).$$

- 2. Given any Lebesgue measurable function $f : \mathbb{R}^d \to \mathbb{C}$ and $E \in \operatorname{Exp}(P)$, $f \circ \psi_E$ is $(\mathcal{L} \times \Sigma)'$ -measurable and the following statements hold:
 - (a) If $f \geq 0$, then

$$\int_{\mathbb{R}^d} f(x) dx = \int_0^\infty \left(\int_S f(r^E \eta) \, \sigma_P(d\eta) \right) r^{\mu_P - 1} dr$$

$$= \int_S \left(\int_0^\infty f(r^E \eta) r^{\mu_P - 1} dr \right) \sigma_P(d\eta). \tag{3.1}$$

(b) When f is complex-valued, we have

$$f \in L^1(\mathbb{R}^d)$$
 if and only if $f \circ \psi_E \in L^1\left((0,\infty) \times S, (\mathcal{L} \times \Sigma)', \lambda_P \times \sigma_P\right)$ and, in this case, (3.1) holds.

From this theorem we have a useful corollary.

Corollary 3.1.2. Given $g: S \to \mathbb{C}$ and $E \in \text{Exp}(P)$, define $f: \mathbb{R}^d \to \mathbb{C}$ by

$$f(x) = \begin{cases} \mu_P \cdot \chi_{(0,1)}(P(x))g(P(x)^{-E}x) & \text{for } x \neq 0\\ 0 & \text{for } x = 0 \end{cases}$$

for $x \in \mathbb{R}^d$ where $\chi_{(0,1)}(\cdot)$ is the indicator function of the interval (0,1). Then $g \in L^1(S, \Sigma_P, \sigma_P)$ if and only if $f \in L^1(\mathbb{R}^d)$ and, in this case,

$$\int_{\mathbb{R}^d} f(x) \, dx = \int_S g(\eta) \sigma_P(d\eta). \tag{3.2}$$

Proof. Observe that, for the $(\mathcal{L} \times \Sigma_P)'$ -measurable function $k(r, \eta) = \chi_{(0,1)}(r)g(\eta)$,

$$k \circ \psi_E^{-1}(x) = k(P(x), P(x)^{-E}x) = f(x)$$

for $x \in \mathbb{R}^d \setminus \{0\}$. By virtue of Theorem 3.1.1, it follows that f is Lebesgue measurable and

$$\int_{\mathbb{R}^d} |f(x)| \, dx = \int_{S} \left(\int_{(0,\infty)} |k(r,\eta)| r^{\mu_P - 1} \, dr \right) \sigma_P(d\eta)
= \left(\int_{S} |g(\eta)| \sigma_P(d\eta) \right) \left(\int_{0}^{1} \mu_P r^{\mu_P - 1} \, dr \right)
= \int_{S} |g(\eta)| \sigma_P(d\eta)$$

and therefore $f \in L^1(\mathbb{R}^d)$ if and only if $g \in L^1(S, \Sigma_P, \sigma_P)$ and $||f||_{L^1(\mathbb{R}^d)} = ||g||_{L^1(S)}$. By an analogous computation (for f instead of |f|), we have

$$\int_{\mathbb{R}^d} f(x) \, dx = \int_S g(\eta) \sigma_P(d\eta),$$

by virtue of Property 2 of Theorem 3.1.1.

Our construction of σ_P is as follows. In Subsection 3.1.1, we fix $E \in \operatorname{Exp}(P)$ and consider the one-parameter contracting group $\{r^E\}$. As the standard isotropic one-parameter group $r \mapsto rI = r^I$ is well-fitted to the unit sphere $\mathbb S$ and allows every non-zero $x \in \mathbb R^d$ to be written uniquely as $x = r\eta$ for $r \in (0,\infty)$ and $\eta \in \mathbb S$, $\{r^E\}$ is well-fitted to S and has the property that every non-zero $x \in \mathbb R^d$ can be written uniquely as $x = r^E\eta$ where $r \in (0,\infty)$ and $\eta \in S$. With this one-parameter group as a tool, we define a surface-carried measure $\sigma_{P,E}$ on S by taking sufficiently nice sets $F \subseteq S$, stretching them into a quasi-conical region of the associated "ball" B with the contracting group $\{r^E\}$, and computing the Lebesgue measure of the result. Figure 3.1 illustrates the action of r^E for $r \in (0,1)$ on a set F (red) in S (blue) to create the set F. Here, we choose $P(x,y) = x^2 + 3xy^2/2 + y^4$ and $E = \operatorname{diag}(1/2,1/4)$. Figure 3.2 illustrates a different example with $P(x,y,z) = x^2 + xy^2 + y^4 + z^4$.

In Subsection 3.1.2, we turn our focus to an associated product measure $\lambda_P \times \sigma_{P,E}$ on $(0,\infty) \times S$ with which we are able to formulate and prove a generalization of (??); this is Theorem 3.1.8. We then derive a number of corollaries of Theorem 3.1.8, including the result that $\sigma_{P,E}$ is a Radon measure on S. In Subsection 3.1.3, we prove that even though Subsections 3.1.1 and 3.1.2 require a choice of $E \in \operatorname{Exp}(P)$, $\sigma_{P,E}$ is independent of the choice of $E \in \operatorname{Exp}(P)$, and so we write $\sigma_P = \sigma_{P,E}$ and this quickly yields a stronger version of Theorem 3.1.8, which is Theorem 3.1.1.

3.1.1 Construction of $\sigma_{P,E}$

Let
$$E\in \operatorname{Exp}(P)$$
. Define $\psi_E:(0,\infty)\times S\to\mathbb{R}^d\setminus\{0\}$ by
$$\psi_E(r,\eta)=r^E\eta \tag{3.3}$$

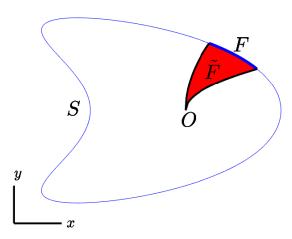


Figure 3.1: Quasi-conical region $\widetilde{F}=\widetilde{F}_E$ (in red) for $F\subseteq S$. Here, S is the unital level set of $P=P_2$ from Example 2 and $E\in \operatorname{Exp}(P_2)$ has the standard representation $\operatorname{diag}(1/2,1/4)$.

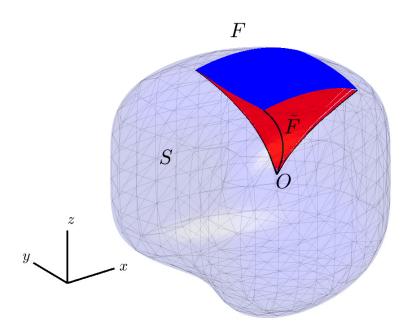


Figure 3.2: Quasi-conical region $\widetilde{F}=\widetilde{F}_E$ (in red) for $F\subseteq S$. Here, S is the unital level set of $P(x,y,z)=x^2+xy^2+y^4+z^4$ and $E\in \operatorname{Exp}(P)$ has standard representation $\operatorname{diag}(1/2,1/4,1/4)$.

for r>0 and $\eta\in S$. As ψ_E is the restriction of the continuous function $(0,\infty)\times\mathbb{R}^d\ni (r,x)\mapsto r^Ex\in\mathbb{R}^d$ to $(0,\infty)\times S$, it is necessarily continuous. We now show that ψ_E is a homeomorphism. This will be particularly useful later on when we are required to regularly "move" back and forth between $\mathbb{R}^d\setminus\{0\}$ and $(0,\infty)\times S$.

Proposition 3.1.3. The map $\psi_E:(0,\infty)\times S\to \mathbb{R}^d\setminus\{0\}$ is a homeomorphism with continuous inverse $\psi_E^{-1}:\mathbb{R}^d\setminus\{0\}\to(0,\infty)\times S$ given by

$$\psi_E^{-1}(x) = (P(x), (P(x))^{-E}x)$$

for $x \in \mathbb{R}^d \setminus \{0\}$.

Proof. Given that P is continuous and positive definite, P(x) > 0 for each $x \in \mathbb{R}^d \setminus \{0\}$ and the map $\mathbb{R}^d \setminus \{0\} \ni x \mapsto (P(x))^{-E}x \in \mathbb{R}^d$ is continuous. Further, in view of the homogeneity of P,

$$P((P(x))^{-E}x) = P(x)^{-1}P(x) = 1$$

for all $x \in \mathbb{R}^d \setminus \{0\}$. It follows from these two observations that

$$\rho(x) = (P(x), (P(x))^{-E}x),$$

defined for $x \in \mathbb{R}^d \setminus \{0\}$, is a continuous function taking $\mathbb{R}^d \setminus \{0\}$ into $(0, \infty) \times S$. We have

$$(\psi_E \circ \rho)(x) = \psi_E(P(x), (P(x))^{-E}x) = (P(x))^E(P(x))^{-E}x = x$$

for every $x \in \mathbb{R}^d \setminus \{0\}$ and

$$(\rho \circ \psi_E)(r,\eta) = \rho(r^E \eta) = (P(r^E \eta), (P(r^E \eta))^{-E}(r^E \eta)) = (rP(\eta), (rP(\eta))^{-E}(r^E \eta)) = (r,\eta)$$

for every $(r, \eta) \in (0, \infty) \times S$. Thus ρ is a (continuous) inverse for ψ_E and so it follows that ψ_E is a homeomorphism and $\rho = \psi_E^{-1}$.

As in the statement of Theorem 3.1.1, for each $F \subseteq S$, define

$$\widetilde{F_E} = \bigcup_{0 < r < 1} \left(r^E F \right) = \{ r^E \eta : 0 < r < 1, \eta \in F \},$$

which we can identify as a quasi-conical region discussed in the previous chapter. Next, we let $\Sigma_{P,E}$ be the collection of subsets F of S for which $\widetilde{F_E} \in \mathcal{M}_d$, i.e.,

$$\Sigma_{P,E} = \{ F \subseteq S : \widetilde{F_E} \in \mathcal{M}_d \}.$$

We now show that this is a σ -algebra $\Sigma_{P,E}$ on S.

Proposition 3.1.4. $\Sigma_{P,E}$ is a σ -algebra on S containing the Borel σ -algebra on S, $\mathcal{B}(S)$.

Proof. Throughout the proof, we write $\Sigma = \Sigma_{P,E}$ and $\widetilde{F} = \widetilde{F_E}$ for each $F \subseteq S$. We first show that Σ is a σ -algebra. Since $\widetilde{S} = B \setminus \{0\}$, it is open in $\mathbb{R}^d \setminus \{0\}$ and therefore Lebesgue measurable. Hence $S \in \Sigma$. Let $G, F \in \Sigma$ be such that $G \subseteq F$. Then,

$$\widetilde{F \setminus G} = \bigcup_{0 < r < 1} r^E \left(F \setminus G \right) = \bigcup_{0 < r < 1} \left(r^E F \setminus r^E G \right) = \left(\bigcup_{0 < r < 1} r^E F \right) \setminus \left(\bigcup_{0 < r < 1} r^E G \right) = \widetilde{F} \setminus \widetilde{G}$$

where we have used the fact that the collection $\{r^E F\}_{0 < r < 1}$ is mutually disjoint to pass the union through the set difference. Consequently $\widetilde{F} \setminus \widetilde{G}$ is Lebesgue measurable and therefore $F \setminus G \in \Sigma$. Now, given a countable collection $\{F_n\} \subseteq \Sigma$, observe that

$$\bigcup_{n=1}^{\infty} F_n = \bigcup_{0 < r < 1} r^E \left(\bigcup_{n=1}^{\infty} F_n \right) = \bigcup_{0 < r < 1} \bigcup_{n=1}^{\infty} r^E F_n = \bigcup_{n=1}^{\infty} \bigcup_{0 < r < 1} r^E F_n = \bigcup_{n=1}^{\infty} \widetilde{F_n} \in \mathcal{M}_d$$

whence $\cup_n F_n \in \Sigma$. Thus Σ is a σ -algebra.

Finally, we show that

$$\mathcal{B}(S) \subseteq \Sigma$$
.

As the Borel σ -algebra is the smallest σ -algebra containing the open subsets of S, it suffices to show that $\mathcal{O} \in \Sigma$ whenever \mathcal{O} is open in S. Armed with Proposition 3.1.3, this is an easy task: Given an open set $\mathcal{O} \subseteq S$, observe that

$$\widetilde{\mathcal{O}} = \{ r^E \eta : 0 < r < 1, \eta \in \mathcal{O} \} = \psi_E((0, 1) \times \mathcal{O}).$$

Upon noting that $(0,1) \times \mathcal{O}$ is an open subset of $(0,\infty) \times S$, Proposition 3.1.3 guarantees that $\widetilde{\mathcal{O}} = \psi_E((0,1) \times \mathcal{O}) \subseteq \mathbb{R}^d \setminus \{0\}$ is open and therefore $\widetilde{\mathcal{O}} \in \mathcal{M}_d$. Thus, $\mathcal{O} \in \Sigma$.

With the σ -algebra $\Sigma_{P,E}$, we now specify a measure on the measurable space $(S, \Sigma_{P,E})$. For each $F \in \Sigma_{P,E}$, define $\sigma_{P,E} : \Sigma_{P,E} \to [0,\infty)$ by

$$\sigma_{P,E}(F) = \mu_P \cdot m(\widetilde{F_E})$$

where m is the Lebesgue measure on \mathbb{R}^d and $\mu_P = \operatorname{tr} E > 0$ is the homogeneous order associated to P.

Proposition 3.1.5. $\sigma_{P,E}$ is a finite measure on $(S, \Sigma_{P,E})$.

Proof. Throughout the proof, we will write $\sigma = \sigma_{P,E}$, $\Sigma = \Sigma_{P,E}$, and, $\widetilde{F} = \widetilde{F_E}$ for each $F \subseteq S$. It is clear that σ is non-negative and $\sigma(\varnothing) = 0$. Let $\{F_n\}_{n=1}^{\infty} \subseteq \Sigma$ be a mutually disjoint collection. We claim that $\{\widetilde{F_n}\}_{n=1}^{\infty} \subseteq \mathcal{M}_d$ is also a mutually disjoint collection. To see this, suppose that $x = r_n^E \eta_n = r_m^E \eta_m \in \widetilde{F_n} \cap \widetilde{F_m}$, where $r_n, r_m \in (0,1)$, $\eta_n \in F_n$, and $\eta_m \in F_m$. Then

$$r_n = P(r_n^E \eta_n) = P(x) = P(r_m^E \eta_m) = r_m,$$

implying that $\eta_n = \eta_m \in F_n \cap F_m$. Because $\{F_n\}_{n=1}^{\infty}$ is mutually disjoint, we must have n=m which verifies our claim. By virtue of the countable additivity of Lebesgue measure, we therefore have

$$\sigma\left(\bigcup_{n=1}^{\infty} F_n\right) = \mu_P \cdot m\left(\bigcup_{n=1}^{\infty} F_n\right) = \mu_P \cdot m\left(\bigcup_{n=1}^{\infty} \widetilde{F_n}\right) = \mu_P \sum_{n=1}^{\infty} m(\widetilde{F_n}) = \sum_{n=1}^{\infty} \sigma(F_n).$$

Therefore σ is a measure on (S, Σ) . In view of Condition (b) of Proposition 2.0.5, $\widetilde{S} = B \setminus \{0\}$ is a bounded subset of $\mathbb{R}^d \setminus \{0\}$ and hence $\sigma(S) = \mu_P \cdot m(B \setminus \{0\}) < \infty$ showing that σ is finite. \square

Corollary 3.1.6. By virtue of the two preceding propositions, $\sigma_{P,E}$ is a finite Borel measure on S.

3.1.2 Product Measure and Point Isomorphism

So far, we have only worked with a surface-carried measure $\sigma_{P,E}$ on S. To obtain the full polar-coordinate integration formula in 3.1.1, we need to construct a **product measure** that combines the surface measure $\sigma_{P,E}$ on the unital level set S and the radial measure λ_P on \mathbb{R} to obtain a "volume" measure. Once we have the product measure, we shall use Fubini's theorem, stated below for complete (product) measure spaces, to obtain our main result for this section.

To start, let $(S, \Sigma_{P,E}, \sigma_{P,E})$ denote the finite measure space of Proposition 3.1.5. Recall that \mathcal{L} denotes the σ -algebra of Lebesgue measurable subsets of $(0, \infty)$ and λ_P denotes the measure on $(0, \infty)$ with $\lambda_P(dr) = r^{\mu_P - 1} dr$, i.e., for each $L \in \mathcal{L}$,

$$\lambda_P(L) = \int_0^\infty \chi_L(r) r^{\mu_P - 1} dr.$$

It is easy to see that λ_P is σ -finite and so, in view of the finiteness of the measure $\sigma_{P,E}$, there exists a unique product measure $\lambda_P \times \sigma_{P,E}$ on $(0,\infty) \times S$ equipped with the product σ -algebra $\mathcal{L} \times \Sigma_{P,E}$ which satisfies

$$(\lambda_P \times \sigma_{P,E})(L \times F) = \lambda_P(L)\sigma_{P,E}(F)$$

for all $L \in \mathcal{L}$ and $F \in \Sigma_{P,E}$ (see Theorem 6.2.31 in the Appendix). We shall denote by $((0,\infty) \times S, (\mathcal{L} \times \Sigma_{P,E})', \lambda_P \times \sigma_{P,E})$ the completion of the measure space $((0,\infty) \times S, \mathcal{L} \times \Sigma_{P,E}, \lambda_P \times \sigma_P)$. With this, we state (without proving) the following Fubini-Tonelli theorem associated to $\lambda_P \times \sigma_{P,E}$, which combines the results of Theorem 6.2.34 and the "measureable slice" theorem in Appendix 6.2. This allows us to exchange the order of integration, which is crucial for future applications.

Theorem 3.1.7 (Analogous to Theorem 8.12 of [14]). Let $g:(0,\infty)\times S\to \mathbb{C}$ be $(\mathcal{L}\times\Sigma_{P,E})'$ -measurable. For each $r\in(0,\infty)$, define $g^r:S\to\mathbb{C}$ by $g^r(\eta)=g(r,\eta)$ for $\eta\in S$ and, for each $\eta\in S$, define $g_\eta:(0,\infty)\to\mathbb{C}$ by $g_\eta(r)=g(r,\eta)$ for $r\in(0,\infty)$.

- 1. For λ_P -almost every r, g^r is $\Sigma_{P,E}$ -measurable and, for $\sigma_{P,E}$ -almost every η , g_{η} is \mathcal{L} -measurable.
- 2. If $g \ge 0$, then:
 - (a) For λ_P -almost every r,

$$H(r) = \int_{S} g^{r}(\eta) \, \sigma_{P,E}(d\eta)$$

exists as a non-negative extended real number.

(b) For $\sigma_{P,E}$ -almost every η ,

$$G(\eta) = \int_0^\infty g_{\eta}(r) r^{\mu_P - 1} dr$$

exists as a non-negative extended real number.

(c) We have

$$\int_0^\infty H(r)r^{\mu_P - 1}dr = \int_{(0,\infty)\times S} g \, d(\lambda_P \times \sigma_{P,E}) = \int_S G(\eta) \, \sigma_{P,E}(d\eta) \tag{3.4}$$

and, in particular,

$$\int_0^\infty \left(\int_S g(r,\eta) \, \sigma_{P,E}(d\eta) \right) r^{\mu_P - 1} \, dr = \int_S \left(\int_0^\infty g(r,\eta) r^{\mu_P - 1} \, dr \right) \, \sigma_{P,E}(d\eta). \quad (3.5)$$

3. If g is complex valued and

$$\int_{S} \left(\int_{0}^{\infty} |g(r,\eta)| r^{\mu_P - 1} \, dr \right) \, \sigma_{P,E}(d\eta) < \infty \ \, \text{or} \ \, \int_{0}^{\infty} \left(\int_{S} |g(r,\eta)| \, \sigma_{P,E}(d\eta) \right) r^{\mu_P - 1} \, dr < \infty,$$

then $g \in L^1((0,\infty) \times S, (\mathcal{L} \times \Sigma_{P,E})', \lambda_P \times \sigma_{P,E}).$

4. If $g \in L^1((0,\infty) \times S, (\mathcal{L} \times \Sigma_{P,E})', \lambda_P \times \sigma_{P,E})$, then $g^r \in L^1(S, \Sigma_{P,E}, \sigma_{P,E})$ for λ_P -almost every $r, g_\eta \in L^1((0,\infty), \mathcal{L}, \lambda_P)$ for $\sigma_{P,E}$ -almost every η , and Equations (3.4) and (3.5) hold.

