

On the Hausdorff volume in sub-Riemannian geometry

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Abstract

For a regular sub-Riemannian manifold we study the Radon-Nikodym derivative of the spherical Hausdorff measure with respect to a smooth volume. We prove that this is the volume of the unit ball in the nilpotent approximation and it is always a continuous function. We then prove that up to dimension 4 it is smooth, while starting from dimension 5, in corank 1 case, it is \mathcal{C}^3 (and \mathcal{C}^4 on every curve) but in general not \mathcal{C}^5 . These results answer to a question addressed by Montgomery about the relation between two intrinsic volumes that can be defined in a sub-Riemannian manifold, namely the Popp and the Hausdorff volume. If the nilpotent approximation depends on the point (that may happen starting from dimension 5), then they are not proportional, in general.

1 Introduction

In this paper, by a sub-Riemannian manifold we mean a triple $\mathbf{S} = (M, \Delta, \mathbf{g})$, where M is a connected orientable smooth manifold of dimension n , Δ is a smooth vector distribution of constant rank $k < n$, satisfying the Hörmander condition and \mathbf{g} is an Euclidean structure on Δ .

A sub-Riemannian manifold has a natural structure of metric space, where the distance is the so called Carnot-Carathéodory distance

$$d(q_0, q_1) = \inf \left\{ \int_0^T \sqrt{\mathbf{g}_{\gamma(t)}(\dot{\gamma}(t), \dot{\gamma}(t))} dt \mid \gamma : [0, T] \rightarrow M \text{ is a Lipschitz curve,} \right. \\ \left. \gamma(0) = q_0, \gamma(T) = q_1, \dot{\gamma}(t) \in \Delta_{\gamma(t)} \text{ a.e. in } [0, T] \right\}.$$

As a consequence of the Hörmander condition this distance is always finite and continuous, and induces on M the original topology (see Chow-Rashevsky theorem).

Since (M, d) is a metric space, for every $\alpha > 0$ one can define the α -dimensional Hausdorff measure on M , and compute the Hausdorff dimension of M .

Define $\Delta^1 := \Delta$, $\Delta^{i+1} := \Delta^i + [\Delta^i, \Delta]$, for every $i = 1, 2, \dots$. Under the hypothesis that the sub-Riemannian manifold is regular, i.e. if the dimension of Δ^i , $i = 1, \dots, m$ do not depend on

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the point, the Hörmander condition guarantees that there exists (a minimal) $m \in \mathbb{N}$, called *step* of the structure, such that $\Delta_q^m = T_q M$, for all $q \in M$. The sequence

$$\mathcal{G}(\mathbf{S}) := (\underbrace{\dim \Delta}_k, \dim \Delta^2, \dots, \underbrace{\dim \Delta^m}_n)$$

is called *growth vector* of the structure.

In this case, the graded vector space $\text{gr}_q(\Delta)$ associated to the filtration $\Delta_q \subset \Delta_q^2 \subset \dots \subset \Delta_q^m = T_q M$ is well defined, namely

$$\text{gr}_q(\Delta) = \bigoplus_{i=1}^m \Delta_q^i / \Delta_q^{i-1}, \quad \text{where } \Delta_q^0 = 0.$$

Moreover, it is well known that the Hausdorff dimension of M is given by the formula (see [34])

$$Q = \sum_{i=1}^m i k_i, \quad k_i := \dim \Delta_q^i / \Delta_q^{i-1}.$$

In particular the Hausdorff dimension is always bigger than the topological dimension of M .

Moreover, the Q -dimensional Hausdorff measure (denoted \mathcal{H}^Q in the following) behaves like a volume. More precisely, in [34] Mitchell proved that if μ is a smooth volume¹ on M , then $d\mu = f_{\mu\mathcal{H}} d\mathcal{H}^Q$, where $f_{\mu\mathcal{H}}$ is a positive measurable function that is locally bounded and locally bounded away from zero, that is the Radon-Nikodym derivative of μ with respect to \mathcal{H}^Q . According to Mitchell terminology, this means that the two measures are *commensurable* one with respect to the other.

Notice that Hausdorff measure on sub-Riemannian manifolds has been intensively studied, see for instance [26, 34]. A deep study of Hausdorff measure for hypersurfaces in sub-Riemannian geometry, in particular in the context of Carnot groups, can be found in [8, 10, 11, 14, 19, 22, 32, 36] and references therein. Hausdorff measure of curves in sub-Riemannian manifolds is also studied in the problem of motion planning and complexity, see [23, 24, 25, 28].