Our primary goal in this subsection is to prove the theorem below, which is a slightly more special statement of the second half of Theorem 3.1.1.

Theorem 3.1.8. Let $((0,\infty)\times S, (\mathcal{L}\times\Sigma_{P,E})', \lambda_P\times\sigma_{P,E})$ be as above.

1. The map $\psi_E: (0,\infty) \times S \to \mathbb{R}^d \setminus \{0\}$, defined by (3.3), is a point isomorphism of the measure spaces $((0,\infty) \times S, (\mathcal{L} \times \Sigma_{P,E})', \lambda_P \times \sigma_{P,E})$ and $(\mathbb{R}^d \setminus \{0\}, \mathcal{M}_d, m)$. That is

$$\mathcal{M}_d = \{ A \subseteq \mathbb{R}^d \setminus \{0\} : \psi_E^{-1}(A) \in (\mathcal{L} \times \Sigma_{P,E})' \}$$

and, for each $A \in \mathcal{M}_d$,

$$m(A) = (\lambda_P \times \sigma_{P,E})(\psi_E^{-1}(A)).$$

- 2. If $f: \mathbb{R}^d \to \mathbb{C}$ is Lebesgue measurable, then $f \circ \psi_E$ is $(\mathcal{L} \times \Sigma_{P,E})'$ -measurable and the following statements hold:
 - (a) If $f \geq 0$, then

$$\int_{\mathbb{R}^d} f(x) dx = \int_0^\infty \left(\int_S f(r^E \eta) \, \sigma_{P,E}(d\eta) \right) r^{\mu_P - 1} dr$$

$$= \int_S \left(\int_0^\infty f(r^E \eta) r^{\mu_P - 1} dr \right) \, \sigma_{P,E}(d\eta). \tag{3.6}$$

(b) When f is complex-valued, we have

$$f \in L^1(\mathbb{R}^d)$$
 if and only if $f \circ \psi_E \in L^1((0,\infty) \times S, (\mathcal{L} \times \Sigma_{P,E})', \lambda_P \times \sigma_{P,E})$ and, in this case, (3.6) holds.

Properties 1 and 2 in Theorem 3.1.8 differ only from Properties 1 and 2 in Theorem 3.1.1 only in the apparent dependence on our choice of $E \in \operatorname{Exp}(P)$. We shall see in Subsection 3.1.3 and Proposition 3.1.16 therein that $\Sigma_{P,E_1} = \Sigma_{P,E_2}$ and $\sigma_{P,E_1} = \sigma_{P,E_2}$ for all $E_1, E_2 \in \operatorname{Exp}(P)$ and so the dependence is indeed superficial. With this result, we shall obtain Properties 1 and 2 of Theorem 3.1.1 immediately from Theorem 3.1.8.

To prove Theorem 3.1.8, we shall first treat several lemmas. These lemmas isolate and generalize several important ideas used in standard proofs of (??) (See, e.g., [6] and [15]). Lemma 3.1.9 establishes the scaling property of $\{r^E\}$ on the Lebesgue measures and measurability of sets in \mathbb{R}^d . Lemma 3.1.10 relates the Lebesgue measure of sets of the form $\psi_E(I \times F)$, where $I \subseteq (0, \infty)$ is sufficiently nice and $F \in \Sigma_{P,E}$, to the product measure $\lambda_P \times \sigma_{P,E}$.

Lemma 3.1.9. Let $A \subseteq \mathbb{R}^d$ and r > 0. A is Lebesgue measurable if and only if $r^E A = \{x = r^E a : a \in A\}$ is Lebesgue measurable and, in this case,

$$m(r^E A) = r^{\mu_P} m(A).$$

Proof. Because $x \mapsto r^E x$ is a linear isomorphism, $r^E A$ is Lebesgue measurable if and only if A is Lebesgue measurable. Observe that $x \in r^E A$ if and only if $r^{-E} x \in A$ and therefore

$$m(r^E A) = \int_{\mathbb{R}^d} \chi_{r^E A}(x) dx = \int_{\mathbb{R}^d} \chi_A(r^{-E} x) dx$$

where χ_{r^EA} and χ_A respectively denote the indicator functions of the sets r^EA and A. Now, by making the linear change of variables $x \mapsto r^Ex$, we have

$$m(r^E A) = \int_{\mathbb{R}^d} \chi_A(x) |\det(r^E)| dx = r^{\mu_P} m(A),$$

because $\det(r^E) = r^{\operatorname{tr} E} = r^{\mu_P} > 0$ by virtue of Proposition 6.1.1 and Corollary 2.0.9.

Lemma 3.1.10. Let $F \in \Sigma_{P,E}$. If $I \subseteq (0,\infty)$ is open, closed, G_{δ} , or F_{σ} , then $\psi_{E}(I \times F) \in \mathcal{M}_{d}$ and

$$m(\psi_E(I \times F)) = (\lambda_P \times \sigma_{P,E})(I \times F) = \lambda_P(I)\sigma_{P,E}(F). \tag{3.7}$$

Proof. To simplify notation, we shall write $\lambda = \lambda_P$ and $\sigma = \sigma_{P,E}$ throughout the proof. We fix $F \in \Sigma_{P,E}$ and consider several cases for I.

Case 1: I = (0, b) for $0 < b \le \infty$. When b is finite, observe that

$$\psi_E(I \times F) = \{r^E \eta : 0 < r < b, \eta \in F\} = b^E \{r^E \eta : 0 < r < 1, \eta \in F\} = b^E \widetilde{F_E}.$$

By virtue of Lemma 3.1.9, it follows that $\psi_E(I \times F) \in \mathcal{M}_d$ and

$$\begin{split} (\lambda \times \sigma)(I \times F) &= \lambda(I)\sigma(F) \\ &= \left(\int_0^b r^{\mu_P - 1} dr\right) \left(\mu_P \cdot m(\widetilde{F_E})\right) \\ &= b^{\mu_P} m(\widetilde{F_E}) \\ &= m(b^E \widetilde{F_E}) \\ &= m(\psi_E(I \times F)). \end{split}$$

When $b = \infty$ i.e., $I = (0, \infty)$, we observe that

$$I = \bigcup_{n=1}^{\infty} (0, n) = \bigcup_{n=1}^{\infty} I_n$$

where the open intervals $I_n = (0, n)$ are nested and increasing. In view of the result above (for finite b = n), we have

$$\psi_E(I \times F) = \psi_E \left(\bigcup_{n=1}^{\infty} (I_n \times F) \right) = \bigcup_{n=1}^{\infty} \psi_E(I_n \times F) \in \mathcal{M}_d.$$

Given that ψ_E is a bijection, $\{\psi_E(I_n \times F)\}$ is necessarily a nested increasing sequence and so, by the continuity of the measures $\lambda \times \sigma$ and m,

$$(\lambda \times \sigma)(I \times F) = \lim_{n \to \infty} (\lambda \times \sigma)(I_n \times F) = \lim_{n \to \infty} m(\psi_E(I_n \times F)) = m(\psi_E(I \times F)).$$

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 $/\!\!/$

Case 2: I = (0, a] for $0 < a < \infty$. We have

$$I = (0, a] = \bigcap_{n=1}^{\infty} (0, a + 1/n) = \bigcap_{n=1}^{\infty} I_n$$

where the open intervals $I_n = (0, a + 1/n)$ are nested and decreasing. By reasoning analogous to that given in Case 1, we have

$$\psi_E(I \times F) = \bigcap_{n=1}^{\infty} \psi_E(I_n \times F) \in \mathcal{M}_d$$

and

$$(\lambda \times \sigma)(I \times F) = \lim_{n \to \infty} (\lambda \times \sigma)(I_n \times F) = \lim_{n \to \infty} m(\psi_E(I_n \times F)) = m(\psi_E(I \times F)).$$

In particular, $m(\psi_E(I \times F)) = \lambda((0, a])\sigma(F) = a^{\mu_P}\sigma(F)/\mu_P < \infty$.

Case 3: I = (a, b) for $0 < a < b \le \infty$. In this case, $I = (0, b) \setminus (0, a]$ and so, in view of Cases 1 and 2, $\psi_E(I \times F) = \psi_E((0, b) \times F) \setminus \psi_E((0, a] \times F) \in \mathcal{M}_d$ and

$$(\lambda \times \sigma)(I \times F) = (\lambda \times \sigma)((0,b) \times F) - (\lambda \times \sigma)((0,a] \times F)$$
$$= m(\psi_E((0,b) \times F)) - m(\psi_E((0,a] \times F))$$
$$= m(\psi_E(I \times F))$$

where we have used the fact that $(\lambda \times \sigma)((0, a] \times F) = m(\psi_E((0, a] \times F)) < \infty$.

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Case 4: $I \subseteq (0, \infty)$ is open. In this case, it is known that I can be expressed as a countable union of disjoint open intervals $\{I_n\}$ and, by virtue of Cases 1 and 3, we have

$$\psi_E(I \times F) = \bigcup_{n=1}^{\infty} \psi_E(I_n \times F) \in \mathcal{M}_d,$$

where this union is disjoint, and

$$m(\psi_E(I \times F)) = \sum_n m(\psi_E(I_n \times F)) = \sum_n (\lambda \times \sigma)(I_n \times F)$$
$$= \sum_n \lambda(I_n)\sigma(F) = \left(\sum_n \lambda(I_n)\right)\sigma(F) = \lambda(I)\sigma(F) = (\lambda \times \sigma)(I \times F).$$

Case 5: $I \subseteq (0, \infty)$ *is closed.* In this case, we have $I = (0, \infty) \setminus O$ where O is open and so

$$\psi_E(I \times F) = \psi_E(F \times ((0, \infty) \setminus O)) = \psi_E((0, \infty) \times F) \setminus \psi_E(O \times F) \in \mathcal{M}_d.$$

At this point, we'd like to use the property that

$$m(\psi_E((0,\infty)\times F)\setminus\psi_E(O\times F))=m(\psi_E((0,\infty)\times F))-m(\psi_E(O\times F)),$$

but this only holds when $m(\psi_E(O \times F))$ is finite. We must therefore proceed differently. For each natural number n, define $O_n = O \cap (0,n)$ and $I_n = (0,n) \setminus O_n$. It is straightforward to show that $\{I_n\}$ and $\{\psi_E(I_n \times F)\}$ are nested and increasing with

$$I = \bigcup_{n=1}^{\infty} I_n$$
 and $\psi_E(I \times F) = \bigcup_{n=1}^{\infty} \psi_E(I_n \times F).$

The results of Cases 1 and 4 guarantee that, for each n,

$$m(\psi_E(O_n \times F)) = (\lambda \times \sigma)(O_n \times F) \le (\lambda \times \sigma)((0, n) \times F) < n^{\mu_P} m(\widetilde{F_E}) < \infty$$

and therefore

$$m(\psi_E(I_n \times F)) = m(\psi_E((0, n) \times F)) - m(\psi_E(O_n \times F))$$

= $(\lambda \times \sigma)((0, n) \times F) - (\lambda \times \sigma)(O_n \times F)$
= $(\lambda \times \sigma)(I_n \times F)$.

Then, by virtue of the continuity of measure,

$$m(\psi_E(I \times F)) = \lim_{n \to \infty} m(\psi_E(I_n \times F)) = \lim_{n \to \infty} (\lambda \times \sigma)(I_n \times F) = (\lambda \times \sigma)(I \times F).$$

Case 6. $I \subseteq (0, \infty)$ is G_{δ} or F_{σ} . Depending on whether I is G_{δ} or F_{σ} , express I as an intersection of nested decreasing open sets or a union of nested increasing closed sets. In both cases, by virtually the same argument given in the previous cases, we find that $\psi_E(I \times F) \in \mathcal{M}_d$,

$$m(\psi_E(I \times F)) = (\lambda \times \sigma)(I \times F).$$

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The following lemma generalizes the result of Lemma 3.1.10 (which only deals with cases of "special rectangles") so that the previous statement holds for any set $L \in \mathcal{L}$, the σ -algebra of Lebesgue measurable subsets of $(0,\infty)$. The result allows us to express the product measure of sufficiently nice sets $L \times F$ in terms of the Lebesgue measure of $\psi_E(L \times F)$, with which we can translate integration over \mathbb{R}^d to integration over $(0,\infty) \times S$ as intended. The proof makes use of F_{σ} and G_{δ} sets and properties of complete measure spaces.

Lemma 3.1.11. For any $L \in \mathcal{L}$ and $F \in \Sigma_{P,E}$, $\psi_E(L \times F) \in \mathcal{M}_d$ and

$$m(\psi_E(L \times F)) = (\lambda_P \times \sigma_{P,E})(L \times F).$$

Proof. Fix $L \in \mathcal{L}$ and $F \in \Sigma_{P,E}$. It is easy to see that λ_P and the Lebesgue measure dr on $(0,\infty)$ are mutually absolutely continuous. It follows that $((0,\infty),\mathcal{L},\lambda_P)$ is a complete measure space and, further, that there exists an F_σ set $L_\sigma \subseteq (0,\infty)$ and a G_δ set $L_\delta \subseteq (0,\infty)$ for which $L_\sigma \subseteq L \subseteq L_\delta$ and $\lambda_P(L_\delta \setminus L_\sigma) = 0$. Note that, necessarily, $\lambda_P(L) = \lambda_P(L_\sigma) = \lambda_P(L_\delta)$. We have

$$\psi_E(L \times F) = \psi_E(L_\sigma \times F) \cup \psi_E((L \setminus L_\sigma) \times F)$$
(3.8)

where, by virtue of the preceding lemma, $\psi_E(L_\sigma \times F) \in \mathcal{M}_d$ and

$$m(\psi_E(L_\sigma \times F)) = (\lambda_P \times \sigma_{P,E})(L_\sigma \times F) = \lambda_P(L_\sigma)\sigma_{P,E}(F) = \lambda_P(L)\sigma_{P,E}(F) = (\lambda_P \times \sigma_{P,E})(L \times F).$$
(3.9)

Observe that

$$\psi_E((L \setminus L_{\sigma}) \times F) \subseteq \psi_E((L_{\delta} \setminus L_{\sigma}) \times F)$$

where, because $L_{\delta} \setminus L_{\sigma}$ is an G_{δ} set, the latter set is a member of \mathcal{M}_d and

$$m(\psi_E((L_\delta \setminus L_\sigma) \times F)) = (\lambda_P \times \sigma_{P,E})((L_\delta \setminus L_\sigma) \times F) = \lambda_P(L_\delta \setminus L_\sigma)\sigma_{P,E}(F) = 0$$

by virtue of the preceding lemma. Using the fact that $(\mathbb{R}^d \setminus \{0\}, \mathcal{M}_d, m)$ is complete, we conclude that $\psi_E((L \setminus L_\sigma) \times F) \in \mathcal{M}_d$ and $m(\psi_E((L \setminus L_\sigma) \times F)) = 0$. It now follows from (3.8) and (3.9) that $\psi_E(L \times F) \in \mathcal{M}_d$ and

$$m(\psi_E(L \times F)) = m(\psi_E(L_\sigma \times F)) + m(\psi_E((L \setminus L_\sigma) \times F)) = (\lambda_P \times \sigma_{P,E})(L \times F),$$

as desired.
$$\Box$$

The following two lemmas provide another key ingredient to proving Theorem 3.1.8. Lemma 3.1.12 is a standard result in real analysis which states that a compact subset of any metric space contains a countably dense set. This results is then used in Lemma 3.1.13 to show that every open set in $\mathbb{R}^d \setminus \{0\}$ can be written as a countable union of sets of the form $\psi_E(\mathcal{U})$ where \mathcal{U} is an open rectangle in $(0, \infty) \times S$. In our proof of Theorem 3.1.8, we will use this result as justification for the fact that $\psi_E(\mathcal{L} \times \Sigma_{P,E})$ contains every open subset of $\mathbb{R}^d \setminus \{0\}$.

Lemma 3.1.12. Let S be a compact subset of a metric space. Then S contains a countably dense set.

Proof. For each $n \in \mathbb{N}$, consider the open cover

$$\{B_{1/n}(x)\cap S, x\in S\}$$

of S. Since S is compact, there exists a finite subcover. Let $x_{j,n}, j=1,2,\ldots N_n$ denote the center of each of the balls, then, we have that S is covered by $\{B_{1/n}(x_{j,n})\cap S, j=1,2,\ldots,N_n\}$. Thus, for each $n\in\mathbb{N}$, we have a finite set $\{x_{j,n}\}$ of centers. The countable union of these finite sets, $\bigcup_{n=1}^{\infty}\{x_{j,n}\}$, is countable. It is also dense because for every point $x\in S$ and $\epsilon>0$, there is always some n such that $|x_{j,n}-x|<1/n<\epsilon$.

Lemma 3.1.13. Every open subset $U \subseteq \mathbb{R}^d \setminus \{0\}$ can be written as a countable union of open sets of the form $\psi_E(\mathcal{U})$ where $\mathcal{U} = I \times \mathcal{O}$ is an open rectangle in $(0, \infty) \times S$.

Proof. Let $\{r_k\}_{k=1}^{\infty}$ and $\{\eta_j\}_{j=1}^{\infty}$ be countably dense subsets of $(0,\infty)$ and S, respectively. For each triple of natural numbers $j,l,n\in\mathbb{N}_+$, consider the open set

$$\mathcal{U}_{j,l,n} = \{|r - r_j| < 1/n\} \times \mathcal{O}_{l,n} \subseteq (0,\infty) \times S$$

where

$$\mathcal{O}_{l,n} = \{ \eta \in S : |\eta - \eta_l| < 1/n \}.$$

Let $U \subseteq \mathbb{R}^d \setminus \{0\}$ be open. We will show that

$$U = \bigcup_{\substack{j,l,n\\\psi_E(\mathcal{U}_{j,l,n}) \subseteq U}} \psi_E(\mathcal{U}_{j,l,n}), \tag{3.10}$$

where, in view of Proposition 3.1.3, each $\psi_E(\mathcal{U}_{j,l,n})$ is open. It is clear that any element of the union on the right hand side of (3.10) belongs to some $\psi_E(\mathcal{U}_{j,l,n}) \subseteq U$ and so the union is a subset of U. To prove (3.10), it therefore suffices to prove that, for each $x \in U$, there exists a triple j, l, n with

$$x \in \psi_E(\mathcal{U}_{j,l,n}) \subseteq U.$$

To this end, fix $x \in U$ and let $\delta > 0$ be such that $\mathbb{B}_{\delta}(x) \subseteq U$. Consider $(r_x, \eta_x) = \psi_E^{-1}(x) \in (0, \infty) \times S$ and set $M = ||r_x^E|| > 0$ and C = ||E|| > 0. Observe that

$$||I - \alpha^E|| = \left\| \sum_{k=1}^{\infty} \frac{(\ln \alpha)^k}{k!} E^k \right\| \le \sum_{k=1}^{\infty} \frac{|\ln \alpha|^k}{k!} ||E||^k = e^{(C|\ln \alpha|)} - 1$$

for all $\alpha > 0$. Since $\alpha \mapsto e^{(C|\ln \alpha|)} - 1$ is continuous and 0 at $\alpha = 1$, we can choose $\delta' > 0$ for which

$$||I - \alpha^E|| < \frac{\delta}{2M(|\eta_x| + 2)}$$

whenever $|\alpha-1|<\delta'$. Fix an integer $n>\max\{1/\delta'r_x,4M/\delta\}$ and, using the density of the collections $\{r_j\}$ and $\{\eta_l\}$, let r_j and η_l be such that $|r_j-r_x|<1/n$ and $|\eta_l-\eta_x|<1/n$. It follows that the corresponding open set $\mathcal{U}_{j,l,n}$ contains $\psi_E^{-1}(x)$, or, equivalently, $x\in\psi_E(\mathcal{U}_{j,l,n})$. Thus, it remains to show that $\psi_E(\mathcal{U}_{j,l,n})\subseteq\mathbb{B}_\delta(x)$. To this end, let $y=\psi_E(r_y,\eta_y)\in\psi_E(\mathcal{U}_{j,l,n})$ and observe that

$$|x - y| \leq |\psi_E(r_x, \eta_x,) - \psi_E(r_x, \eta_y)| + |\psi_E(r_x, \eta_y) - \psi_E(r_y, \eta_y)|$$

$$= |r_x^E(\eta_x - \eta_y)| + |(r_x^E - r_y^E)\eta_y|$$

$$\leq M|\eta_x - \eta_y| + ||r_x^E - r_y^E|||\eta_y|.$$

Since both $(\eta_x, r_x), (\eta_y, r_y) \in \mathcal{U}_{j,l,n}$, we have

$$|\eta_x - \eta_y| \le |\eta_x - \eta_j| + |\eta_j - \eta_y| < \frac{2}{n}$$
 and $|\eta_y| \le |\eta_y - \eta_x| + |\eta_x| < |\eta_x| + \frac{2}{n}$.