Let us recall that there are two common non-equivalent definitions of Hausdorff measure. The standard Hausdorff measure, where arbitrary coverings can be used, and the spherical Hausdorff measure, where only ball-coverings appear (see Definition 19). However it is well known that, if \mathcal{S}^Q denotes the spherical Q -Hausdorff measure, then \mathcal{H}^Q is commensurable with \mathcal{S}^Q .² As a consequence, \mathcal{S}^Q is commensurable with μ , i.e.

$$d\mu = f_{\mu\mathcal{S}} d\mathcal{S}^Q,$$

for a positive measurable function $f_{\mu\mathcal{S}}$ that is locally bounded and locally bounded away from zero. In this paper, we are interested to the properties of the function $f_{\mu\mathcal{S}}$. In particular, we would like to get informations about its regularity.

The reason why we study the spherical Hausdorff measure and not the standard Hausdorff measure is that the first one appears to be more natural. Indeed, as explained later, $f_{\mu\mathcal{S}}$ is determined by

¹In the following by a smooth volume on M we mean a measure μ associated to a smooth non-vanishing n -form $\omega_\mu \in \Lambda^n M$, i.e. for every measurable subset $A \subset M$ we set $\mu(A) = \int_A \omega_\mu$.

²Indeed they are absolutely continuous one with respect to the other. In particular, for every $\alpha > 0$, we have $2^{-\alpha} \mathcal{S}^\alpha \leq \mathcal{H}^\alpha \leq \mathcal{S}^\alpha$ (see for instance [21]).

the volume of the unit sub-Riemannian ball of the nilpotent approximation of the sub-Riemannian manifold, that can be explicitly described in a certain number of cases (see Theorem 1 below). On the other hand nothing is known on how to compute $f_{\mu\mathcal{H}}$. We conjecture that $f_{\mu\mathcal{H}}$ is given by the μ -volume of certain isodiametric sets, i.e. the maximum of the μ -volume among all sets of diameter 1 in the nilpotent approximation (see [31, 38] and reference therein for a discussion on isodiametric sets). This quantity is not very natural in sub-Riemannian geometry and is extremely difficult to compute.

Our interests in studying $f_{\mu\mathcal{S}}$ comes from the following question:

Q1 How can we define an intrinsic volume in a sub-Riemannian manifold?

Here by intrinsic we mean a volume which depends neither on the choice of the coordinate system, nor on the choice of the orthonormal frame, but only on the sub-Riemannian structure.

This question was first pointed out by Brockett, see [17], and by Montgomery in his book [35]. Having a volume that depends only on the geometric structure is interesting by itself, however, it is also necessary to define intrinsically a Laplacian in a sub-Riemannian manifold. We recall that the Laplacian is defined as the divergence of the gradient and the definition of the divergence needs a volume since it measures how much the flow of a vector field increases or decreases the volume.

Before talking about the question **Q1** in sub-Riemannian geometry, let us briefly discuss it in the Riemannian case. In a n -dimensional Riemannian manifold there are three common ways of defining an invariant volume. The first is defined through the Riemannian structure and it is the so called Riemannian volume, which in coordinates has the expression $\sqrt{g} dx^1 \dots dx^n$, where g is the determinant of the metric. The second and the third ones are defined via the Riemannian distance and are the n -dimensional Hausdorff measure and the n -dimensional spherical Hausdorff measure. These three volumes are indeed proportional (the constant of proportionality depending on the normalization, see e.g. [20, 21]).

As we said, a regular sub-Riemannian manifold is a metric space, hence it is possible to define the Hausdorff volume \mathcal{H}^Q and the spherical Hausdorff volume \mathcal{S}^Q . Also, there is an equivalent of the Riemannian volume, the so called Popp's volume \mathcal{P} , introduced by Montgomery in his book [35] (see also [2]). The Popp volume is a smooth volume and was used in [2] to define intrinsically the Laplacian (indeed a sub-Laplacian) in sub-Riemannian geometry.

In his book, Montgomery proposed to study whether these invariant volumes are proportional as it occurs in Riemannian geometry. More precisely, he addressed the following question:

Q2 Is Popp's measure equal to a constant multiple (perhaps depending on the growth vector) of the Hausdorff measure?