Also, since $|r_x - r_y| < 1/n$, it follows that $r_y = \alpha r_x$ where $|1 - \alpha| < 1/n r_x < \delta'$ by our choice of n. Consequently,

$$|x - y| < \frac{2}{n}M + \left(|\eta_x| + \frac{2}{n}\right) ||r_x^E - r_x^E \alpha^E||$$

$$< \frac{2}{n}M + (|\eta_x| + 2)M||I - \alpha^E||$$

$$< \frac{2}{n}M + \frac{\delta M(|\eta_x| + 2)}{2M(|\eta_x| + 2)}$$

$$< \frac{\delta}{2} + \frac{\delta}{2} = \delta$$

and so we have established (3.10). Finally, upon noting that $\{\mathcal{U}_{j,l,n}\}_{(j,l,n)\in\mathbb{N}_+^3}$ is a countable collection of open rectangles, the union in (3.10) is necessarily countable and we are done with the proof.

In our final lemma preceding the proof of Theorem 3.1.8, we treat a general measure-theoretic statement which gives sufficient conditions concerning two measure spaces to ensure that their completions are isomorphic. This lemma is used in the proof of Theorem 3.1.8 to show that the (complete) measure spaces $((0,\infty), (\mathcal{L} \times \Sigma_{P,E})', \lambda_P \times \sigma_{P,E})$ and $(\mathbb{R}^d \setminus \{0\}, \mathcal{M}_d, m)$ are isomorphic with point isomorphism ψ_E . Here, $((0,\infty), (\mathcal{L} \times \Sigma_{P,E})', \lambda_P \times \sigma_{P,E})$ is the completion of $((0,\infty), \mathcal{L} \times \Sigma_{P,E}, \lambda_P \times \sigma_{P,E})$, and $(\mathbb{R}^d \setminus \{0\}, \mathcal{M}_d, m)$ is the completion of $(\mathbb{R}^d \setminus \{0\}, \mathcal{B}(\mathbb{R}^d \setminus \{0\}), m)$.

Lemma 3.1.14. Let (X_1, Σ_1, ν_1) and (X_2, Σ_2, ν_2) be measure spaces, let $\varphi : X_1 \to X_2$ be a bijection and denote by (X_i, Σ_i', ν_i') the completion of the measure space (X_i, Σ_i, ν_i) for i = 1, 2. Assume that the following two properties are satisfied:

- 1. For each $A_1 \in \Sigma_1$, $\varphi(A_1) \in \Sigma_2'$ and $\nu_2'(\varphi(A_1)) = \nu_1(A_1)$.
- 2. For each $A_2 \in \Sigma_2$, $\varphi^{-1}(A_2) \in \Sigma_1'$ and $\nu_1'(\varphi^{-1}(A_2)) = \nu_2(A_2)$.

Then the measure spaces (X_1, Σ_1', ν_1') and (X_2, Σ_2', ν_2') are isomorphic with point isomorphism φ . Specifically,

$$\Sigma_2' = \{ A_2 \subseteq X_2 : \varphi^{-1}(A_2) \in \Sigma_1' \}$$
(3.11)

and

$$\nu_2'(A_2) = \nu_1'(\varphi^{-1}(A_2)) \tag{3.12}$$

for all $A_2 \in \Sigma_2'$.

Proof. Let us first assume that $A_2 \in \Sigma_2'$. By definition, $A_2 = G_2 \cup H_2$ where $G_2 \in \Sigma_2$ and $H_2 \subseteq G_{2,0} \in \Sigma_2$ with $\nu_2'(A_2) = \nu_2(G_2)$ and $\nu_2'(H_2) = \nu_2(G_{2,0}) = 0$. Consequently, $\varphi^{-1}(A_2) = \varphi^{-1}(G_2) \cup \varphi^{-1}(H_2)$ and $\varphi^{-1}(H_2) \subseteq \varphi^{-1}(G_{2,0})$. In view of Property 2, $\varphi^{-1}(G_2), \varphi^{-1}(G_{2,0}) \in \Sigma_1'$ and we have

$$\nu_1'(\varphi^{-1}(G_2)) = \nu_2(G_2) = \nu_2'(A_2)$$
 and $\nu_1'(\varphi^{-1}(G_{2,0})) = \nu_2(G_{2,0}) = 0.$

In view of the fact that $(X_1', \Sigma_1', \nu_1')$ is complete, $\varphi^{-1}(H_2) \in \Sigma_1'$ and $\nu_1'(\varphi^{-1}(H_2)) = 0$. Consequently, we obtain $\varphi^{-1}(A_2) = \varphi^{-1}(G_2) \cup \varphi^{-1}(H_2) \in \Sigma_1'$ and

$$\nu_2'(A_2) = \nu_1'(\varphi^{-1}(G_2)) \le \nu_1'(\varphi^{-1}(A_2)) \le \nu_1'(\varphi^{-1}(G_2)) + \nu_1'(\varphi^{-1}(H_2)) = \nu_2(G_2) + 0 = \nu_2'(A_2).$$

From this we obtain that $\Sigma_2' \subseteq \{A_2 \subseteq X_2 : \varphi^{-1}(A_2) \in \Sigma_1'\}$ and, for each $A_2 \in \Sigma_2'$, $\nu_2'(A_2) = \nu_1'(\varphi^{-1}(A_2))$. It remains to prove that

$${A_2 \subseteq X_2 : \varphi^{-1}(A_2) \in \Sigma_1'} \subseteq \Sigma_2'.$$

To this end, let A_2 be a subset of X_2 for which $\varphi^{-1}(A_2) \in \Sigma_1'$. By the definition of Σ_1' , we have $\varphi^{-1}(A_2) = G_1 \cup H_1$ where $G_1 \in \Sigma_1$, $H_1 \subseteq G_{1,0} \in \Sigma_1$ and $\nu_1'(H_1) = \nu_1(G_{1,0}) = 0$. In view of Property 1, $\varphi(G_1) \in \Sigma_2'$, $\varphi(H_1) \subseteq \varphi(G_{1,0}) \in \Sigma_2'$ and $\nu_2'(\varphi(G_{1,0})) = \nu_1(G_{1,0}) = 0$. Because $(X_2', \Sigma_2', \nu_2')$ is complete, we have $\varphi(H_1) \in \Sigma_2'$ and so

$$A_1 = \varphi(\varphi^{-1}(A_2)) = \varphi(G_1) \cup \varphi(H_1) \in \Sigma_2',$$

as desired. \Box

We are finally in a position to prove Theorem 3.1.8. Our approach can be described as follows. We first argue that $\mathcal{L} \times \Sigma_{P,E}$ is contained in the monotone class \mathcal{C} comprised of the sets $G \subseteq (0,\infty) \times S$ for which $\psi_E(G) \in \mathcal{M}_d$ and $m(\psi_E(G)) = (\lambda_P \times \sigma_{P,E})(G)$ via the monotone class lemma. This allows us to show that every set G in $\mathcal{L} \times \Sigma_{P,E}$ satisfies the two properties above, and that the σ -algebra of Borel subsets of $\mathbb{R}^d \setminus \{0\}$ is contained in $\psi_E(\mathcal{L} \times \Sigma_{P,E})$. Property 1 follows from recognizing that the σ -algebras and measure spaces in the hypothesis of Theorem 3.1.8 are completions of those in preceding arguments. To prove Property 2, we let a measureable function $f: \mathbb{R}^d \to \mathbb{C}$ be given and use Property 2, the monotone convergence theorem (in case that $f \geq 0$), and Fubini's theorem.

Proof of Theorem 3.1.8. Denote by $\mathcal C$ the collection of sets $G\subseteq (0,\infty)\times S$ for which $\psi_E(G)\in \mathcal M_d$ and $m(\psi_E(G))=(\lambda_P\times\sigma_{P,E})(G)$. By virtue of Lemma 3.1.11, it follows that $\mathcal C$ contains all elementary sets, i.e., finite unions of disjoint measurable rectangles. Using the continuity of measure (applied to the measures m and $\lambda_P\times\sigma_{P,E}$) and the fact that ψ_E is a bijection, it is straightforward to verify that $\mathcal C$ is a monotone class. By the monotone class lemma (Theorem 8.3 of [14]), it immediately follows that $\mathcal L\times\Sigma_{P,E}\subseteq\mathcal C$. In other words, for each $G\in\mathcal L\times\Sigma_{P,E}$,

$$\psi_E(G) \in \mathcal{M}_d$$
 and $m(\psi_E(G)) = (\lambda_P \times \sigma_{P,E})(G).$ (3.13)

We claim that, for each Borel subset A of $\mathbb{R}^d \setminus \{0\}$, $\psi_E^{-1}(A) \subseteq \mathcal{L} \times \Sigma_{P,E}$. To this end, we write

$$\psi_E(\mathcal{L} \times \Sigma_{P,E}) = \{ \psi_E(G) : G \in \mathcal{L} \times \Sigma_{P,E} \}$$

for the σ -algebra on $\mathbb{R}^d \setminus \{0\}$ induced by ψ_E . In view of Lemma 3.1.13, $\psi_E(\mathcal{L} \times \Sigma_{P,E})$ contains every open subset of $\mathbb{R}^d \setminus \{0\}$ and therefore

$$\mathcal{B}(\mathbb{R}^d \setminus \{0\}) \subseteq \psi_E(\mathcal{L} \times \Sigma_{P,E}).$$

where $\mathcal{B}(\mathbb{R}^d \setminus \{0\})$ denotes the σ -algebra of Borel subsets of $\mathbb{R}^d \setminus \{0\}$ thus proving our claim.

Together, the results of the two preceding paragraphs show that, for each $A \in \mathcal{B}(\mathbb{R}^d \setminus \{0\})$, $\psi_E^{-1}(A) \subseteq \mathcal{L} \times \Sigma_{P,E}$ and $m(A) = (\lambda_P \times \sigma_{P,E})(\psi_E^{-1}(A))$. Upon noting that $\mathcal{L} \times \Sigma_{P,E} \subseteq (\mathcal{L} \times \Sigma_{P,E})'$, we immediately obtain the following statement: For each $A \in \mathcal{B}(\mathbb{R}^d \setminus \{0\})$,

$$\psi_E^{-1}(A) \in (\mathcal{L} \times \Sigma_{P,E})'$$
 and $m(A) = (\lambda_P \times \sigma_{P,E})(\psi_E^{-1}(A)).$ (3.14)

In comparing (3.13) and (3.14) with Properties 1 and 2 of Lemma 3.1.14 and, upon noting that $((0,\infty)\times S, (\mathcal{L}\times\Sigma_{P,E})', \lambda_P\times\sigma_{P,E})$ is the completion of $((0,\infty)\times S, \mathcal{L}\times\Sigma_{P,E}, \lambda_P\times\sigma_{P,E})$ and $(\mathbb{R}^d\setminus\{0\},\mathcal{M}_d,m)$ is the completion of $(\mathbb{R}^d\setminus\{0\},\mathcal{B}(\mathbb{R}^d\setminus\{0\}),m)$, Property 1 of Theorem 3.1.8 follows immediately from Lemma 3.1.14.

It remains to prove Property 2. To this end, let $f: \mathbb{R}^d \to \mathbb{C}$ be Lebesgue measurable. Because $\mathcal{M}_d = \{A \subseteq \mathbb{R}^d \setminus \{0\} : \psi_E^{-1}(A) \in (\mathcal{L} \times \Sigma_{P,E})'\}$, it follows that $f \circ \psi_E$ is $(\mathcal{L} \times \Sigma_{P,E})'$ -measurable. In the case that $f \geq 0$, we may approximate f monotonically by simple functions and, by invoking Property 1 and the monotone convergence theorem, we find that

$$\int_{\mathbb{R}^d} f(x) \, dx = \int_{\mathbb{R}^d \setminus \{0\}} f(x) \, dx = \int_{(0,\infty) \times S} f \circ \psi_E \, d(\lambda_P \times \sigma_{P,E}). \tag{3.15}$$

From this, (3.6) follows from Fubini's theorem (see, e.g., Part (a) of Theorem 8.8 and Theorem 8.12 of [14]). Finally, by applying the above result to $|f| \ge 0$, we obtain $f \in L^1(\mathbb{R}^d)$ if and only if $f \circ \psi_E \in L^1((0,\infty) \times S, (\mathcal{L} \times \Sigma_{P,E})', \lambda_P \times \sigma_{P,E})$. In this case, by applying (3.15) to $\operatorname{Re}(f)_+, \operatorname{Re}(f)_-, \operatorname{Im}(f)_+$ and $\operatorname{Im}(f)_-$, we find that (3.15) holds for our integrable f and, by virtue Fubini's theorem (see, e.g., Part (b) of Theorem 8.8 and Theorem 8.12 of [14]), the desired result follows.

Our next result, Proposition 3.1.15, guarantees that, in particular, $\sigma_{P,E}$ is a Radon measure.

Proposition 3.1.15. *The following statements are true:*

- 1. $(S, \Sigma_{P,E}, \sigma_{P,E})$ is the completion of the measure space $(S, \mathcal{B}(S), \sigma_{P,E})$. In particular, the $(S, \Sigma_{P,E}, \sigma_{P,E})$ is complete and every $F \in \Sigma_{P,E}$ is of the form $F = G \cup H$ where G is a Borel set and H is a subset of a Borel set Z with $\sigma_{P,E}(Z) = 0$.
- 2. For each $F \in \Sigma_{P,E}$,

$$\sigma_{P,E}(F) = \inf\{\sigma_{P,E}(\mathcal{O}) : F \subseteq \mathcal{O} \subseteq S \text{ and } \mathcal{O} \text{ is open}\}$$
 (3.16)

and

$$\sigma_{P,E}(F) = \sup\{\sigma_{P,E}(K) : K \subseteq F \subseteq S \text{ and } K \text{ is compact}\}.$$
 (3.17)

Remark 2. This proposition can be seen as an application of Proposition 3.1.4 and Theorem 2.18 of [14]. The proof we give here is distinct and, we believe, nicely illustrates the utility of (3.6) of Theorem 3.1.8.

Proof. Throughout the proof, we shall write $\sigma = \sigma_{P,E}$, $\Sigma = \Sigma_{P,E}$ and, for each $F \subseteq S$, $\widetilde{F} = \widetilde{F_E}$. We remark that, by standard arguments using G_δ and F_σ sets, Item 1 follows immediately from Item 2. Also, given that S is compact and σ is finite, it suffices to prove (3.16), i.e., it suffices to prove the statement: For each $F \in \Sigma$ and $\epsilon > 0$, there is an open subset \mathcal{O} of S containing F for which

$$\sigma(\mathcal{O}\setminus F)<\epsilon.$$

To this end, let $F \in \Sigma$ and $\epsilon > 0$. Given that $\widetilde{F} \in \mathcal{M}_d$ and m is outer regular, there exists an open set $U \subseteq \mathbb{R}^d \setminus \{0\}$ for which $\widetilde{F} \subseteq U$ and $m(U \setminus \widetilde{F}) < \epsilon/(2\mu_P)$. Since \widetilde{F} is a subset of the open set $B \setminus \{0\}$, we may assume without loss of generality that $U \subseteq B \setminus \{0\}$ and so $m(\widetilde{F}) \leq m(U) < \infty$ and

$$m(U \setminus \widetilde{F}) = m(U) - m(\widetilde{F}) < \epsilon/(2\mu_P).$$
 (3.18)

For each 0 < r < 1, consider the open set

$$\mathcal{O}_r = S \cap \left(r^{-E}U\right)$$

in S. Observe that, for each $x \in F$, $r^E x \in \widetilde{F} \subseteq U$ and therefore $x \in \mathcal{O}_r$. Hence, for each 0 < r < 1, \mathcal{O}_r is an open subset of S containing F.

We claim that there is at least one $r_0 \in (0,1)$ for which

$$m(\widetilde{\mathcal{O}_{r_0}}) < m(U) + \epsilon/(2\mu_P).$$
 (3.19)

To prove the claim, we shall assume, to reach a contradiction, that

$$m(\widetilde{\mathcal{O}_r}) \ge m(U) + \epsilon/(2\mu_P)$$

for all 0 < r < 1. By virtue of (3.6) of Theorem 3.1.8,

$$m(U) = \int_0^\infty \left(\int_S \chi_U(r^E \eta) \, \sigma(d\eta) \right) r^{\mu_P - 1} \, dr.$$

Upon noting that $U \subseteq B \setminus \{0\}$, it is easy to see that

$$U = \bigcup_{0 < s < 1} s^E \mathcal{O}_s \qquad \text{ and } \qquad r^E \eta \in \bigcup_{0 < s < 1} s^E \mathcal{O}_s$$

if and only if 0 < r < 1 and $\eta \in \mathcal{O}_r$. Consequently,

$$m(U) = \int_0^1 \left(\int_S \chi_{\mathcal{O}_r}(\eta) \, \sigma(d\eta) \right) \, r^{\mu_P - 1} \, dr = \int_0^1 \sigma(\mathcal{O}_r) r^{\mu_P - 1} \, dr = \int_0^1 \mu_P \cdot m(\widetilde{\mathcal{O}_r}) \, r^{\mu_P - 1} \, dr.$$

Upon making use of our supposition, we have

$$\int_0^1 \mu_P \cdot m(\widetilde{\mathcal{O}_r}) r^{\mu_P - 1} \, dr \ge \int_0^1 \mu_P \cdot (m(U) + \epsilon/(2\mu_P)) r^{\mu_P - 1} \, dr = m(U) + \epsilon/(2\mu_P)$$

and so

$$m(U) \ge m(U) + \epsilon/(2\mu_P),$$

which is impossible. Thus, the stated claim is true.

Given any such r_0 for which (3.19) holds, set $\mathcal{O} = \mathcal{O}_{r_0}$. As previously noted, \mathcal{O} is an open subset of S which contains F. In view of (3.18) and (3.19), we have

$$m(\widetilde{\mathcal{O}}) - m(\widetilde{F}) < m(U) - m(\widetilde{F}) + \epsilon/(2\mu_P) < \epsilon/(2\mu_P) + \epsilon/(2\mu_P) = \epsilon/\mu_P$$

and therefore

$$\sigma(\mathcal{O} \setminus F) = \sigma(\mathcal{O}) - \sigma(F) = \mu_P(m(\widetilde{\mathcal{O}}) - m(\widetilde{F})) < \epsilon,$$

as desired.

3.1.3 The construction is independent of $E \in \text{Exp } P$

In this subsection, we show that the measure $\sigma_{P,E}$ is independent of the choice of $E \in \text{Exp}(P)$ and complete the proof of Theorem 3.1.1. We let $E_1, E_2 \in \text{Exp}(P)$ and consider the measure spaces $(S, \Sigma_{P,E_1}, \sigma_{P,E_1})$ and $(S, \Sigma_{P,E_2}, \sigma_{P,E_2})$ produced via the construction in Subsection 3.1.1.

Proposition 3.1.16. $\Sigma_{P,E_1} = \Sigma_{P,E_2}$ and $\sigma_{P,E_1} = \sigma_{P,E_2}$.

Proof. Throughout the proof, we will write $\Sigma_i = \Sigma_{P,E_i}$ and $\sigma_i = \sigma_{P,E_i}$ for i = 1, 2. In view of the Proposition 3.1.15, it suffices to show that

$$\sigma_1(F) = \sigma_2(F)$$

for all $F \in \mathcal{B}(S) \subseteq \Sigma_1 \cap \Sigma_2$. To this end, we let $F \in \mathcal{B}(S)$ be arbitrary but fixed.

Given $n \in \mathbb{N}_+$, using the regularity of the measures σ_1 and σ_2 , select open sets $\mathcal{O}_{n,1}, \mathcal{O}_{n,2}$ and compact sets $K_{n,1}, K_{n,2}$ for which

$$K_{n,j} \subseteq F \subseteq \mathcal{O}_{n,j}$$
 and $\sigma_j(\mathcal{O}_{n,j} \setminus K_{n,j}) < 1/n$

for j=1,2. Observe that $K_n=K_{n,1}\cup K_{n,2}$ is a compact set, $\mathcal{O}_n=\mathcal{O}_{n,1}\cap \mathcal{O}_{n,2}$ is an open set, and $K_n\subseteq F\subseteq \mathcal{O}_n$. Furthermore,

$$\sigma_j(\mathcal{O}_n \setminus K_n) \le \sigma_j(\mathcal{O}_{n,j} \setminus K_{n,j}) < 1/n$$

for j=1,2. Given that \mathcal{O}_n is open in S, $\mathcal{O}_n=S\cup U_n$ where U_n is an open subset of \mathbb{R}^d and, because that S is compact, $K_n=K_n\cap S$ is a compact subset of \mathbb{R}^d . By virtue of Urysohn's lemma, let $\phi_n:\mathbb{R}^d\to[0,1]$ be a continuous function which is compactly supported in U_n and for which $\phi_n(x)=1$ for all $x\in K_n$. Using this sequence of functions $\{\phi_n\}$, we establish the following useful lemma.

Lemma 3.1.17. For j = 1, 2 and $n \in \mathbb{N}$, define $g_{n,j} : (0, \infty) \to \mathbb{R}$ by

$$g_{n,j}(r) = \int_{S} \phi_n(r^{E_j}\eta) \, \sigma_j(d\eta).$$

for r > 0. Then $g_{n,j}$ is continuous for each $n \in \mathbb{N}$ and j = 1, 2 and

$$\sigma_j(F) = \lim_{n \to \infty} g_{n,j}(1)$$

for j = 1, 2.

Proof. First, we note that, for each $r \in (0,\infty)$, the above integral makes sense because $\eta \mapsto \phi_n(r^{E_j}\eta)$ is Borel measurable (because it's continuous on S) and non-negative. Let $\epsilon>0$ and $r_0\in (0,\infty)$ be arbitrary but fixed. It is clear that the function $(0,\infty)\times S\ni (r,\eta)\mapsto \phi_n(r^{E_j}\eta)$ is continuous on its domain and therefore, in view of the compactness of S, we can find a $\delta>0$ for which

$$|\phi_n(r^{E_j}\eta) - \phi_n(r_0^{E_j}\eta)| \le \frac{\epsilon}{2\sigma_j(S)}$$
 whenever $|r - r_0| < \delta$

for all $\eta \in S$. The triangle inequality guarantees that

$$|g_{n,j}(r) - g_{n,j}(r_0)| \le \int_S |\phi_n(r^{E_j}\eta) - \phi_n(r_0^{E_j}\eta)| \, \sigma_j(d\eta) \le \epsilon/2 < \epsilon$$

whenever $|r - r_0| < \delta$. Thus, $g_{n,j}$ is continuous.

We observe that

$$g_{n,j}(1) = \int_{S} \phi_n(\eta) \, \sigma_j(d\eta)$$

because $1^{E_j} = I$. By construction, we have $\chi_{K_n}(\eta) \le \phi_n(\eta) \le \chi_{\mathcal{O}_n}(\eta)$ for all $\eta \in S$ and $n \in \mathbb{N}_+$ and therefore

$$\sigma_j(K_n) \le g_{n,j}(1) \le \sigma_j(\mathcal{O}_n)$$

by the monotonicity of the integral. Since

$$\sigma_j(F) = \lim_{n \to \infty} \sigma_j(K_n) = \lim_{n \to \infty} \sigma_j(\mathcal{O}_n)$$

in view of our choice of \mathcal{O}_n and K_n , the remaining result follows immediately from the preceding inequality (and the squeeze theorem).

Let us now complete the proof of Proposition 3.1.16. Given any 0 < s < 1 < t and $n \in \mathbb{N}$, consider the function $f = f_{n,s,t} : \mathbb{R}^d \to [0,1]$ given by

$$f(x) = \phi_n(x)\chi_{[s,t]}(P(x))$$

for $x \in \mathbb{R}^d$. It is clear that f is Lebesgue measurable on \mathbb{R}^d and non-negative. By virtue of Theorem 3.1.8 (applied to the two measures σ_1 and σ_2), we have

$$\int_0^\infty \int_S f(r^{E_1}\eta) \,\sigma_1(d\eta) r^{\mu_P - 1} \,dr = \int_{\mathbb{R}^d} f(x) \,dx = \int_0^\infty \int_S f(r^{E_2}\eta) \,\sigma_2(d\eta) r^{\mu_P - 1} \,dr \tag{3.20}$$

Upon noting that

$$f(r^{E_j}\eta) = \phi_n(r^{E_j}\eta)\chi_{[s,t]}\left((P(r^{E_j}\eta)\right) = \phi_n(r^{E_j}\eta)\chi_{[s,t]}(rP(\eta)) = \chi_{[s,t]}(r)\phi_n(r^{E_j}\eta)$$

for $r \in (0, \infty)$, $\eta \in S$, and j = 1, 2, we have

$$\int_0^\infty \int_S f(r^{E_j} \eta) \, \sigma_j(d\eta) r^{\mu_P - 1} \, dr = \int_{[s,t]} \int_S \phi_n(r^{E_j} \eta) \, \sigma_j(d\eta) r^{\mu_P - 1} \, dr = \int_{[s,t]} g_{n,j}(r) r^{\mu_P - 1} \, dr$$

for j=1,2. By virtue of the Lemma 3.1.17, $r\mapsto g_{n,j}(r)r^{\mu_P-1}$ is continuous and necessarily bounded on [s,t] and so the final integral above can be interpreted as a Riemann integral. In this interpretation, we have

$$\int_{s}^{t} g_{n,j}(r) r^{\mu_P - 1} dr = \int_{0}^{\infty} \int_{S} f(r^{E_j} \eta) \,\sigma_j(d\eta) r^{\mu_P - 1} dr \tag{3.21}$$

for j = 1, 2 and 0 < s < 1 < t. In view of (3.20) and (3.21), we conclude that

$$\int_{s}^{t} g_{n,1}(r)r^{\mu_{P}-1} dr = \int_{s}^{t} g_{n,2}(r)r^{\mu_{P}-1} dr$$

for all 0 < s < 1 < t. In view of continuity of the integrands, an application of the fundamental theorem of calculus now guarantees that $g_{n,1}(1) = g_{n,2}(1)$ for each $n \in \mathbb{N}$. Therefore

$$\sigma_1(F) = \lim_{n \to \infty} g_{n,1}(1) = \lim_{n \to \infty} g_{n,2}(1) = \sigma_2(F)$$

by virtue of Lemma 3.1.17.

In view of Proposition 3.1.16, we will denote by Σ_P and σ_P the unique σ -algebra and measure on S which, respectively, satisfy

$$\Sigma_P = \Sigma_{P,E}$$
 and $\sigma_P = \sigma_{P,E}$

for all $E \in \text{Exp}(P)$. We will henceforth assume this notation.

Proposition 3.1.18. For any $O \in \operatorname{Sym}(P)$ and $F \in \Sigma_P$, $OF \in \Sigma_P$ and

$$\sigma_P(OF) = \sigma_P(F).$$

That is, the measure σ_P is invariant under the action by Sym(P).

Proof. Let $O \in \text{Sym}(P)$, $F \in \Sigma_P$ and, for $E \in \text{Exp}(P)$, define $E' = O^*EO$. In view of Proposition 2.0.10, we note that $E' \in \text{Exp}(P)$. Observe that

$$\widetilde{(OF)_E} = \bigcup_{0 < r < 1} r^E(OF) = \bigcup_{0 < r < 1} O\left(O^*r^EOF\right) = O\left(\bigcup_{0 < r < 1} r^{E'}F\right) = O\widetilde{F_{E'}}$$
(3.22)

thanks to Proposition 6.1.1. In view of Proposition 3.1.16, we have $OF_{E'} \in \mathcal{M}_d$ because $F \in \Sigma_P = \Sigma_{P,E'}$ and O is linear. Using (3.22), we find that $OF_E \in \mathcal{M}_d$ and therefore $OF \in \Sigma_{P,E} = \Sigma_P$ by virtue of Proposition 3.1.16. In view of the fact that O is orthogonal,

$$\sigma_{P\!,E}(OF) = \mu_P \cdot m\left(\widetilde{(OF)_E}\right) = \mu_P \cdot m\left(\widetilde{OF_{E'}}\right) = \mu_P \cdot m\left(\widetilde{F}_{E'}\right) = \sigma_{P\!,E'}(F)$$

and therefore, a final appeal to Proposition 3.1.16 guarantees that

$$\sigma_P(OF) = \sigma_{P,E}(OF) = \sigma_{P,E'}(F) = \sigma_P(F),$$

as desired. \Box

This is the final result required to completely prove Theorem 3.1.1:

Proof of Theorem 3.1.1. Together, the results of Propositions 3.1.15, 3.1.16 and 3.1.18, guarantee that σ_P is a Radon measure satisfying Properties 1 and 2. Property 3 follows directly from Proposition 3.1.16 and the definition of σ_P in terms of $\sigma_{P,E}$ for any $E \in \operatorname{Exp}(P)$. Similarly, Properties 1 and 2 follow from Theorem 3.1.8 by virtue of Proposition 3.1.16.

3.2 Using a smooth structure on S to compute σ

In this section, we consider the special case in which a positive-homogeneous function P on \mathbb{R}^d is smooth¹². Under this additional assumption, we show in Proposition 3.2.1 that ∇P is everywhere non-vanishing on S and so S is a smooth compact embedded hypersurface in \mathbb{R}^d . For a detailed treatment of smooth manifold theory, the author encourages the reader to refer to [9]. For an elementary application of theory to be presented here, the reader may refer to Section ??.

We first set up some notation: For a smooth manifold M, we denote by $\mathcal{A}(M)$ its unique maximal atlas. Also on M, the collection of smooth vector fields is denoted by $\mathfrak{X}(M)$ and, for each $k=1,2,\ldots$, the set of (smooth) differential k-forms on M will be denoted by $\Omega^k(M)$. In this section, we integrate non-smooth differential forms and, for the generality needed here, we shall refer the reader to [11] for background (Another perspective is given in [1]). To this end, let us denote the Lebesgue σ -algebra of measurable sets on M by $\mathcal{L}(M)$. We note that $F \in \mathcal{L}(M)$ if and only if,

$$\varphi(F \cap \mathcal{U}) \in \mathcal{M}_d$$

for every chart $(\varphi, \mathcal{U}) \in \mathcal{A}(M)$. An n-form ω on M, with $n = \dim(M)$, is said to be (Lebesgue) measurable, if in each coordinate system (φ, \mathcal{U}) , the local representation

$$\omega = h_{\varphi}(x)dx^1 \wedge dx^2 \wedge \dots \wedge dx^n$$

in the coordinates $\varphi=(x^1,x^2,\ldots,x^n)$ has $h_{\varphi}(x)$ a Lebesgue measurable function on $U=\varphi(\mathcal{U})\subseteq\mathbb{R}^n$. The collection of measurable n-forms on M is denoted by $\mathcal{L}(\Lambda^d(M))$ and, naturally, $\Omega^d(M)\subseteq\mathcal{L}(\Lambda^d(M))$. As standard, we shall use Einstein's summation convention.

We view \mathbb{R}^d as smooth oriented Riemannian manifold with its standard Euclidean metric \overline{g} , oriented smooth atlas $\mathcal{A}_+(\mathbb{R}^d)$, and Riemannian volume form $d\operatorname{Vol}_{\mathbb{R}^d}$. Given any $E\in\operatorname{Exp}(P)$, consider $\mathcal{E}_E\in\mathfrak{X}(\mathbb{R}^d)$ defined, at each $x\in\mathbb{R}^d$, by

$$(\mathcal{E}_E)_x(f) = \frac{d}{dt} f(x + t(Ex)) \big|_{t=0}$$

for $f \in C^{\infty}(\mathbb{R}^d)$. In the standard (global) chart with coordinates $x = (x^1, x^2, \dots, x^d)$, $(\mathcal{E}_E)_x \in T_x(\mathbb{R}^d)$ is given by

$$(\mathcal{E}_E)_x = (Ex)^{\alpha} \partial_{x^{\alpha}} = E^{\alpha}_{\beta} x^{\beta} \partial_{x^{\alpha}}$$

where (E_{α}^{β}) is the standard matrix representation for E and $\partial_{x^{\alpha}} = \partial/\partial x^{\alpha}$. By an abuse of notation, we shall write ∇P to denote both the function

$$\mathbb{R}^d \ni x \mapsto \nabla P(x) = \left(\frac{\partial P}{\partial x^1}, \frac{\partial P}{\partial x^2}, \dots, \frac{\partial P}{\partial x^d}\right) \in \mathbb{R}^d,$$

where $\frac{\partial P}{\partial x^{\alpha}} = \frac{\partial P}{\partial x^{\alpha}}|_{x}$ for $\alpha = 1, 2, ..., d$, and its canonical identification $\nabla P \in \mathfrak{X}(\mathbb{R}^{d})$ given by

$$\nabla P_x = \overline{g}^{\alpha\beta} \frac{\partial P}{\partial x^{\alpha}} \partial_{x^{\beta}} = \delta^{\alpha\beta} \frac{\partial P}{\partial x^{\alpha}} \partial_{x^{\beta}} = \sum_{i=1}^d \frac{\partial P}{\partial x^{\alpha}} \partial_{x^{\alpha}}$$

¹Many of the results in this section remain valid (with appropriate modification) under the weaker assumption that $P \in C^k(\mathbb{R}^d)$ for $k = 1, 2, \ldots$ In this setting, S is easily seen to be a C^k manifold. Because working in the smooth category is sufficient for our purposes, we shall not pursue the greater level of generality but invite the reader to do so. ²To avoid trivialities, we assume that d > 1 throughout.

in standard Euclidean coordinates $x=(x^{\alpha})$. Of course, for each $x\in\mathbb{R}^d$, the Riemannian norm $|\nabla P_x|_{\overline{g}}$ of $\nabla P_x\in T_x(\mathbb{R}^d)$ coincides with the Euclidean norm $|\nabla P(x)|$ of $\nabla P(x)\in\mathbb{R}^d$. These equivalent quantities (functions) will be henceforth denoted by $|\nabla P|$.

Proposition 3.2.1. *For each* $\eta \in S$,

$$\overline{g}(\nabla P, \mathcal{E})_{\eta} = \nabla P(\eta) \cdot (E\eta) = 1.$$

In particular, $\nabla P \neq 0$ on S and so S is a compact embedded hypersurface in \mathbb{R}^d .

Proof. Given that $E \in \text{Exp}(P)$ and $P \in C^{\infty}(\mathbb{R}^d)$, we differentiate the identity $rP(x) = P(r^Ex)$ to find that

$$P(x) = \frac{d}{dr}P(r^{E}x) = \nabla P(r^{E}x) \cdot ((r^{E-I}E)x)$$

for r > 0 and $x \in \mathbb{R}^d$. In particular, when r = 1 and $x = \eta \in S$, we have

$$1 = \nabla P(\eta) \cdot (E\eta) = \frac{\partial P}{\partial x^{\alpha}} (E\eta)^{\alpha} = \overline{g}_{\alpha\beta} (\nabla P_{\eta})^{\alpha} (\mathcal{E}_{\eta})^{\beta} = \overline{g} (\nabla P, \mathcal{E})_{\eta}.$$

Thus, our (necessarily) compact level set S is a smooth embedded hypersurface of \mathbb{R}^d in view of the regular level set theorem.

We shall denote by $\iota: S \hookrightarrow \mathbb{R}^d$ the canonical inclusion map and set d' = d - 1. As an embedded submanifold of \mathbb{R}^d , S is a Riemannian submanifold of \mathbb{R}^d with metric g^S given by

$$g^S(X,Y) = \overline{g}(\iota_*(X),\iota_*(Y))$$

for $X,Y\in\mathfrak{X}(S)$; here, for each $\eta\in S$, $\iota_*:T_\eta(S)\to T_\eta(\mathbb{R}^d)$ is the pushforward of ι . In view of the preceding proposition, $N:=\mathbf{\nabla}P/|\mathbf{\nabla}P|\in\mathfrak{X}(\mathbb{R}^d)$ is a smooth unit normal vector field along S and it determines an orientation on the Riemannian manifold S. Equipped with this orientation, (S,g^S) is an oriented Riemannian manifold and we shall denote by $d\operatorname{Vol}_S$ the Riemannian volume form and by $\mathcal{A}_+(S)$ its corresponding (maximal) oriented atlas. By virtue of Proposition 15.21 of [9], $d\operatorname{Vol}_S=(N\sqcup d\operatorname{Vol}_{\mathbb{R}^d})|_{S}$, i.e.,

$$d \operatorname{Vol}_{S}(X_{1}, X_{2}, \dots, X_{d'}) = d \operatorname{Vol}_{\mathbb{R}^{d}}(N, \iota_{*}(X_{1}), \iota_{*}(X_{2}), \dots, \iota_{*}(X_{d'}))$$

$$= \frac{1}{|\nabla P|} d \operatorname{Vol}_{\mathbb{R}^{d}}(\nabla P, \iota_{*}(X_{1}), \iota_{*}(X_{2}), \dots, \iota_{*}(X_{d'}))$$

for any collection $\{X_1, X_2, \dots, X_{d'}\} \in \mathfrak{X}(S)$. Beyond $d\operatorname{Vol}_S \in \Omega^{d'}(S)$, we consider the following smooth d'-form(s): Given $E \in \operatorname{Exp}(P)$, define $d\sigma_{P,E} \in \Omega^{d'}(S)$ by

$$d\sigma_{P,E}(X_1, X_2, \dots, X_{d'}) = d \operatorname{Vol}_{\mathbb{R}^d}(\mathcal{E}_E, \iota_*(X_1), \iota_*(X_2), \dots, \iota_*(X_{d'}))$$

for $X_1, X_2, \dots, X_{d'} \in \mathfrak{X}(S)$. We have

Proposition 3.2.2. *For any* $E \in \text{Exp}(P)$ *,*

$$d\sigma_{P,E} = \frac{1}{|\nabla P|} d\operatorname{Vol}_S.$$

In particular, $d\sigma_{P,E}$ *is positively oriented and is independent of* $E \in \text{Exp}(P)$ *.*

Before proving the proposition, we first treat a lemma of a purely linear algebraic nature.

Lemma 3.2.3. Let $v_1, v_2, \ldots, v_{d'}$ be linearly independent vectors in \mathbb{R}^d and suppose that $w \in \mathbb{R}^d \setminus \{0\}$ is such that $w \perp v_i$ for all i. Then, for any $z \in \mathbb{R}^d$ for which $z \cdot w = 1$,

$$\det(z, v_1, v_2, \dots, v_{d'}) = \frac{1}{|w|} \det(n, v_1, v_2, \dots, v_{d'}) = \frac{1}{|w|^2} \det(w, v_1, v_2, \dots, v_{d'}).$$

where n := w/|w|.

Proof. Given $z \in \mathbb{R}^d$ such that $z \cdot w = 1$, it follows that

$$z = \frac{1}{|w|}n + a_1v_1 + a_2v_2 + \cdots + a_{d'}v_{d'}.$$

By the multilinearity of the determinant map, we have

$$\det(z, v_1, v_2, \dots, v_{d'}) = \det\left(\frac{1}{|w|}n + a_1v_2 + a_2v_2 + \dots + a_{d'}v_{d'}, v_1, v_2, \dots, v_{d'}\right)
= \frac{1}{|w|}\det(n, v_1, v_2, \dots, v_{d'}) + \det(a_1v_1 + \dots + a_{d'}v_{d'}, v_1, v_2, \dots, v_{d'})
= \frac{1}{|w|}\det(n, v_1, v_2, \dots, v_{d'}) + 0$$

where we have used the fact that the columns of the matrix $(a_1v_1 + \cdots + a_{d'}v_{d'}, v_1, v_2, \dots, v_{d'})$ are linearly dependent to conclude that the final determinant is zero.