Montgomery noted that the answer to this question is positive for left-invariant sub-Riemannian structures on Lie groups, since the Hausdorff (both the standard and spherical one) and the Popp volumes are left-invariant and hence proportional to the left Haar measure. But this question is nontrivial when there is no group structure.

One of the main purpose of our analysis is to answer to question **Q2** for the spherical Hausdorff measure, i.e. to the question if the function $f_{\mathcal{P}\mathcal{S}}$ (defined by $d\mathcal{P} = f_{\mathcal{P}\mathcal{S}} d\mathcal{S}^Q$) is constant or not. More precisely, we get a positive answer for regular sub-Riemannian manifolds of dimension 3 and 4, while a negative answer starting from dimension 5, in general.

Once a negative answer to **Q2** is given, it is natural to ask

Q3 What is the regularity of $f_{\mathcal{P}\mathcal{S}}$?

This question is important since the definition of an intrinsic Laplacian via \mathcal{S}^Q require $f_{\mathcal{P}\mathcal{S}}$ to be at least \mathcal{C}^1 .

Notice that since the Popp measure is a smooth volume, then $f_{\mu\mathcal{S}}$ is $\mathcal{C}^k, k = 0, 1, \dots, \infty$ if and only if $f_{\mathcal{P}\mathcal{S}}$ is as well.

We prove that $f_{\mu\mathcal{S}}$ is a continuous function and that for $n \leq 4$ it is smooth. In dimension 5 it is \mathcal{C}^3 but not smooth, in general. Moreover, we prove that the same result holds all corank 1 cases (see Section 2 for a precise definition).

Our main tool is the nilpotent approximation (or the symbol) of the sub-Riemannian structure. Recall that, under the regularity hypothesis, the sub-Riemannian structure $\mathbf{S} = (M, \Delta, \mathbf{g})$ induces a structure of nilpotent Lie algebra on $\text{gr}_q(\Delta)$. The nilpotent approximation at q is the nilpotent simply connected Lie group $\text{Gr}_q(\Delta)$ generated by this Lie algebra, endowed with a suitable left-invariant sub-Riemannian structure $\widehat{\mathbf{S}}_q$ induced by \mathbf{S} , as explained in Section 3.

Recall that there exists a canonical isomorphism of 1-dimensional vector spaces (see [2] for details)

$$\bigwedge^n (T_q^* M) \simeq \bigwedge^n (\text{gr}_q(\Delta)^*). \quad (1)$$

Given a smooth volume μ on M , we define the induced volume $\widehat{\mu}_q$ on the nilpotent approximation at point q as the left-invariant volume on $\text{Gr}_q(\Delta)$ canonically associated to $\omega_\mu(q) \in \wedge^n (T_q^* M)$ by the above isomorphism.

The first result concerns an explicit formula for $f_{\mu\mathcal{S}}$.

Theorem 1. *Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular sub-Riemannian manifold. Let μ a volume on M and $\widehat{\mu}_q$ the induced volume on the nilpotent approximation at point $q \in M$. If A is an open subset of M , then*

$$\mu(A) = \frac{1}{2^Q} \int_A \widehat{\mu}_q(\widehat{B}_q) d\mathcal{S}^Q,$$

where \widehat{B}_q is the unit ball in the nilpotent approximation at point q , i.e.

$$f_{\mu\mathcal{S}}(q) = \frac{1}{2^Q} \widehat{\mu}_q(\widehat{B}_q).$$

Starting from this formula we prove our first result about regularity of the density:

Corollary 2. *Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular sub-Riemannian manifold and let μ be a smooth volume on M . Then the density $f_{\mu\mathcal{S}}$ is a continuous function.*

Theorem 1, specified for the Popp measure \mathcal{P} , permits to answer the Montgomery's question. Indeed, the measure $\widehat{\mathcal{P}}_q$ induced by \mathcal{P} on the nilpotent approximation at point q coincides with the Popp measure built on $\widehat{\mathbf{S}}_q$, as a sub-Riemannian structure. In other words, if we denote $\mathcal{P}_{\widehat{q}}$ the Popp measure on $\widehat{\mathbf{S}}_q$, we get

$$\widehat{\mathcal{P}}_q = \mathcal{P}_{\widehat{q}}. \quad (2)$$

Hence, if the nilpotent approximation does not depend on the point, then $f_{\mathcal{P}\mathcal{S}}$ is constant. In other words we have the following corollary.