Proof of Proposition 3.2.2. We fix $E \in \text{Exp}(P)$ and note that the assertion at hand is a local one. Thus, it suffices to verify that, for any $\eta \in S$ and $X_1, X_2, \dots, X_{d'} \in T_\eta(S)$,

$$d\sigma_{P,E}(X_1, X_2, \dots, X_{d'}) = \frac{1}{|\nabla P_{\eta}|} d\operatorname{Vol}_{\mathbb{R}^d}(N_{\eta}, \iota_*(X_1), \iota_*(X_2), \dots, \iota_*(X_{d'})).$$

Fix $\eta \in S$ and let (\mathcal{O}, φ) be a coordinate chart centered at η with local coordinates $u = (u^{\alpha})$. As usual, denote by $x = (x^{\alpha})$ the Euclidean coordinates on \mathbb{R}^d . For each $i = 1, 2, \dots, d'$,

$$X_i = X_i^{lpha} \partial_{u^{lpha}} \qquad ext{ and } \qquad \iota_*(X_i) = v_i^{eta} \partial_{x^{eta}}$$

where

$$v_i^{\beta} = X_i^{\alpha} \frac{\partial x^{\beta}}{\partial u^{\alpha}}.$$

For each $i=1,2,\ldots,d'$, we set $v_i=(v_i^1,v_i^2,\ldots,v_i^{d'})\in\mathbb{R}^d$. Also, let $w=\nabla P(\eta)\in\mathbb{R}^d$ with $|w|=|\nabla P(\eta)|=|\nabla P_\eta|$, set $n=w/|w|\in\mathbb{R}^d$ and note that

$$N_{\eta} = \frac{1}{|\nabla P_{\eta}|} \sum_{k=1}^{d} \frac{\partial P}{\partial x^{k}} \partial_{x^{k}} = n^{\mu} \partial_{x^{\mu}}.$$

Given that ∇P is normal to S, we have

$$v_i \cdot w = \overline{g}_{\mu,\nu} v_i^{\mu} w^{\nu} = \overline{g}(\iota_*(X_i), \nabla P)_{\eta} = 0$$

and therefore $w \perp v_i$ for each $i=1,2,\ldots,d'$. Upon recalling that $(\mathcal{E}_E)_{\eta}=(E\eta)^{\alpha}\partial_{x^{\alpha}}$, set $z=((E\eta)^1,(E\eta)^2,\ldots,(E\eta)^d)\in\mathbb{R}^d$ and observe that $z\cdot w=1$ by virtue of Proposition 3.2.1. An appeal to the lemma guarantees that

$$d\sigma_{P,E}(X_{1}, X_{2}, \dots, X_{d'}) = d\operatorname{Vol}_{\mathbb{R}^{d}}(\mathcal{E}_{E}, \iota_{*}(X_{1}), \iota_{*}(X_{2}), \dots, \iota_{*}(X_{d'}))$$

$$= \det(z, v_{1}, v_{2}, \dots, v_{d'})$$

$$= \frac{1}{|w|} \det(n, v_{1}, v_{2}, \dots, v_{d'})$$

$$= \frac{1}{|\nabla P_{\eta}|} d\operatorname{Vol}_{\mathbb{R}^{d}}(N_{\eta}, \iota_{*}(X_{1}), \iota_{*}(X_{2}), \dots, \iota_{*}(X_{d'}))$$

$$= \frac{1}{|\nabla P|} d\operatorname{Vol}_{S}(X_{1}, X_{2}, \dots, X_{d'}).$$

By virtue of the preceding proposition, we shall denote by $d\sigma_P$ the unique smooth d'-form on S which satisfies

$$d\sigma_P = d\sigma_{P,E} = \frac{1}{|\nabla P|} d \operatorname{Vol}_S$$
(3.23)

for all $E \in \text{Exp}(P)$. In this notation, we have this section's central result.

Theorem 3.2.4. Let P be a smooth positive-homogeneous function and let $S = \{ \eta \in \mathbb{R}^d : P(\eta) = 1 \}$. Then S is a compact smooth embedded hypersurface of \mathbb{R}^d . Viewing \mathbb{R}^d as an oriented Riemannian manifold with its usual orientation and metric \overline{g} , $N = \nabla P/|\nabla P|$ is a smooth unit normal vector field along S. As a submanifold of \mathbb{R}^d , S is a oriented Riemannian manifold of dimension d' = d - 1 with its induced Riemannian metric g^S , volume form $d\operatorname{Vol}_S \in \Omega^{d'}(S)$ and orientation determined by N. The σ -algebras Σ_P and $\mathcal{L}(S)$ on S coincide and the smooth d'-form $d\sigma_P \in \Omega^{d'}(S)$, defined by (3.23), coincides with the measure σ_P in the sense that

$$\int_{S} g(\eta) \, \sigma_{P}(d\eta) = \int_{S} g \, d\sigma_{P} \tag{3.24}$$

for all $g \in L^1(S, \Sigma_P, \sigma_P)$; here, the left hand side represents the Lebesgue integral of g with respect to σ_P and the right hand side is the integral of the measurable d'-form g $d\sigma_P$. Furthermore, the measure σ_P and the canonical Riemannian volume measure Vol_S on S are mutually absolutely continuous.

The following two lemmas will take us to the proof of this theorem. The first lemma lets us relate the smooth d'-form $d\sigma_{P,E}=d\sigma_P$ to the Jacobian induced by some parameterization φ^{-1} via the Riemannian volume form $d\operatorname{Vol}_{\mathbb{R}^d}$ and $d\operatorname{Vol}_S$.

Lemma 3.2.5. Let $(\mathcal{U}, \varphi) \in \mathcal{A}(S)$ and $E \in \operatorname{Exp}(P)$. Set $U = \varphi(\mathcal{U}) \subseteq \mathbb{R}^{d'}$, $V = (0, 1) \times U$ and define $\rho_{E,\varphi} : V \to \mathbb{R}^d$ and $h_{E,\varphi} : U \to \mathbb{R}$, respectively, by

$$\rho_{E,\varphi}(y) = \psi_E(r,\varphi^{-1}) = r^E \varphi^{-1}(u)$$

for
$$y = (r, u) \in V$$
 and

$$h_{E,\varphi}(u) = \det \left(E\varphi^{-1}(u) \middle| D_u \varphi^{-1}(u) \right)$$

for $u \in U$; here, the vertical bar separates the first column of the (necessarily) $d \times d$ matrix from the rightmost $d \times d'$ submatrix and D_u denotes the Jacobian in the coordinates $u = (u^1, u^2, \dots, u^{d'}) \in U$. Then, $\rho_{E,\varphi}$ is a diffeomorphism onto its image $\rho_{E,\varphi}(V) = \widetilde{\mathcal{U}}_E$ and its Jacobian matrix $D\rho_{E,\varphi}$ has

$$\det(D\rho_{E,\varphi}(y)) = r^{\mu_P - 1} h_{E,\varphi}(u) \tag{3.25}$$

for all $y = (r, u) \in V$. Furthermore, $h_{E,\varphi}$ is everywhere non-zero, smooth and

$$d\sigma_P = h_{E,\varphi}(u) du^1 \wedge du^2 \wedge \dots \wedge du^{d'}$$
(3.26)

in the coordinates $u=(u^1,u^2,\ldots,u^{d'})\in U.$

Proof. The map $\rho_{E,\varphi}$ is smooth because φ is smooth. By virtually the same argument made in the proof of Proposition 3.1.3, which here uses the fact that $P \in C^{\infty}(\mathbb{R}^d)$, we conclude that $\rho_{E,\varphi}$ is a diffeomorphism onto

$$\rho_{E,\varphi}(V) = \bigcup_{0 < r < 1} r^E \mathcal{U} = \widetilde{\mathcal{U}_E}$$

with inverse $\rho_{E,\varphi}^{-1}(x)=(P(x),\varphi((P(x))^{-E}x)).$ For $y=(r,u)\in V$, observe that

$$D\rho_{E,\varphi}(y) = \left(\frac{d}{dr}(r^{E}\varphi^{-1}(u))\middle|D_{u}\left[r^{E}\varphi^{-1}(u)\right]\right)$$
$$= \left(r^{E-I}E\varphi^{-1}(u)\middle|r^{E}D_{u}\varphi^{-1}(u)\right)$$
$$= \left.r^{E}\left(\frac{1}{r}E\varphi^{-1}(u)\middle|D_{u}\varphi^{-1}(u)\right).$$

Using properties of the determinant, we have

$$\det(D\rho_{E}(y)) = \det(r^{E}) \det\left(\frac{1}{r}E\varphi^{-1}(u) \middle| D_{u}\varphi^{-1}(u)\right)$$
$$= r^{\operatorname{tr} E}r^{-1} \det\left(E\varphi^{-1}(u) \middle| D_{u}\varphi^{-1}(u)\right)$$
$$= r^{\mu_{P}-1}h_{E,\varphi}(u)$$

for all $y=(r,u)\in V$ thus proving (3.25). It is clear that $h_{E,\varphi}$ is smooth and, by virtue of the fact that $\rho_{E,\varphi}:(0,1)\times U\to \widetilde{\mathcal{U}_E}$ is a diffeomorphism, (3.25) guarantees that $h_{E,\varphi}$ is everywhere non-vanishing. Finally, in view of Proposition 3.2.2 (and (3.23)), we find that

$$h_{E,\varphi}(u) = \det \left(E\varphi^{-1}(u) \middle| D_u\varphi^{-1}(u) \right)$$

$$= (d\operatorname{Vol}_{\mathbb{R}^d})_{\varphi^{-1}(u)} (\mathcal{E}_E, \iota_*(\partial_{u^1}), \iota_*(\partial_{u^2}), \dots, \iota_*(\partial_{u^{d'}}))$$

$$= (d\sigma_P)_{\varphi^{-1}(u)} (\partial_{u^1}, \partial_{u^2}, \dots, \partial_{u^{d'}})$$

for $u \in U$ and so (3.26) is satisfied. ³

³For more details on conventional notation in smooth manifold theory, the reader may reference [9] or see Appendix 3.2.

Lemma 3.2.6. Let $g \in L^1(S, \Sigma_P, \sigma_P)$ be supported on the domain of some chart on S. Then, for any $(\mathcal{U}, \varphi) \in \mathcal{A}^+(S)$ such that $\operatorname{Supp}(g) \subseteq \mathcal{U}$, the pushforward $(\varphi^{-1})^*(g) = g \circ \varphi^{-1}$ is Lebesgue measurable on $U = \varphi(\mathcal{U}) \subseteq \mathbb{R}^{d'}$ and

$$\int_{S} g(\eta)\sigma_{P}(d\eta) = \int_{U} (\psi^{-1})^{*}(g \, d\sigma_{P}).$$

Proof. Let (\mathcal{U}, φ) be a chart on S for which $\mathrm{Supp}(g) \subseteq \mathcal{U}$ and assume the notation of Lemma 3.2.5. Given $E \in \mathrm{Exp}(P)$, observe that

$$f(x) = \mu_P \chi_{(0,1)}(P(x))g(P(x)^{-E}x),$$

defined for $x \in \mathbb{R}^d \setminus \{0\}$, is supported on $\widetilde{\mathcal{U}}_E = \rho_{E,\varphi}(V)$. An appeal to Corollary 3.1.2 guarantees that f is absolutely integrable on $\rho_{E,\varphi}(V)$ and

$$\int_{S} g(\eta) \, \sigma_{P}(d\eta) = \int_{\rho_{E,\varphi}(V)} f(x) \, dx. \tag{3.27}$$

Given that $\rho_{E,\varphi}$ is a diffeomorphism, an appeal to Theorem 15.11 of [2] guarantees that

$$\int_{\rho_{E,\varphi}(V)} f(x) dx = \int_{V} f(\rho_{E,\varphi}(y)) |\det(D\rho_{E,\varphi}(y))| dy; \tag{3.28}$$

in particular, $V \ni y \mapsto f(\rho_{E,\varphi}(y))|\det(D\rho_{E,\varphi}(y))$ is Lebesgue measurable on \mathbb{R}^d . Let us now view the Lebesgue measure dy on \mathbb{R}^d as the completion of the product measure $dr \times du$ on the product space $\mathbb{R} \times \mathbb{R}^{d'}$. By virtue of Lemma 3.2.5, we have

$$f(\rho_{E,\varphi}(y))|\det(D\rho_{E,\varphi}(y))| = f(r^{E}\varphi^{-1}(u)) r^{\mu_{P}-1}|h_{E,\varphi}(u)|$$

$$= \mu_{P} \chi_{(0,1)}(r)g(\varphi^{-1}(u)) r^{\mu_{P}-1}|h_{E,\varphi}(u)|$$

$$= \mu_{P} r^{\mu_{P}-1}g(\varphi^{-1}(u))|h_{E,\varphi}(u)|$$
(3.29)

for $y = (r, u) \in V$. By virtue of Fubini's theorem, dr-almost every $r \in (0, 1)$, the r-section

$$U \ni u \mapsto \mu_P r^{\mu_P - 1} g(\varphi^{-1}(u)) |h_{E,\varphi}(u)|$$

is Lebesgue measurable and, upon recalling that $h_{E,\varphi}$ is smooth and everywhere nonzero, we conclude that $(\varphi^{-1})^*(g) = g \circ \varphi^{-1}$ is Lebesgue measurable on U. In view of (3.29), Fubini's theorem also guarantees that

$$\int_{V} f(\rho_{E,\varphi}(y)) |\det(D\rho_{E,\varphi}(y))| \, dy = \int_{(0,1)} \int_{U} \mu_{P} \, r^{\mu_{P}-1} g(\varphi^{-1}(u)) |h_{E,\varphi}(u)| \, du \, dr$$

$$= \int_{0}^{1} \mu_{P} \, r^{\mu_{P}-1} \left(\int_{U} g(\varphi^{-1}(u)) |h_{E,\varphi}(u)| \, du \right) \, dr$$

$$= \int_{U} g(\varphi^{-1}(u)) |h_{E,\varphi}(u)| \, du \tag{3.30}$$

By combining (3.27), (3.28) and (3.30), we have shown that $(\varphi^{-1})^*g = g \circ \varphi^{-1}$ is Lebesgue measurable on $U = \varphi(\mathcal{U})$ and

$$\int_{S} g(\eta) \sigma_{P}(d\eta) = \int_{U} g(\varphi^{-1}(u)) |h_{E,\varphi}(u)| du.$$

Finally, if $(\mathcal{U}, \varphi) \in \mathcal{A}_+(S)$, an appeal to Proposition 3.2.2 and (3.26) of Lemma 3.2.5 guarantees that

$$|h_{E,\varphi}(u)| = h_{E,\varphi}(u) = (d\sigma_P)_{\varphi^{-1}(u)}(\partial_{u^1}, \partial_{u^2}, \dots, \partial_{u^{d'}}) > 0$$

for all $u = (u^1, u^2, \dots, u^{d'}) \in U$ and thus

$$\int_{S} g(\eta)\sigma_{P}(d\eta) = \int_{U} g(\varphi^{-1}(u))h_{E,\varphi}(u) du$$

$$= \int_{U} g(\varphi^{-1}(u)) h_{E,\varphi}(u) du^{1} \wedge du^{2} \wedge \cdots \wedge du^{d'}$$

$$= \int_{U} (\varphi^{-1})^{*}(g \cdot d\sigma_{P}),$$

as desired.

Remark 3. In studying the proof of Lemma 3.2.6, we deduce a slightly more general statement: Given any chart $(\mathcal{U}, \varphi) \in \mathcal{A}(S)$ and $g \in L^1(S, \Sigma_P, \sigma_P)$ for which $\operatorname{Supp}(g) \in \mathcal{U}$, $(\varphi^{-1})^*g$ is Lebesgue measurable on $U = \varphi(\mathcal{U})$ and

$$\int_{S} g(\eta) \sigma_{P}(d\eta) = \begin{cases} \int_{U} (\varphi^{-1})^{*}(g \, d\sigma_{P}) & \text{if } (\mathcal{U}, \varphi) \in \mathcal{A}_{+}(S) \\ -\int_{U} (\varphi^{-1})^{*}(g \, d\sigma_{P}) & \text{if } (\mathcal{U}, \varphi) \in \mathcal{A}_{-}(S) \coloneqq \mathcal{A}(S) \setminus \mathcal{A}_{+}(S). \end{cases}$$

Proof of Theorem 3.2.4. In view of Proposition 3.2.1 and the discussion following its proof, it remains to prove the assertions in the last two sentences in the statement of the theorem. Given any $F \in \Sigma_P$ and chart $(\mathcal{U}, \varphi) \subseteq \mathcal{A}(S)$, we have $\chi_{\varphi(F \cap \mathcal{U})} = (\varphi^{-1})^* (\chi_{F \cap \mathcal{U}})$ is Lebesgue measurable on $U = \varphi(\mathcal{U})$ and so it follows (Exercise 4.6.2 of [11]) that $F \in \mathcal{L}(S)$. Consequently, $\mathcal{B}(S) \subseteq \Sigma_P \subseteq \mathcal{L}(S)$.

Let $g \in L^1(S, \Sigma_P, \sigma_P)$, which is necessarily Lebesgue measurable on S in view of the results of the previous paragraph. Now, let $\{(\mathcal{U}_j, \varphi_j)\} \subseteq \mathcal{A}_+(S)$ be a countable atlas on S and let $\{\kappa_j\}$ be a smooth partition of unity subordinate to the cover $\{\mathcal{U}_j\}$. For each $j \in \mathbb{N}$, observe that $\kappa_j g \in L^1(S, \Sigma_P, \sigma_P)$ and has $\operatorname{Supp}(\kappa_j g) \subseteq \mathcal{U}_j$. By virtue of Lemma 3.2.6, we have $(\varphi_j^{-1})^*(\kappa_j g \, d\sigma_P)$ is integrable on $U_j = \varphi(\mathcal{U}_j)$ and

$$\int_{S} \kappa_{j}(\eta)g(\eta)\sigma_{P}(d\eta) = \int_{U_{j}} (\varphi_{j}^{-1})^{*}(\kappa_{j}g\,d\sigma_{P})$$
(3.31)

for each $j \in \mathbb{N}$. With the help of Proposition 3.2.2, it is easy to see that $g d\sigma_P$ and $\kappa_j g d\sigma_P$ (for $j \in \mathbb{N}$) are Lebesgue measurable d'-forms on S. In view of (4.4.6) of [11], (3.31) ensures that, for each $j \in \mathbb{N}$, $\kappa_j g d\sigma_P$ is integrable (in the sense of forms) on S and

$$\int_{S} \kappa_{j}(\eta)g(\eta)\sigma_{P}(d\eta) = \int_{S} \kappa_{j}g \, d\sigma_{P}.$$

By the monotone convergence theorem, it follows that

$$\sum_{j=1}^{\infty} \left| \int_{S} \kappa_{j} g \, d\sigma_{P} \right| = \sum_{j=1}^{\infty} \left| \int_{S} \kappa_{j}(\eta) g(\eta) \sigma(d\eta) \right| \leq \sum_{j=1}^{\infty} \int_{S} \kappa_{j}(\eta) |g(\eta)| \sigma_{P}(d\eta) = \|g\|_{L^{1}(S, \Sigma_{P}, \sigma_{P})} < \infty.$$

Therefore, in view of the construction on p. 242 of [11], we conclude that the d'-form $g d\sigma_P$ is integrable and

$$\int_{S} g(\eta)\sigma_{P}(d\eta) = \sum_{j=1}^{\infty} \int_{S} \kappa_{j}(\eta)g(\eta)\sigma_{P}(d\eta) = \sum_{j=1}^{\infty} \int_{S} \kappa_{j}g \, d\sigma_{P} = \int_{S} g \, d\sigma_{P}$$

by virtue of the dominated convergence theorem; this is (3.24).

Finally, given that $|\nabla P|$ is continuous and non-vanishing on the compact set S,

$$C_1 := \inf_S \frac{1}{|\nabla P|}$$
 and $C_2 := \sup_S \frac{1}{|\nabla P|}$

are both positive real numbers. For each $F \in \Sigma_P \subseteq \mathcal{L}(S)$, (3.24) guarantees that

$$\sigma_P(F) = \int_S \chi_F(\eta) \sigma_P(d\eta) = \int_S \chi_F d\sigma_P$$

By virtue of Proposition 3.2.2, it follows that

$$C_1 \operatorname{Vol}_S(F) = C_1 \int_S \chi_F \, d\operatorname{Vol}_S \le \int_S \frac{\chi_F}{|\nabla P|} d\operatorname{Vol}_S = \int_S \chi_F \, d\sigma_P = \sigma_P(F)$$

and

$$\sigma_P(F) = \int_S \chi_F d\sigma_P = \int_S \frac{\chi_F}{|\nabla P|} d\operatorname{Vol}_S \le C_2 \int_S \chi_F d\operatorname{Vol}_S = C_2 \operatorname{Vol}_S(F)$$

where we have used the definition of the Riemannian volume measure on S, c.f., [1]. In short, there are positive constants C_1 and C_2 for which

$$C_1 \operatorname{Vol}_S(F) \le \sigma_P(F) \le C_2 \operatorname{Vol}_S(F)$$
 (3.32)

for all $F \in \Sigma_P$. In particular, (3.24) holds for all $F \in \mathcal{B}(S)$ and so it follows that the completions of the σ -algebra $\mathcal{B}(S)$ with respect to σ_P and Vol_S coincide. We know, however, that Vol_S , which is defined on $\mathcal{L}(S)$, is a Radon measure (Proposition 1.5 [1, Chapter XII]) and, by virtue of Proposition 3.1.15, it follows that $\Sigma_P = \mathcal{L}(S)$ and (3.32) holds for all F in this common σ -algebra. Thus σ_P and Vol_S are mutually absolutely continuous and the theorem is proved.