Corollary 3. *Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular sub-Riemannian manifold and $\widehat{\mathbf{S}}_q$ its nilpotent approximation at point $q \in M$. If $\widehat{\mathbf{S}}_{q_1}$ is isometric to $\widehat{\mathbf{S}}_{q_2}$ for any $q_1, q_2 \in M$, then $f_{\mathcal{P}\mathbf{S}}$ is constant. In particular this happens if the sub-Riemannian structure is free.*

For the definition of free structure see [35].

Notice that, in the Riemannian case, nilpotent approximations at different points are isometric, hence the Hausdorff measure is proportional to the Riemannian volume (see [20, 21]).

When the nilpotent approximation contains parameters that are function of the point, then, in general, $f_{\mathcal{P}\mathbf{S}}$ is not constant. We have analyzed in details all growth vectors in dimension less or equal than 5:

- dimension 3: (2,3),
- dimension 4: (2,3,4), (3,4),
- dimension 5: (2,3,5), (3,5), (4,5) and the non generic cases (2,3,4,5), (3,4,5).

In all cases the nilpotent approximation is unique, except for the (4,5) case. As a consequence, we get:

Theorem 4. *Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular sub-Riemannian manifold of dimension $n \leq 5$. Let μ be a smooth volume on M and \mathcal{P} be the Popp measure. Then*

- (i) *if $\mathcal{G}(\mathbf{S}) \neq (4, 5)$, then $f_{\mathcal{P}\mathbf{S}}$ is constant. As a consequence $f_{\mu\mathbf{S}}$ is smooth.*
- (ii) *if $\mathcal{G}(\mathbf{S}) = (4, 5)$, then $f_{\mu\mathbf{S}}$ is \mathcal{C}^3 (and \mathcal{C}^4 on smooth curves) but not \mathcal{C}^5 , in general.*

Actually the regularity result obtained in the (4,5) case holds for all corank 1 structures, as specified by the following theorem.

Theorem 5. *Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular corank 1 sub-Riemannian manifold of dimension $n \geq 5$. Let μ be a smooth volume on M . Then $f_{\mu\mathbf{S}}$ is \mathcal{C}^3 (and \mathcal{C}^4 on smooth curves) but not \mathcal{C}^5 , in general.*

Recall that for a corank 1 structure one has $\mathcal{G}(\mathbf{S}) = (n - 1, n)$ (see also Section 2).

Notice that Theorem 5 apply in particular for the Popp measure. The loss of regularity of $f_{\mu\mathbf{S}}$ is due to the presence of what are called *resonance points*. More precisely, the parameters appearing in the nilpotent approximation are the eigenvalues of a certain skew-symmetric matrix which depends on the point. Resonances are the points in which these eigenvalues are multiple.

To prove Theorem 5, we have computed explicitly the optimal synthesis (i.e. all curves that minimize distance starting from one point) of the nilpotent approximation and, as a consequence, the volume of nilpotent balls \widehat{B}_q .

Another byproduct of our analysis is

Proposition 6. *Under the hypothesis of Theorem 5, if there are no resonance points then $f_{\mu\mathbf{S}}$ is smooth.*

The structure of the paper is the following. In Section 2 we recall basic facts about sub-Riemannian geometry and about Hausdorff measures. In Section 3 we provide normal forms for nilpotent structures in dimension less or equal than 5. In Section 4 we prove Theorem 1 and its corollaries, while in Section 5 we study the differentiability of the density for the corank 1 case. In the last Section we prove Theorem 4.

2 Basic Definitions

2.1 Sub-Riemannian manifolds

We start recalling the definition of sub-Riemannian manifold.

Definition 7. A *sub-Riemannian manifold* is a triple $\mathbf{S} = (M, \Delta, \mathbf{g})$, where

- (i) M is a connected orientable smooth manifold of dimension $n \geq 3$;
- (ii) Δ is a smooth distribution of constant rank $k < n$ satisfying the *Hörmander condition*, i.e. a smooth map that associates to $q \in M$ a k -dimensional subspace Δ_q of T_qM and we have

$$\text{span}\{[X_1, [\dots [X_{j-1}, X_j]]](q) \mid X_i \in \overline{\Delta}, j \in \mathbb{N}\} = T_qM, \quad \forall q \in M, \quad (3)$$

where $\overline{\Delta}$ denotes the set of *horizontal smooth vector fields* on M , i.e.

$$\overline{\Delta} = \{X \in \text{Vec}(M) \mid X(q) \in \Delta_q \quad \forall q \in M\}.$$

- (iii) \mathbf{g}_q is a Riemannian metric on Δ_q which is smooth as function of q . We denote the norm of a vector $v \in \Delta_q$ with $|v|$, i.e. $|v| = \sqrt{\mathbf{g}_q(v, v)}$.

A Lipschitz continuous curve $\gamma : [0, T] \rightarrow M$ is said to be *horizontal* (or *admissible*) if

$$\dot{\gamma}(t) \in \Delta_{\gamma(t)} \quad \text{for a.e. } t \in [0, T].$$

Given an horizontal curve $\gamma : [0, T] \rightarrow M$, the *length of γ* is

$$l(\gamma) = \int_0^T |\dot{\gamma}(t)| dt. \quad (4)$$

The *distance* induced by the sub-Riemannian structure on M is the function

$$d(q_0, q_1) = \inf\{l(\gamma) \mid \gamma(0) = q_0, \gamma(T) = q_1, \gamma \text{ horizontal}\}. \quad (5)$$

The hypothesis of connectedness of M and the Hörmander condition guarantees the finiteness and the continuity of $d(\cdot, \cdot)$ with respect to the topology of M (Chow-Rashevsky theorem, see, for instance, [6]). The function $d(\cdot, \cdot)$ is called the *Carnot-Carathéodory distance* and gives to M the structure of metric space (see [12, 26]).

Remark 8. It is a standard fact that $l(\gamma)$ is invariant under reparameterization of the curve γ . Moreover, if an admissible curve γ minimizes the so-called *action functional*

$$J(\gamma) := \frac{1}{2} \int_0^T |\dot{\gamma}(t)|^2 dt.$$

with T fixed (and fixed initial and final point), then $|\dot{\gamma}(t)|$ is constant and γ is also a minimizer of $l(\cdot)$. On the other side, a minimizer γ of $l(\cdot)$ such that $|\dot{\gamma}(t)|$ is constant is a minimizer of $J(\cdot)$ with $T = l(\gamma)/v$.

Locally, the pair (Δ, \mathbf{g}) can be given by assigning a set of k smooth vector fields spanning Δ and that are orthonormal for \mathbf{g} , i.e.

$$\Delta_q = \text{span}\{X_1(q), \dots, X_k(q)\}, \quad \mathbf{g}_q(X_i(q), X_j(q)) = \delta_{ij}. \quad (6)$$

In this case, the set $\{X_1, \dots, X_k\}$ is called a *local orthonormal frame* for the sub-Riemannian structure.

Definition 9. Let Δ be a distribution. Its *flag* is the sequence of distributions $\Delta^1 \subset \Delta^2 \subset \dots$ defined through the recursive formula

$$\Delta^1 := \Delta, \quad \Delta^{i+1} := \Delta^i + [\Delta^i, \Delta].$$

A sub-Riemannian manifold is said to be *regular* if for each $i = 1, 2, \dots$ the dimension of $\Delta_{q_0}^i$ does not depend on the point $q_0 \in M$.

Remark 10. In this paper we always deal with regular sub-Riemannian manifolds. In this case Hörmander condition can be rewritten as follows:

$$\exists \text{ minimal } m \in \mathbb{N} \quad \text{such that} \quad \Delta_q^m = T_q M, \quad \forall q \in M.$$

The sequence $\mathcal{G}(\mathbf{S}) := (\dim \Delta, \dim \Delta^2, \dots, \dim \Delta^m)$ is called *growth vector*. Under the regularity assumption $\mathcal{G}(\mathbf{S})$ does not depend on the point and m is said the *step* of the structure. The minimal growth is $(k, k+1, k+2, \dots, n)$. When the growth is maximal the sub-Riemannian structure is called *free* (see [35]).

A sub-Riemannian manifold is said to be *corank 1* if its growth vector satisfies $\mathcal{G}(\mathbf{S}) = (n-1, n)$. A sub-Riemannian manifold \mathbf{S} of odd dimension is said to be *contact* if $\Delta = \ker \omega$, where $\omega \in \Lambda^1 M$ and $d\omega|_{\Delta}$ is non degenerate. A sub-Riemannian manifold M of even dimension is said to be *quasi-contact* if $\Delta = \ker \omega$, where $\omega \in \Lambda^1 M$ and satisfies $\dim \ker d\omega|_{\Delta} = 1$.

Notice that contact and quasi-contact structures are regular and corank 1.