We immediately obtain the following corollary which allows us to compute the Lebesgue integral with respect to σ_P in coordinates.

Corollary 3.2.7. Let $g \in L^1(S, \mathcal{L}(S), \sigma_P)$. Then, given any countable (or finite) atlas $\{(\mathcal{U}_j, \varphi_j)\}\subseteq \mathcal{A}_+(S)$, smooth partition of unity $\{\kappa_i\}$ subordinate to $\{\mathcal{U}_i\}$, and $E \in \operatorname{Exp}(P)$,

$$\int_{S} g(\eta) \sigma_{P}(d\eta) = \sum_{i} \int_{S} \kappa_{i} g \, d\sigma_{P} = \sum_{i} \int_{U_{j}} \kappa_{j}(\varphi^{-1}(u)) g(\varphi^{-1}(u)) h_{E,\varphi_{j}}(u) \, du$$

where, for each
$$j$$
, $U_j = arphi_j(\mathcal{U}_j) \subseteq \mathbb{R}^{d'}$ and

$$h_{E,\varphi_j}(u) = \det(E\varphi^{-1}(u)|D_u\varphi^{-1})$$

for
$$u = (u^1, u^2, \dots, u^{d'}) \in U_j$$
.

Chapter 4

Van der Corput Lemmas and Estimating Oscillatory Integrals

In this chapter, we are interested in some technical results related to estimating one-dimensional oscillatory integrals of the form

$$\int_{a}^{b} e^{i\Phi(x)} \psi(x) dx, \tag{4.1}$$

where Φ, ψ are functions which satisfy certain hypotheses. In particular, Φ is a real-valued function. Here, $\Phi(x)$ is often referred to as the **phase** function, while $\psi(x)$ the **amplitude**. Integrals of this form often appear in Fourier analysis and are treated extensively in this context (see Chapter 8 of [16], for example). In Section 4.1, we present a number of propositions and lemmas and prove a version of the Van der Corput lemma (Lemma 4.1.6), which is a key ingredient for obtaining our sup-norm estimate of convolution powers. The main approach we take to prove our preliminary results is fairly standard; however, we also take a more general approach by considering sub-level set estimates inspired by [12], which then gives the Van der Corput lemma as a special case. In Section 4.2, we apply our result to obtain decay estimates for the Fourier transform of the surface measure σ_P associated with a positive homogeneous polynomial P along the coordinate axes. These estimates are a partial improvement of the results in [7].

4.1 Estimating Oscillatory Integrals & the Van der Corput lemma

The following result helps us isolate the influences of the amplitude ψ and the phase Φ on the decay of (4.1); its proof makes use of integration by parts.

Lemma 4.1.1. Let $h \in L^1([a,b])$ and $g \in C^1([a,b])$ be complex-valued. Then for any M such that

$$\left| \int_{a}^{x} h(u) \, du \right| \le M$$

for all $x \in [a, b]$ we have

$$\left| \int_{a}^{b} g(u)h(u) du \right| \leq M \left(\left\| g' \right\|_{1} + \left\| g \right\|_{\infty} \right).$$

Proof. Since $h \in L^1([a,b])$, the function

$$f(x) = \int_{a}^{x} h(u) \, du$$

is absolutely continuous and f'(x) = h(x) almost everywhere. Further, because $|f(x)| \le M$ for all $x \in [a, b]$, we have, by integration by parts,

$$\int_{a}^{b} g(u)h(u) \, du = g(u)f(u) \Big|_{a}^{b} - \int_{a}^{b} g'(u)f(u) \, du.$$

It follows that

$$\left| \int_{a}^{b} g(u)h(u) \, du \right| \le |f(b)g(b)| + 0 + \int_{a}^{b} |f(u)| |g'(u)| \, du \le M \left(||g||_{\infty} + ||g'||_{1} \right).$$

Remark 4. The assumption that g is once continuously differentiable can be weakened slightly to ask that g belongs to the Sobolev space $W^{1\,1}([a,b])$. For the essential details, we refer the reader to [10] and specifically Theorem 7.16.

Proposition 4.1.2. Let $\Phi \in C^2([a,b])$ be a real-valued function for which Φ' is monotonic and $|\Phi'(x)| \ge \delta > 0$ for all $x \in [a,b]$, then

$$\left| \int_{a}^{b} e^{i\Phi(x)} \, dx \right| \le \frac{4}{\delta}.$$

Proof. Observe that

$$\int_{a}^{b} e^{i\Phi(x)} dx = \int_{a}^{b} \frac{1}{i\Phi'(x)} \left(\frac{d}{dx}e^{i\Phi}\right) dx$$
$$= \frac{1}{i\Phi'(x)} e^{i\Phi(x)} \Big|_{a}^{b} - \int_{a}^{b} e^{i\Phi(x)} \frac{d}{dx} \left(\frac{1}{i\Phi'(x)}\right) dx$$

where we have integrated by parts and used the fact that Φ' never vanishes. Consequently,

$$\left| \int_a^b e^{i\Phi(x)} \, dx \right| \le \left| \frac{1}{\Phi'(b)} - \frac{1}{\Phi'(a)} \right| + \int_a^b \left| \frac{d}{dx} \left(\frac{1}{\Phi'(x)} \right) \right| \, dx.$$

Since Φ' is monotonic, the absolute value in the second term of the right hand side can be brought outside of the integral. Subsequently, we apply fundamental theorem of calculus to this term to obtain the first assertion by taking

$$C = \left| \frac{1}{\Phi'(b)} - \frac{1}{\Phi'(a)} \right| + \int_a^b \left| \frac{d}{dx} \left(\frac{1}{\Phi'(x)} \right) \right| dx = 2 \left| \frac{1}{\Phi'(b)} - \frac{1}{\Phi'(a)} \right| \le \frac{4}{\delta}.$$

We note that the function Φ in the preceding proposition has no critical points on [a, b]. In what follows, we will extend the result of Proposition 4.1.2 by obtaining an estimate for

$$\left| \int_{a}^{b} e^{i\Phi(x)} \, dx \right|$$

where Φ is a real-valued function on [a,b] such that for some $k \geq 2$ and $\delta > 0$ we have $\Phi \in C^k([a,b])$ and $|\Phi^{(k)}(x)| \geq \delta$ for all $x \in [a,b]$. To this end, we require the help from the forthcoming results (Proposition 4.1.3 and Lemma 4.1.4) to count and bound the (Lebesgue) measures of sublevel sets of Φ , which have the form

$$E_{\alpha} = \{ t \in [a, b] : |\Phi(t)| \le \alpha \}$$

for each $\alpha > 0$. Proposition 4.1.3 is inspired by Proposition 1.2 of [12]. The proof of Proposition 4.1.3 is fairly involved and makes use of the Lagrange interpolation polynomials, while the proof of Lemma 4.1.4 is a straightforward induction argument on the number of critical points of f.

Proposition 4.1.3. For an integer $k \ge 1$, suppose that $\Phi \in C^k([a,b])$ is real-valued and, for some $\delta > 0$, $|\Phi^{(k)}(t)| \ge \delta$ for all $t \in [a,b]$. Then, for each $\alpha > 0$, the sub-level set $E_\alpha := \{t \in [a,b] : |\Phi(t)| \le \alpha\}$ has

$$m(E_{\alpha}) \le 2k \left(\frac{\alpha}{\delta}\right)^{1/k},$$

where m denotes the Lebesgue measure on \mathbb{R} .

Proof. Since Φ is continuous, there is an integer N for which $E_{\alpha} = \bigcup_{i=1}^{N} E_i$ where E_i for i = 1, 2, ..., N are disjoint compact intervals for which

$$\max E_1 \le \min E_2 \le \max E_2 \le \cdots \le \max E_{N-1} \le \min E_N.$$

Slide each E_i along \mathbb{R} such that $\max E_i = \min E_{i+1}$ for $i=1,2,\ldots k-1$ to create a single closed interval I with $m(I) = m(E_{\alpha})$. Pick k+1 equally spaced points x'_0, x'_1, \ldots, x'_k in I and move the intervals E_i back to their original positions. The selected points are now x_0, x_1, \ldots, x_k and satisfy

$$|x_j - x_l| \ge m(E_\alpha) \frac{|j - l|}{k}$$
, for $j, l = 0, 1, \dots, k$ (4.2)

Let us now consider the Lagrange interpolation polynomial

$$h(x) = \sum_{j=0}^{k} \Phi(x_j) p_j(x)$$

where

$$p_j(x) = \prod_{\substack{l=0\\l\neq j}}^k \frac{(x-x_j)}{(x_l-x_j)}$$

for $j=0,1,\ldots,k$. As designed, each p_j is a kth-order polynomial and h interpolates Φ at the points x_0,x_1,\ldots,x_k . In particular, $F(x):=h(x)-\Phi(x)$ vanishes at x_0,x_1,\ldots,x_k and is necessarily

k-times continuously differentiable on [a,b]. It follows that there are k points $y_1 < y_2 < \cdots < y_k$ such that $x_0 < y_1 < x_1 < y_2 < x_2 < \cdots < y_k < x_k$ and

$$F'(y_1) = F'(y_2) = \dots = F'(y_k) = 0.$$

Upon repeating this argument an additional k-1 times, we find that $F^{(k)}(\xi)=0$ at some $\xi\in(a,b)$ and here,

$$\Phi^{(k)}(\xi) = h^{(k)}(\xi) = \sum_{j=0}^{k} \Phi(x_j) p_j^{(k)}(\xi).$$

By virtue of (4.2), for each j = 0, 1, 2, ... k,

$$|p_j^{(k)}(\xi)| = k! \prod_{\substack{l=0\\l\neq j}}^k \frac{1}{|x_l - x_j|} \le k! \left(\frac{k^k}{m(E_\alpha)^k} \prod_{\substack{l=0\\l\neq j}}^k \frac{1}{|j-l|} \right) = \frac{k^k}{m(E_\alpha)^k} \frac{k!}{j!(k-j)!} = \frac{k^k}{m(E_\alpha)^k} \binom{k}{j}.$$

Because $\delta \leq |\Phi^{(k)}(\xi)|$ and $|\Phi(x_j)| \leq \alpha$ for each $j = 0, 1, 2, \dots, k$, it follows that

$$\delta \le \left| \sum_{j=0}^{k} \Phi(x_j) p_j^{(k)}(\xi) \right| \le \sum_{j=0}^{k} \alpha |p_j^{(k)}(\xi)| \le \frac{\alpha k^k}{m(E_\alpha)^k} \sum_{j=0}^{k} \binom{k}{j} = \frac{\alpha (2k)^k}{m(E_\alpha)^k}$$

where we have made use of the binomial theorem. Consequently,

$$m(E_{\alpha}) \le 2k \left(\frac{\alpha}{\delta}\right)^{1/k}$$
.

Lemma 4.1.4. Given an integer $k \ge 1$, suppose that $f \in C^k([a,b])$ is a real-valued function for which $|f^{(k)}(t)| > 0$ for all $t \in [a,b]$. Then [a,b] has a minimal cover of compact subintervals I_1, I_2, \ldots, I_l with $l \le k$ and for which the restriction of f to each subinterval is strictly monotonic.

Proof. We will prove this lemma by induction on k. When k=1, the hypothesis ensures that f is strictly monotonic on [a,b], as asserted. We now suppose that, for $k\geq 1$, $f\in C^{(k+1)}([a,b])$ and $f^{(k)}(t)>0$ for all $t\in [a,b]$; the case in which $f^{(k)}(t)<0$ is similar. Our inductive hypothesis guarantees a minimal cover I_1,I_2,\ldots,I_l of [a,b] where $l\leq k$ and each I_j is a compact interval on which f' is strictly monotonic. Given that a strictly monotonic function can intersect 0 at most once on its domain, f must have, at most, f critical points in the interval f0, one for each interval f1, f2, ..., f3. If there are no critical points, then f1 is necessarily strictly monotonic on f2, f3 and the desired result holds. Otherwise, denote by f3, f4, f5, f7, f8. In this case, f8 is strictly monotonic on the f8. In this case, f8 is strictly monotonic on the f8 is one constant f8. In this case, f8 is strictly monotonic on the f8 in the assertion holds for f8. In this case, f8 is strictly monotonic on the f8 is assertion holds for f8. In this case, f8 is a strictly monotonic on the f8 is assertion holds for f8. In this case, f9 is strictly monotonic on the f9 is assertion holds for f9 in the hypothesis ensures that f9 is assertion holds for f8 in the hypothesis ensures that f9 is a strictly monotonic on the f9 is assertion holds for f9 in the hypothesis ensures that f9 is at f1 in the hypothesis ensures that f1 is a strictly monotonic on the f1 in the hypothesis ensures that f1 is a strictly monotonic ensures that f2 is a strictly monotonic ens

Now we are ready to prove the following extension of Proposition 4.1.2 to deal with cases where $\Phi(x)$ has critical points. From here, Lemma 4.1.1 gives us an estimate for more general oscillatory integrals where $\Phi(x)$ is not necessarily the constantly 1 functions. A special case of this (where k=2) is a version of the Van der Corput lemma which we will state as Lemma 4.1.6 and prove.

Proposition 4.1.5. Let Φ be a real-valued function on [a,b] and suppose that, for some $k \geq 2$, $\Phi \in C^k([a,b])$ and $|\Phi^{(k)}(x)| \geq \delta$ for all $x \in [a,b]$ where $\delta > 0$. Then

$$\left| \int_a^b e^{i\Phi(x)} \, dx \right| \le \frac{2k \cdot (2(k-1))^{\frac{1}{k}}}{\delta^{\frac{1}{k}}}.$$

Proof. Define, for $\alpha > 0$,

$$I(\alpha) = \int_{\{x \in [a,b]: |\Phi'(x)| > \alpha\}} e^{i\Phi(x)} dx$$

and observe that

$$\left| \int_{a}^{b} e^{i\Phi(x)} dx \right| \leq \left| \int_{\{x \in [a,b]: |\Phi'(x)| \ge \alpha\}} e^{i\Phi(x)} dx \right| + \left| \int_{\{x \in [a,b]: |\Phi'(x)| < \alpha\}} e^{i\Phi(x)} dx \right|$$

$$\leq |I(\alpha)| + m(E_{\alpha})$$

$$\leq |I(\alpha)| + 2(k-1) \left(\frac{\alpha}{\delta}\right)^{\frac{1}{k-1}}$$

for all $\alpha > 0$ by virtue of Proposition 4.1.3.

Let's estimate $I(\alpha)$. In view of Lemma 4.1.4, there is a minimal cover of [a,b] containing at most (k-1) compact intervals on which Φ' is strictly monotonic. For any $\alpha>0$, by removing the set E_{α} from [a,b], we obtain at most 2(k-1) compact intervals I_1,I_2,\ldots,I_l covering $\{x\in [a,b]: |\Phi'|\geq \alpha\}=[a,b]\setminus E_{\alpha}$ and on each of which Φ' is strictly monotonic. Applying Proposition 4.1.2 to each such interval and using the triangle inequality we find

$$|I(\alpha)| \le l\frac{2}{\alpha} \le 2(k-1)\frac{2}{\alpha} = \frac{4(k-1)}{\alpha}$$

which holds for all $\alpha > 0$. Consequently,

$$\left| \int_a^b e^{i\Phi(x)} \, dx \right| \le \frac{4(k-1)}{\alpha} + 2(k-1) \left(\frac{\alpha}{\delta} \right)^{1/(k-1)} = 2(k-1) \left(\frac{2}{\alpha} + \left(\frac{\alpha}{\delta} \right)^{1/(k-1)} \right)$$

for all $\alpha > 0$ and therefore

$$\left| \int_{a}^{b} e^{i\Phi(x)} dx \right| \le 2(k-1) \inf_{\alpha > 0} \left(\frac{2}{\alpha} + \left(\frac{\alpha}{\delta} \right)^{1/(k-1)} \right).$$

It is easily shown that this infimum is attained at

$$\alpha_0 = \delta^{\frac{1}{k}} [2(k-1)]^{\frac{k-1}{k}}$$

and so it follows that,

$$\left| \int_{a}^{b} e^{i\Phi(x)} \, dx \right| \le 2(k-1) \left(\frac{2}{\alpha_0} + \left(\frac{\alpha_0}{\delta} \right)^{\frac{1}{k-1}} \right) = \frac{2k \cdot (2(k-1))^{\frac{1}{k}}}{(\delta)^{\frac{1}{k}}},$$

as desired. It follows that for any choice of $\alpha > 0$,

$$\left| \int_a^b e^{i\Phi(x)} \, dx \right| \le \frac{4(k-1)}{\alpha} + 2(k-1) \left(\frac{\alpha}{\delta} \right)^{\frac{1}{k-1}} = 2(k-1) \left[\frac{2}{\alpha} + \left(\frac{\alpha}{\delta} \right)^{\frac{1}{k-1}} \right] =: f(\alpha),$$

Now,

$$f'(\alpha) = \frac{2}{\alpha^2} \left[\alpha^{\frac{k}{k-1}} \delta^{\frac{-1}{k-1}} - 2(k-1) \right] = 0 \iff \alpha = \delta^{\frac{1}{k}} \left[2(k-1) \right]^{\frac{k-1}{k}}.$$

With this choice of α , we have that $f(\alpha)$ attains a minimum:

$$f''(\alpha) = \frac{2(k-1)}{\alpha^3} \left[4 - \frac{k-2}{(k-1)^2} \alpha^{\frac{k}{k-1}} \delta^{\frac{-1}{k-1}} \right]$$

$$= \frac{2(k-1)}{\alpha^3} \left[4 - \frac{k-2}{(k-1)^2} 2(k-1) \delta^{\frac{1}{k-1}} \delta^{\frac{-1}{k-1}} \right]$$

$$= \frac{4(k-1)}{\alpha^3} \left(2 - \frac{k-2}{k-1} \right)$$

$$\geq 0$$

and that

$$\begin{split} \left| \int_{a}^{b} e^{i\Phi(x)} \, dx \right| & \leq 2(k-1) \left[\frac{2}{\alpha} + \left(\frac{\alpha}{\delta} \right)^{\frac{1}{k-1}} \right] \\ & = 2(k-1) \left[2\delta^{\frac{-1}{k}} \left[2(k-1) \right]^{\frac{1-k}{k}} + \left\{ \delta^{\frac{1}{k}-1} \left[2(k-1) \right]^{\frac{k-1}{k}} \right\}^{\frac{1}{k-1}} \right] \\ & = 2 \left[\frac{2(k-1)}{\delta} \right]^{\frac{1}{k}} + 2(k-1) \left[\frac{2(k-1)}{\delta} \right]^{\frac{1}{k}} \\ & = \frac{2k \cdot \left[2(k-1) \right]^{\frac{1}{k}}}{\delta^{\frac{1}{k}}} \end{split}$$

as desired. \Box

Finally, the following lemma is a version of the Van der Corput lemma which we will use to obtain sup-norm estimates for convolution powers in the next chapter.

Lemma 4.1.6. Let $g \in C^1([a,b])$ be complex-valued and let $\Phi \in C^2([a,b])$ be real-valued such that $\Phi''(x) \neq 0$ for all $x \in [a,b]$. Then

$$\left| \int_a^b g(u)e^{i\Phi(u)} du \right| \le \min \left\{ \frac{4}{\delta}, \frac{8}{\sqrt{\rho}} \right\} \left(\left\| g' \right\|_1 + \left\| g \right\|_{\infty} \right),$$

where $\delta = \inf_{x \in [a,b]} |\Phi'(x)|$ and $\rho = \inf_{x \in [a,b]} |\Phi''(x)|$.

Proof. Since Φ'' never changes sign on [a,b], Φ' must be monotonic on [a,b]. As a result, we can combine Proposition 4.1.2 and Proposition 4.1.5 (with k=2) to find

$$\left| \int_a^x e^{i\Phi(u)} \, du \right| \le \min\left\{ \frac{4}{\delta}, \frac{8}{\sqrt{\rho}} \right\}$$

for any $x \in [a,b]$. Setting $h(u) = e^{i\Phi(u)}$ we note that the functions g and h are those in Lemma 4.1.1, from which the desired result follows immediately.

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4.2 An application

An important application of the facts regarding oscillatory integrals in the preceding section is in the study of averaging operators and the decay of the Fourier transform of surface-carried measures. Here, we will follow the beginning of Chapter 8 of [16] to give some motivational context.

In \mathbb{R}^d , the averaging operator A that gives for each function f its average over the unit hypersphere centered at x is defined by

$$\mathcal{A}(f) = \frac{1}{\sigma(\mathbb{S}_d)} \int_{\mathbb{S}_{d-1}} f(x - y) \,\Theta(dy)$$

where Θ is the induced measure on the unit hypersphere \mathbb{S}_{d-1} . \mathcal{A} smooths f in several senses, the simplest one is given by Proposition 1.1 of [16], which states that the mapping $f \to \mathcal{A}(f)$ from $L^2(\mathbb{R}^d)$ to the Sobolev space $L^2_k(\mathbb{R}^d)$ with k = (d-1)/2 is bounded. The proof of this statement relies on a crucial fact, which is that the Fourier transform of the canonical spherical measure Θ , denoted by $\widehat{\Theta}$, can be written as

$$\widehat{\Theta}(\xi) = 2\pi |\xi|^{-d/2+1} J_{d/2}(2\pi |\xi|),$$

from which and the known asymptotic behavior of the Bessel functions J_m of order m one finds

$$\left|\widehat{\Theta}(\xi)\right| \le O(|\xi|^{-(d-1)/2}) \tag{4.3}$$

as $|\xi| \to \infty$. We thus see that the knowledge of how the Fourier transform of the spherical measure decays allows us to learn something about the averaging operator \mathcal{A} . The natural next step is to generalize the estimate in (4.3). Chapter 8.3 of [16] provides an excellent summary of results relating the curvature of hypersurfaces and the associated decay estimates. While this is interesting in its own right, for the remainder of this section, we shall turn our attention to directly obtaining decay estimates of the Fourier transform of a measure σ associated with a positive homogeneous polynomial P along the coordinate axes, using only the facts about oscillatory integrals provided in the preceding section. We will then compare our results to [7], which studies various mathematical objects derived from estimates on Fourier transforms of hypersurface measures, and demonstrate that when restricted to only finding decay estimates along the coordinate axes, our result is slightly more optimal.

To start, let us consider the polynomial P of the form

$$P(x_1, x_2) = x_1^2 + x_2^4.$$

An appeal to Corollary 3.2.7 gives us an expression for the Fourier transform of the measure $\widehat{\sigma_P}(\xi) = \widehat{\sigma_P}(\xi_1, \xi_2)$:

$$2\pi\widehat{\sigma_P}(\xi) = \int_S e^{-i\xi\cdot\eta} \sigma_P(d\eta)$$

where of course

$$S = \{ \eta \in \mathbb{R}^d : P(\eta) = 1 \}.$$

From Theorem 1.1 of [7], we have that for any $\xi \in \mathbb{R}^d$ with sufficiently large $|\xi|$, there is a constant C for which

$$|\widehat{\sigma_P}(\xi)| \le \frac{C}{|\xi|^{1/4}}.\tag{4.4}$$

¹For definitions of these mathematical objects, the reader may refer to Chapter 1 of [16].

With this result in mind, we shall restrict ourselves to only the coordinate axes, i.e., we now study how $|\widehat{\sigma_P}(\xi_1,0)|$ and $|\widehat{\sigma_P}(0,\xi_2)|$ decay. Given S, we consider a finite atlas of charts $\{(U_\alpha,\varphi_\alpha^{-1})\}$ with $\alpha=1,2,3,4$ given by

$$\begin{cases} \varphi_1: (-a, a) \to S & \varphi_1(u) = (\sqrt{1 - u^4}, u) \\ \varphi_2: (-a, a) \to S & \varphi_2(u) = -\varphi_1(u) \\ \varphi_3: (-b, b) \to S & \varphi_3(u) = (u, -(1 - u^2)^{1/4}) \\ \varphi_4: (-b, b) \to S & \varphi_4(u) = -\varphi_3(u). \end{cases}$$

where 0 < a < 1 and $(1 - a^2)^{1/4} = b$. Further, set $S_{\alpha} = \varphi_{\alpha}((-a, a))$ for $\alpha = 1, 2$ and $S_{\alpha} = \varphi_{\alpha}((-b, b))$ for $\alpha = 3, 4$. Figure 4.1 shows the parameterization associated with this atlas.

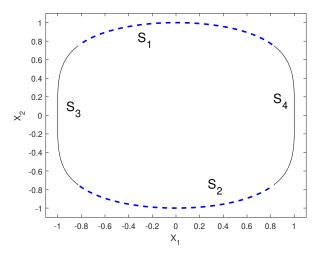


Figure 4.1: *S*

Note that the atlas $\{(U_{\alpha}, \varphi_{\alpha}^{-1})\}$ is defined so that it is positively oriented. We find the P-adapted surface area scaling factor to be

$$h_{E,\varphi_1}(u) = h_{E,\varphi_2}(u) = \det \begin{pmatrix} \frac{1}{2}\sqrt{1-u^4} & -\frac{2u^3}{\sqrt{1-u^4}} \\ \frac{1}{4}u & 1 \end{pmatrix} = \frac{1}{2\sqrt{1-u^4}}$$

for $u \in (-a, a)$ and

$$h_{E,\varphi_3}(u) = h_{E,\varphi_4}(u) = \det \begin{pmatrix} \frac{1}{2}u & 1\\ -\frac{1}{4}(1-u^2)^{1/4} & \frac{u}{2(1-u^2)^{3/4}} \end{pmatrix} = \frac{1}{4(1-u^2)^{3/4}}$$

for $u \in (-b,b)$. We want to break this integral up into a sum over the S_{α} 's. To this end, putting A = (-a,a), B = (-b,b), we let a partition of unity $\{\rho_{\alpha}\}_{\alpha}$ where $\alpha = 1,2,3,4$ subordinate to the

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cover $\{\varphi_1(A), \varphi_2(A), \varphi_3(B), \varphi_4(B)\}$ associated with the atlas above. With this, we have

$$2\pi\widehat{\sigma_P}(\xi) = \int_S e^{-i\xi\cdot\eta} \sigma_P(d\eta)$$

$$= \sum_{\alpha=1}^4 \int_S \rho_\alpha e^{-i\xi\cdot\eta} \sigma_P(d\eta)$$

$$= \sum_{\alpha=1}^4 \int_{S_\alpha} e^{-i\xi\cdot\eta} \sigma_P(d\eta)$$

$$= \int_{-a}^a e^{-i\xi\cdot\varphi_1(u)} h_{E,\varphi_1}(u) du + \int_{-a}^a e^{-i\xi\cdot\varphi_2(u)} h_{E,\varphi_2}(u) du$$

$$+ \int_{-b}^b e^{-i\xi\cdot\varphi_3(u)} h_{E,\varphi_3}(u) du + \int_{-b}^b e^{-i\xi\cdot\varphi_4(u)} h_{E,\varphi_4}(u) du.$$

4.2.1 Decay along the ξ_2 -axis

When $\xi = (0, \xi_2)$ with $\xi_2 \neq 0$ we have:

$$2\pi\widehat{\sigma_P}(\xi) = \int_{-a}^a e^{-i\xi_2 u} \frac{1}{2\sqrt{1-u^4}} du + \int_{-a}^a e^{+i\xi_2 u} \frac{1}{2\sqrt{1-u^4}} du + \int_{-b}^b e^{+i\xi_2(1-u^2)^{1/4}} \frac{1}{4(1-u^2)^{3/4}} du + \int_{-b}^b e^{-i\xi_2(1-u^2)^{1/4}} \frac{1}{4(1-u^2)^{3/4}} du + \int_{-b}^b e^{-i\xi_2(1-u^2)^{1/4}} \frac{1}{4(1-u^2)^{3/4}} du$$

$$\coloneqq I_1(\xi_2) + I_2(\xi_2) + I_3(\xi_2) + I_4(\xi_2).$$

For each of I_1 and I_2 , the phase function is real-valued, twice-differentiable and has monotonic first derivative whose absolute value is greater than or equal to $|\xi_2| > 0$ on (-a, a). Moreover, the amplitude function is once-differentiable on (-a, a). Thus, we can appeal to Lemma 4.1.1 and Proposition 4.1.2 to find a positive constant C for which

$$|I_1(\xi_2)| \le \frac{C}{|\xi_2|}$$
 and $|I_2(\xi_2)| \le \frac{C}{|\xi_2|}$.

To estimate $|I_3(\xi_2)|$ and $|I_4(\xi_2)|$, we first let $\Phi(u)=(1-u^2)^{1/4}$. We have that $\Phi'(u)=0$ if and only if u=0, and that the smallest integer k for which $\Phi^{(k)}(u)\neq 0$ for all $u\in [-a,a]$ is k=2. Another appeal to Lemma 4.1.6 yields

$$|I_3(\xi_2)| \leq rac{C}{|\xi_2|^{1/2}} \quad ext{and} \quad |I_4(\xi_2)| \leq rac{C}{|\xi_2|^{1/2}}.$$

It follows that there is a positive constant *C* for which

$$|\widehat{\sigma_P}(\xi)| = |\widehat{\sigma_P}(0, \xi_2)| \le \frac{C}{|\xi_2|^{1/2}}.$$

for all $\xi = (0, \xi_2)$ where $\xi_2 \in \mathbb{R} \setminus \{0\}$. In the case that $\xi_2 = 0$, we can evaluate the integrals, and no estimates are needed. We see that this estimate is sharper than in estimate in (4.4).

4.2.2 Decay along the ξ_1 -axis

When $\xi = (\xi_1, 0)$ we have

$$2\pi\widehat{\sigma_P}(\xi) = \int_{-a}^a e^{-i\xi_1\sqrt{1-u^4}} \frac{1}{2\sqrt{1-u^4}} du + \int_{-a}^a e^{+i\xi_1\sqrt{1-u^4}} \frac{1}{2\sqrt{1-u^4}} du + \int_{-b}^b e^{-i\xi_1 u} \frac{1}{4(1-u^2)^{3/4}} du + \int_{-b}^b e^{+i\xi_1 u} \frac{1}{4(1-u^2)^{3/4}} du$$

$$\coloneqq I_1(\xi_1) + I_2(\xi_1) + I_3(\xi_1) + I_4(\xi_1).$$

In view of Lemma 4.1.6, we have

$$|I_3(\xi_1)| \le \frac{C}{|\xi_1|}$$
 and $|I_4(\xi_1)| \le \frac{C}{|\xi_1|}$.

for some positive constant C. To estimate $|I_1(\xi_2)|$ and $|I_2(\xi_2)|$, we first let $\Phi(u) = \sqrt{1 - u^4}$. We have that $\Phi'(u) = 0$ if and only if u = 0, and that the smallest integer k for which $\Phi^{(k)}(u) \neq 0$ for all $u \in [-b,b]$ is k=4. Another appeal to Lemma 4.1.6 yields

$$|I_1(\xi_1)| \le rac{C}{|\xi_1|^{1/4}} \quad ext{and} \quad |I_2(\xi_1)| \le rac{C}{|\xi_1|^{1/4}}.$$

It follows that there is a positive constant *C* for which

$$|\widehat{\sigma_P}(\xi)| = |\widehat{\sigma_P}(\xi_1, 0)| \le \frac{C}{|\xi_1|^{1/4}}.$$

for all $\xi = (\xi_1, 0)$ where $\xi_1 \in \mathbb{R}$. This estimate is consistent with (4.4).

Chapter 5

Application: Estimating convolution powers of complex-valued functions on \mathbb{Z}^d

Chapter 6

Appendix

Continuous one-parameter subgroups of $Gl(\mathbb{R}^d)$ 6.1

Proposition 6.1.1 (See Section 8 of [13]). Let $E, G \in \text{End}(\mathbb{R}^d)$ and $A \in \text{Gl}(\mathbb{R}^d)$. Also, let E^* denote the adjoint of E. Then, for all t, s > 0, the following statements hold:

$$\bullet 1^E = I$$

$$\bullet \ t^{E^*} = (t^E)^*$$

•
$$t^{E^*} = (t^E)^*$$
 • $(t^E)^{-1} = t^{-E}$

$$\bullet At^E A^{-1} = t^{AEA^{-1}}$$

•
$$\det(t^E) = t^{\operatorname{tr} E}$$
 • $(st)^E = s^E t^E$

$$(st)^E = s^E t^E$$

• If
$$EG = GE$$
, then $t^E t^G = t^{E+G}$

Definition 6.1.2. A continuous one-parameter group $T_t = t^E$ is said to be contracting if

$$\lim_{t\to 0}||T_t||=0.$$

Proposition 6.1.3. Let $\{T_t\}$ be a continuous one-parameter group. Then $\{T_t\}$ is contracting if and only if

$$\lim_{t \to 0} |T_t x| = 0 \tag{6.1}$$

for all $x \in \mathbb{R}^d$.

Proof. It is clear that (6.1) is a necessary condition for $\{T_t\}$ to be contracting. We must therefore prove (6.1) also sufficient. To this end, we assume that the continuous one-parameter group $\{T_t\}$ satisfies (6.1). By virtue of the continuity of $\{T_t\}$ and (6.1), we have

$$\sup_{0 < t \le 1} |T_t x| < \infty$$

for each $x \in \mathbb{R}^d$. From the Banach-Steinhaus theorem, it follows that $||T_t|| \leq C$ for all $0 < t \leq 1$. Now, suppose that $\{T_t\}$ is not contracting. In this case, one can find a sequence $\{\eta_n\}\subseteq\mathbb{S}$ and a sequence $t_n \to 0$ for which $\lim_n |T_{t_n}\eta_n| > 0$. But because the unit sphere is compact, $\{\eta_n\}$ has a convergence subsequence $\eta_{n_k} \to \eta$ with $|\eta| = 1$. Observe that, for all n,

$$|T_{t_n}(\eta - \eta_n)| \le C|\eta - \eta_n|$$

and so it follows that

$$\lim_{k \to \infty} |T_{t_{n_k}} \eta| = \lim_{k \to \infty} |T_{t_{n_k}} \eta_{n_k}| > 0,$$

a contradiction.

Lemma 6.1.4. Let $\{T_t\} \subseteq Gl_d(\mathbb{R})$ be a continuous one-parameter group and let $E \in \operatorname{End}(\mathbb{R}^d)$ be its generator, i.e., $T_t = t^E$ for all t > 0. If $\{T_t\}$ is contracting, then $E \in \operatorname{Gl}(\mathbb{R}^d)$ and there is a positive constant C for which

$$||T_t|| \le C + t^{||E||}$$

for all t > 0.

Proof. If for some non-zero vector η , $E\eta=0$, then $t^E\eta=\eta$ for all t>0 and this would contradict our assumption that $\{T_t\}$ is contracting. Hence $E\in \mathrm{Gl}(\mathbb{R}^d)$ and, in particular, $\|E\|>0$. From the representation $T_t=t^E$, it follows immediately that $\|T_t\|\leq t^{\|E\|}$ for all $t\geq 1$ and so it remains to estimate $\|T_t\|$ for t<1. Given that $\{T_t\}$ is continuous and contracting, the map $t\mapsto \|T_t\|$ is continuous and approaches 0 as $t\to 0$ and so it is necessarily bounded for $0< t\leq 1$.

Lemma 6.1.5. Let $E \in Gl_d(\mathbb{R})$ be diagonalizable with strictly positive spectrum. Then $\{t^E\}$ is a continuous one-parameter contracting group. Moreover, there is a positive constant C such that

$$||t^E|| \le Ct^{\lambda_{max}}$$

for all $t \ge 1$ and

$$||t^E|| \le Ct^{\lambda_{min}}$$

for all 0 < t < 1, where $\lambda_{max} = \max(\operatorname{Spec}(E))$ and $\lambda_{min} = \min(\operatorname{Spec}(E))$.

Proof. Let $A \in \operatorname{Gl}_d(\mathbb{R})$ be such that $A^{-1}EA = D = \operatorname{diag}(\lambda_1, \lambda_2, \dots, \lambda_d)$ where necessarily $\operatorname{Spec}(E) = \operatorname{Spec}(D) = \{\lambda_1, \lambda_2, \dots, \lambda_d\} \subseteq (0, \infty)$. It follows from the spectral mapping theorem that $\operatorname{Spec}(t^D) = \{t^{\lambda_1}, t^{\lambda_2}, \dots, t^{\lambda_d}\}$ for all t > 0 and moreover, because t^D is symmetric,

$$||t^D|| \le \max(\{t^{\lambda_1}, t^{\lambda_2}, \dots, t^{\lambda_d}\}) = \begin{cases} t^{\lambda_{\max}} & \text{if } t \ge 1\\ t^{\lambda_{\min}} & \text{if } t < 1. \end{cases}$$

By the property of t^E , we have

$$\|t^E\| = \|At^DA^{-1}\| \leq \|A\|\|t^D\|\|A^{-1}\| \leq C\|t^D\| = C \times \begin{cases} t^{\lambda_{\max}} & \text{if } t \geq 1 \\ t^{\lambda_{\min}} & \text{if } t < 1 \end{cases}$$

for t > 0 where $C = ||A|| ||A^{-1}||$; in particular, $\{t^E\}$ is contracting because $\lambda_{\min} > 0$.

Proposition 6.1.6. Let $\{T_t\}_{t>0} \subseteq Gl(\mathbb{R}^d)$ be a continuous one-parameter contracting group. Then, for all non-zero $x \in \mathbb{R}^d$,

$$\lim_{t\to 0} |T_t x| = 0 \quad \text{and} \quad \lim_{t\to \infty} |T_t x| = \infty.$$

Proof. The validity of the first limit is clear. Upon noting that $|x| = |T_{1/t}T_tx| \le ||T_{1/t}|||T_tx|$ for all t > 0, the second limit follows at once.

Proposition 6.1.7. Let $\{T_t\}_{t>0}$ be a continuous one-parameter contracting group. There holds the following:

(a) For each non-zero $x \in \mathbb{R}^d$, there exists t > 0 and $\eta \in \mathbb{S}$ for which $T_t \eta = x$. Equivalently,

$$\mathbb{R}^d \setminus \{0\} = \{T_t \eta : t > 0 \text{ and } \eta \in \mathbb{S}\}.$$

(b) For each sequence $\{x_n\} \subseteq \mathbb{R}^d$ such that $\lim_n |x_n| = \infty$, $x_n = T_{t_n} \eta_n$ for each n, where $\{\eta_n\} \subseteq \mathbb{S}$ and $t_n \to \infty$ as $n \to \infty$.

(c) For each sequence $\{x_n\} \subseteq \mathbb{R}^d$ such that $\lim_n |x_n| = 0$, $x_n = T_{t_n} \eta_n$ for each n, where $\{\eta_n\} \subseteq \mathbb{S}$ and $t_n \to 0$ as $n \to \infty$.

Proof. In view of Proposition 6.1.6, the assertion (a) is a straightforward application of the intermediate value theorem. For (b), suppose that $\{x_n\} \subseteq \mathbb{R}^d$ is such that $|x_n| \to \infty$ as $n \to \infty$. In view of (a), take $\{\eta_n\} \subseteq S$ and $\{t_n\} \subseteq (0,\infty)$ for which $x_n = T_{t_n}\eta_n$ for each n. In view of Lemma 6.1.4,

$$\infty = \liminf_{n} |x_n| \le \liminf_{n} (C + t_n^M) |\eta_n| \le C + \liminf_{n} t_n^M,$$

where C, M > 0 and therefore $t_n \to \infty$. If instead $\lim_n x_n = 0$,

$$\infty = \lim_{n \to \infty} \frac{|\eta_n|}{|x_n|} = \lim_{n \to \infty} \frac{|T_{1/t_n} x_n|}{|x_n|} \le \limsup_n ||T_{1/t_n}|| \le \limsup_n (C + (1/t_n)^M)$$

from which we see that $t_n \to 0$, thus proving (c).

Proposition 6.1.8. Let $\{T_t\}$ be a continuous contracting one-parameter group. Then for any open neighborhood $\mathcal{O} \subseteq \mathbb{R}^d$ of the origin and any compact set $K \subseteq \mathbb{R}^d$, $K \subseteq T_t(\mathcal{O})$ for sufficiently large t.

Proof. Assume, to reach a contradiction, that there are sequences $\{x_n\} \subseteq K$ and $t_n \to \infty$ for which $x_n \notin T_{t_n}(\mathcal{O})$ for all n. Because K is compact, $\{x_n\}$ has a subsequential limit and so by relabeling, let us take sequences $\{\zeta_k\} \subseteq K$ and $\{r_k\} \subseteq (0,\infty)$ for which $\zeta_k \to \zeta$, $r_k \to \infty$ and $\zeta_k \notin T_{r_k}(\mathcal{O})$ for all k. Setting $s_k = 1/r_k$ and using the fact that $\{T_t\}$ is a one-parameter group, we have $T_{s_k}\zeta_k \notin \mathcal{O}$ for all k and so $\liminf_k |T_{s_k}\zeta_k| > 0$, where $s_k \to 0$. This is however impossible because $\{T_t\}$ is contracting and so

$$\lim_{k \to \infty} |T_{s_k} \zeta_k| \le \lim_{k \to \infty} |T_{s_k} (\zeta_k - \zeta)| + \lim_{k \to \infty} |T_{s_k} \zeta| \le C \lim_{k \to \infty} |\zeta_k - \zeta| + 0 = 0$$

in view of Lemma 6.1.4.

6.2 Some basic measure theory

Theorem 6.2.1. Every open subset \mathcal{O} of \mathbb{R} can be written uniquely as a countable union of disjoint open intervals. Every open subset \mathcal{O} of \mathbb{R}^d , $d \geq 1$, can be written as a countable union of almost disjoint closed cubes.

Definition 6.2.2 (Exterior measure). *If* $E \subseteq \mathbb{R}^d$, then the exterior measure of E is

$$m_*(E) = \inf \sum_{n=1}^{\infty} |Q_j| \in [0, \infty]$$

where the infimum is taken over all countable coverings $E \subseteq \bigcup_{j=1}^{\infty} Q_j$ by closed cubes.

Proposition 6.2.3. (Monotonicity) If $E_1 \subset E_2$ then $m_*(E_1) \leq m_*(E_2)$.

Proposition 6.2.4. (Countable Sub-additivity) If $E = \bigcup_{j=1}^{\infty} E_j$ then $m_*(E) \leq \sum_{j=1}^{\infty} m_*(E_j)$.

Proposition 6.2.5. *If* $E \subseteq \mathbb{R}^d$, $m_*(E) \leq \inf m_*(\mathcal{O})$ *where the infimum is taken over all open* $\mathcal{O} \supseteq E$.

Proposition 6.2.6. If $E = E_1 \cup E_2$ and $d(E_1, E_2) > 0$ then $m_*(E) = m_*(E_1) + m_*(E_2)$.

Proposition 6.2.7. If E is a countable union of almost disjoint cubes $E = \bigcup_{i=1}^{\infty} Q_i$ then $m_*(E) = \sum_{i=1}^{\infty} |Q_i|$.

Definition 6.2.8 (Lebesgue measurable). A subset E of \mathbb{R}^d is Lebesgue measurable if for any $\epsilon > 0$ there exists an open set \mathcal{O} containing E such that

$$m_*(\mathcal{O} \setminus E) \leq \epsilon$$
,

in which case, the Lebesgue measure of E is given by

$$m(E) = m_*(E).$$

Proposition 6.2.9. *Every open set in* \mathbb{R}^d *is measurable.*

Proposition 6.2.10. *If* $m_*(E) = 0$ *then* E *is measurable. In particular, if* $F \subseteq E$ *with* $m_*(E) = 0$ *then* F *is measurable.*

Proposition 6.2.11. A countable union of measurable sets is measurable.

Proposition 6.2.12. *Closed sets are measurable.*

Proposition 6.2.13. *The complement of a measurable set is measurable.*

Proposition 6.2.14. A countable intersection of measurable sets is measurable.

Theorem 6.2.15. If E_1, E_2, \ldots are disjoint measurable sets, and $E = \bigcup_{i=1}^{\infty} E_i$ then

$$m(E) = \sum_{i=1}^{\infty} m(E_i)$$

Corollary 6.2.16. Suppose E_1, E_2, \ldots are measurable subsets of \mathbb{R}^d ,

- If $E_k \uparrow E$ then $m(E) = \lim_{n \to \infty} m(E_n)$.
- If $E_k \downarrow E$ and $m(E_k) < \infty$ for some k then $m(E) = \lim_{n \to \infty} m(E_n)$.

Theorem 6.2.17. If E is a measurable subset of \mathbb{R}^d then for every $\epsilon > 0$,

- There exists an open set \mathcal{O} with $E \subseteq \mathcal{O}$ and $m(\mathcal{O} \setminus E) \le \epsilon$.
- There exists a closed set F with $F \subseteq E$ and $m(E \setminus F) \le \epsilon$.
- If m(E) is finite, there exists a compact set K with $K \subseteq E$ and $m(E \setminus K) \le \epsilon$.
- If m(E) is finite, then there exists a finite union $F = \bigcup_{i=1}^{N} Q_i$ of closed cubes such that $m(E\Delta F) \le \epsilon$, where the notation $E\Delta F$ stands for the symmetric difference between the sets E and F:

$$E\Delta F = (E \setminus F) \cup (F \setminus E).$$

Corollary 6.2.18. A subset E of \mathbb{R}^d is measurable

• if and only if E differs from a G_{δ} set by a set of measure zero. Here a G_{δ} set is a countable intersection of open sets,

• if and only if E differs from a F_{σ} by a set of measure zero. Here an F_{σ} set is a countable union of closed sets.

Definition 6.2.19 (Algebra in a set). Let X be a set. An algebra in X is a non-empty collection of subsets of X that is closed under complements, finite unions, and finite intersections.

Definition 6.2.20 (σ -algebra). Consider a set X an its power set $\mathcal{P}(X)$. A set $\mathcal{A} \subseteq \mathcal{P}(X)$ is a σ -algebra if

- $\varnothing, X \in \mathcal{A}$
- $A \in \mathcal{A} \implies A^c = X \setminus A \in \mathcal{A}$.
- $A_i \in \mathcal{A}, i \in \mathbb{N} \implies \bigcup_{i=1}^{\infty} A_i \in \mathcal{A}.$

Any set $A \in \mathcal{A}$ is called an \mathcal{A} -measurable set.

Remark 5. If A_i is a σ -algebra on X, then for $i \in I$ where I is any index set, then $\bigcap_{i \in I} A_i$ is also a σ -algebra on X. This is important especially when we want to construct some σ -algebra that has all the properties of other σ -algebras.

Proposition 6.2.21. *For* $\mathcal{M} \subset \mathcal{P}(X)$ *, there is a smallest* σ *-algebra that contains* \mathcal{M} :

$$\sigma(M) := \bigcap_{\mathcal{A} \supset \mathcal{M}} \mathcal{A}, \quad \mathcal{A} \text{ is a } \sigma\text{-algebra}$$

called the σ -algebra generated by \mathcal{M} .

Definition 6.2.22 (Borel σ -algebra). Let X be a topological space (or a metric space, or a subset of \mathbb{R}^n), so that we have "open sets." The Borel σ -algebra $\mathcal{B}(X)$ is the smallest σ -algebra generated by the open sets.

Definition 6.2.23 (Measure(able) space). *Consider a set* X *and a* σ -algebra A *on* X. (X, A) *is a measurable space. The map* $\mu: A \to [0, \infty] = [0, \infty) \cup \{\infty\}$ *is called a measure if it satisfies:*

- $\mu(\varnothing) = 0$
- σ-additivity

$$\mu\left(\bigcup_{i=1}^{\infty} A_i\right) = \sum_{i=1}^{\infty} \mu(A_i)$$

whenever $A_i \cap A_j = \emptyset$ for $i \neq j$ and $A_i \in \mathcal{A}$ for all i.

Once the measure μ is defined on the measurable space (X, A), the triple (X, A, μ) is called a measure space.

Definition 6.2.24 (Complete measure space). *A measure space* (X, Σ, μ) *is a complete measure space if and only if for any null set (set of measure zero)* $N \in \Sigma$ *, if* $S \subseteq N$ *then* $S \in \Sigma$.

Definition 6.2.25 (σ -finite measure). Let (X, \mathcal{A}, μ) be a measure space. The measure μ is called a σ -finite measure if it satisfies one of the following equivalent criteria:

• The set X can be covered with at most countably many measurable sets with finite measure, i.e., there are sets $A_1, \ldots, \in \mathcal{A}$ with $\mu(A_n) < \infty$ for all n such that $\bigcup_{n \in \mathbb{N}} A_n = X$.

- The set X can be covered with at most countable many measurable disjoint sets with finite measure, i.e., there are sets $B_1, \ldots, \in \mathcal{A}$ with $\mu(B_n) < \infty$ for all n and $B_i \cap B_j = \varnothing$ for $i \neq j$ that satisfy $\bigcup_{n \in \mathbb{N}} B_n = X$.
- The set X can be covered with a monotone sequence of measurable sets with finite measure, i.e., there are sets $C_1, \ldots, \in \mathcal{A}$ with $C_1 \subseteq C_2 \subseteq \ldots$ and $\mu(C_n) < \infty$ for all n that satisfy $\bigcup_{n \in \mathbb{N}} C_n = X$.
- There exists a strictly positive measurable function f whose integral is finite, i.e., f(x) > 0 for all $x \in X$ and $\int_X f d\mu < \infty$.

If μ is a σ -finite measure, the measure space (X, \mathcal{A}, μ) is called a σ -finite measure space.

Definition 6.2.26 (Premeasure). Let A be an algebra in X. A premeasure on an algebra A is a function $\mu_0: A \to [0, \infty]$ that satisfies

- $\mu_0(\emptyset) = 0$.
- If E_1, \ldots is a countable collection of disjoint sets in A with $\bigcup_{n=1}^{\infty} E_n \in A$ then

$$\mu_0\left(\bigcup_{k=1}^{\infty} E_k\right) = \sum_{k=1}^{\infty} \mu_0(E_k)$$

In particular, μ_0 *is finitely additive on* \mathcal{A} *.*

Definition 6.2.27 (Monotone class). A monotone class on a space X is a subset $C \subset \mathcal{P}(X)$ that is closed under increasing unions and countable decreasing intersections. That is, if $E_j \in C$ and $E_1 \subset E_2 \subset \ldots$ then $\bigcup E_j \in C$, and likewise for intersections.

Remark 6. Every σ -algebra is a monotone class.

Remark 7. The intersection of any family of monotone classes is a monotone class.

Definition 6.2.28 (Monotone class generated by a subset of $\mathcal{P}(X)$). For any $\mathcal{E} \subset \mathcal{P}(X)$, there is a unique smallest monotone class containing \mathcal{E} , called the monotone class generated by \mathcal{E} .

Definition 6.2.29 (σ -algebra generated by a family of subsets). Let F be an arbitrary family of subsets of X. Then there exists a unique smallest σ -algebra which contains every set in F. It is, in fact, the intersection of all σ -algebras containing F. This σ -algebra is denoted $\sigma(F)$ and is called the σ -algebra generated by F.

Lemma 6.2.30 (Monotone Class Lemma). *If* A *is an algebra of subsets of* X, *then the monotone class* C *generated by* A *coincides with the* σ *-algebra* M *generated by* A.

Theorem 6.2.31 (Product measure). Let (X, \mathcal{M}, μ) and (Y, \mathcal{N}, ν) be measure spaces, then we have the product σ -algebra $\mathcal{M} \otimes \mathcal{N}$ on $X \times Y$. If μ, ν are σ -finite, then $\mu \times \nu$ is also σ -finite. In this case, $\mu \times \nu$ is a unique measure on $\mathcal{M} \otimes \mathcal{N}$ such that

$$\mu \times \nu(A \times B) = \mu(A)\nu(B)$$

for all rectangles $A \times B$ where $A \in \mathcal{M}$ and $B \in \mathcal{N}$.

Definition 6.2.32 (Slices). Let $(X, \mathcal{M}, \mu), (Y, \mathcal{N}, \nu)$ be measure spaces. If $E \subset X \times Y$, for $x \in X, y \in Y$ we define the x-section E_x and y-section E^y of E by

$$E_x = \{ y \in Y : (x, y) \in E \} \subseteq Y, \quad E^y = \{ x \in X : (x, y) \in E \} \subseteq X$$

Also, if f is a function on $X \times Y$ we define the x-section f_x and y-section f^y of f by

$$f_x(y) = f^y(x) = f(x, y).$$

Proposition 6.2.33 (Measurable slice proposition).

- If $E \in \mathcal{M} \otimes \mathcal{N}$, then $E_x \in \mathcal{N}$ for all $x \in X$ and $E^y \in \mathcal{M}$ for all $y \in Y$.
- If f is $M \otimes N$ -measurable, then f_x is N-measurable for all $x \in X$ and f^y is M-measurable for all $y \in Y$

Theorem 6.2.34 (Fubini-Tonelli's Theorem for finite measure spaces). *Suppose that* (X, \mathcal{M}, μ) *and* (Y, \mathcal{N}, ν) *are* σ -finite measure spaces

• (Tonelli's) If $f \in L^+(X \times Y)$ then the functions $g(x) = \int f_x d\nu$ and $h(y) = \int f^y d\mu$ are in $L^+(X), L^+(Y)$, respectively, and

$$\int f d(\mu \times \nu) = \int \left[\int f(x, y) d\nu(y) \right] d\mu(x) = \int \left[\int f(x, y) d\mu(x) \right] d\nu(y).$$

• (Fubini's) If $f \in L^1(\mu \times \nu)$ then $f_x \in L^1(\nu)$ for a.e. $x \in X$, $f^y \in L^1(\mu)$ for a.e. $y \in Y$, the a.e.-defined functions $g(x) = \int f_x d\nu$ and $h(y) = \int f^y d\mu$ are in $L^1(\mu), L^1(\nu)$ respectively, and

$$\int f d(\mu \times \nu) = \int \left[\int f(x, y) d\nu(y) \right] d\mu(x) = \int \left[\int f(x, y) d\mu(x) \right] d\nu(y).$$

Theorem 6.2.35 (Fubini's Theorem for Complete Measures). Suppose that (X, \mathcal{M}, μ) and (Y, \mathcal{N}, ν) are σ -finite and complete measure spaces, and let $(X \times Y, L, \lambda)$ be the completion of $(X \times Y, \mathcal{M} \otimes \mathcal{N}, \mu \times \nu)$. If f is L-measurable and either (a) $f \geq 0$ or (b) $f \in L^1(\lambda)$, then f_x, f^y are also integrable for a.e. x, y. Moreover, $x \mapsto \int f_x d\nu, y \mapsto \int f^y d\mu$ are measurable, and in case (b) also integrable, and

$$\int f d\lambda = \iint f(x, y) d\mu(x) d\nu(y) = \iint f(x, y) d\nu(y) d\mu(x).$$

Theorem 6.2.36 (Monotone Convergence Theorem). Let a measure space (X, \mathcal{A}, μ) and measurable functions $f_n, f: X \to [0, \infty)$ be given for all $n \in \mathbb{N}$ with

- $f_1 \le f_2 \le f_3 \le \dots \mu$ -a.e.
- $\lim_{n\to\infty} f_n(x) = f(x) \mu$ -a.e.

then

$$\lim_{n \to \infty} \int_X f_n \, d\mu = \int \lim_{n \to \infty} f_n \, d\mu = \int_X f \, d\mu.$$

Corollary 6.2.37. Suppose $(g_n)_{n\in\mathbb{N}}$ with $g_n:X\to [0,\infty]$ measurable for all n be given. Then $\sum_{i=1}^\infty g_n:X\to [0,\infty]$ is measurable. By the MCT

$$\int_X \sum_{i=1}^{\infty} g_n d\mu = \sum_{i=1}^{\infty} \int_X g_n d\mu.$$

Corollary 6.2.38. Let (X_j, Σ_j, μ_j) for j = 1, 2 be measure spaces which are isomorphic in the sense that there is a bijection $\varphi : X_1 \to X_2$ for which

$$\Sigma_2 = \{ A \in X_2 : \varphi^{-1}(A) \in \Sigma_1 \}$$

and, for each $A \in \Sigma_2$,

$$\mu_2(A) = \mu_1(\varphi^{-1}(A)).$$

Let $f: X_2 \to \mathbb{C}$. Then f is Σ_2 -measurable if and only if $f \circ \varphi: X_1 \to \mathbb{C}$ is Σ_1 measurable. Further, if $f \geq 0$,

$$\int_{X_2} f \, d\mu_2 = \int_{X_1} f \circ \varphi \, d\mu_1.$$

If f is generally complex-valued, then $f \in L^1(\mu_1)$ if and only if $f \circ \psi \in L^1(\mu_2)$ and in this case we also have

$$\int_{X_2} f \, d\mu_2 = \int_{X_1} f \circ \varphi \, d\mu_1.$$

Theorem 6.2.39 (Lebesgue's Dominated Convergence Theorem). *Let a measure space* (X, \mathcal{A}, μ) *be given. Consider the set of all Lebesgue-integrable functions:*

$$L^1(\mu) \coloneqq \{f: X \to \mathbb{R} \; \text{measureable} | \int_X |f|^1 \, d\mu < \infty \}.$$

For $f \in L^1(\mu)$, write $f = f^+ - f^-$ with $f^+, f^- \ge 0$ and define

$$\int_X f \, d\mu \coloneqq \int_X f^+ \, d\mu - \int_X f^- \, d\mu.$$

Consider $f_n: X \to \mathbb{R}$ a sequence of measurable functions and $f: X \to \mathbb{R}$ with $\lim_{n \to \infty} f_n(x) = f(x)$ for $x \in X$, μ -a.e. where $|f_n| \le g$ with $g \in L^1(\mu)$ for all n. Then $f_1, \ldots, \in L^1(\mu)$, $f \in L^1(\mu)$ and

$$\lim_{n \to \infty} \int_X f_n \, d\mu = \int_X f \, d\mu.$$

6.3 Notational details for Subsection 3.2

Let's write $\partial_{u^1}, \partial_{u^2}, \dots, \partial_{u^{d'}}$ for the frame on S, i.e., each is a vector field on S, and $\frac{\partial}{\partial u^1}, \frac{\partial}{\partial u^2}, \dots, \frac{\partial}{\partial u^{d'}}$ for the frame on $U \in \mathbb{R}^d$. These are connected via the following relation (which actually defines the former): For $f \in C^{\infty}(S)$,

$$\partial_{u^{\mu}}(f)(\eta) = \frac{\partial}{\partial u^{\mu}} \left(f \circ \varphi^{-1} \right) (u)$$

where $u = \varphi(\eta)$. In other words,

$$\partial_{u^{\mu}} = \varphi_*^{-1} \frac{\partial}{\partial u^{\mu}}.$$

Now, by definition of the pullback $(\varphi^{-1})^* d\sigma_P$ is a d'-form on the d'-dimensional manifold U and every such form is given by $f(u)du^1 \wedge du^2 \wedge \cdots \wedge du^{d'}$, i.e., there is a smooth function $f: \mathbb{R}^d \to \mathbb{R}$ for which

$$(\varphi^{-1})^* d\sigma_P = f du^1 \wedge du^2 \wedge \dots \wedge du^{d'}$$

Then, by definition,

$$\int_{U} (\varphi^{-1})^* d(g \cdot d\sigma_P) = \int_{U} g(\varphi^{-1}(u)) f(u) du^1, du^2, \dots, du^{d'}$$

(we've erased the wedges to integrate). Thus, we really just need to know f. We have

$$f(u) = f(u) \left(du^{1} \wedge du^{2} \wedge \dots \wedge du^{d'} \right) \left(\frac{\partial}{\partial u^{1}}, \frac{\partial}{\partial u^{2}}, \dots, \frac{\partial}{\partial u^{d'}} \right)$$

$$= \left((\varphi^{-1})^{*} d\sigma_{P} \right) \left(\frac{\partial}{\partial u^{1}}, \frac{\partial}{\partial u^{2}}, \dots, \frac{\partial}{\partial u^{d'}} \right)$$

$$= d\sigma_{P} \left(\varphi_{*}^{-1} \left(\frac{\partial}{\partial u^{1}} \right), \varphi_{*}^{-1} \left(\frac{\partial}{\partial u^{2}} \right), \dots, \varphi_{*}^{-1} \left(\frac{\partial}{\partial u^{d'}} \right) \right)$$

$$= d\sigma_{P}(\partial_{u^{1}}, \partial_{u^{2}}, \dots, \partial_{u^{d'}})$$

$$= h_{E,\varphi}(u).$$

Therefore

$$\int_U (\varphi^{-1})^*(g \cdot d\sigma_P) = \int_U g(\varphi^{-1}(u)) h_{E,\varphi}(u) du.$$

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