A sub-Riemannian manifold is said to be *nilpotent* if there exists an orthonormal frame for the structure $\{X_1, \dots, X_k\}$ and $j \in \mathbb{N}$ such that $[X_{i_1}, [X_{i_2}, \dots, [X_{i_{j-1}}, X_{i_j}]]] = 0$ for every commutator of length j .

2.2 Geodesics

In this section we briefly recall some facts about sub-Riemannian geodesics. In particular, we define the sub-Riemannian Hamiltonian.

Definition 11. A *geodesic* for a sub-Riemannian manifold $\mathbf{S} = (M, \Delta, \mathbf{g})$ is a curve $\gamma : [0, T] \rightarrow M$ such that for every sufficiently small interval $[t_1, t_2] \subset [0, T]$, the restriction $\gamma|_{[t_1, t_2]}$ is a minimizer of $J(\cdot)$. A geodesic for which $\mathbf{g}_{\gamma(t)}(\dot{\gamma}(t), \dot{\gamma}(t))$ is (constantly) equal to one is said to be parameterized by arclength.

Let us consider the cotangent bundle T^*M with the canonical projection $\pi : T^*M \rightarrow M$, and denote the standard pairing between vectors and covectors with $\langle \cdot, \cdot \rangle$. The Liouville 1-form $s \in \Lambda^1(T^*M)$ is defined as follows: $s_\lambda = \lambda \circ \pi_*$, for every $\lambda \in T^*M$. The canonical symplectic structure on T^*M is defined by the closed 2-form $\sigma = ds$. In canonical coordinates (ξ, x)

$$s = \sum_{i=1}^n \xi_i dx_i, \quad \sigma = \sum_{i=1}^n d\xi_i \wedge dx_i.$$

We denote the Hamiltonian vector field associated to a function $h \in C^\infty(T^*M)$ with \vec{h} . Namely we have $dh = \sigma(\cdot, \vec{h})$ and in coordinates we have

$$\vec{h} = \sum_i \frac{\partial h}{\partial \xi_i} \frac{\partial}{\partial x_i} - \frac{\partial h}{\partial x_i} \frac{\partial}{\partial \xi_i}$$

The sub-Riemannian structure defines an Euclidean norm $|\cdot|$ on the distribution $\Delta_q \subset T_qM$. As a matter of fact this induces a dual norm

$$\|\lambda\| = \max_{\substack{v \in \Delta_q \\ |v|=1}} \langle \lambda, v \rangle, \quad \lambda \in T_q^*M,$$

which is well defined on $\Delta_q^* \simeq T_q^*M/\Delta_q^\perp$, where $\Delta_q^\perp = \{\lambda \in T_q^*M \mid \langle \lambda, v \rangle = 0, \forall v \in \Delta_q\}$ is the annihilator of the distribution.

The *sub-Riemannian Hamiltonian* is the smooth function on T^*M , which is quadratic on fibers, defined by

$$H(\lambda) = \frac{1}{2} \|\lambda\|^2, \quad \lambda \in T^*M.$$

If $\{X_1, \dots, X_k\}$ is a local orthonormal frame for the sub-Riemannian structure it is easy to see that

$$H(\lambda) = \frac{1}{2} \sum_{i=1}^k \langle \lambda, X_i(q) \rangle^2, \quad \lambda \in T^*M, \quad q = \pi(\lambda).$$

Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a sub-Riemannian manifold and fix $q_0 \in M$. We define the *endpoint map* (at time 1) as

$$F : \mathcal{U} \rightarrow M, \quad F(\gamma) = \gamma(1),$$

where \mathcal{U} denotes the set of admissible trajectories starting from q_0 and defined in $[0, 1]$. If we fix a point $q_1 \in M$, the problem of finding shortest paths from q_0 to q_1 is equivalent to the following one

$$\min_{F^{-1}(q_1)} J(\gamma), \tag{7}$$

where J is the action functional (see Remark 8). Then Lagrange multipliers rule implies that any $\gamma \in \mathcal{U}$ solution of (7) satisfies one of the following equations

$$\lambda_1 D_\gamma F = d_\gamma J, \tag{8}$$

$$\lambda_1 D_\gamma F = 0, \tag{9}$$

for some nonzero covector $\lambda_1 \in T_{\gamma(1)}^*M$ associated to γ . The following characterization is a corollary of Pontryagin Maximum Principle (PMP for short, see for instance [6, 15, 29, 37]):

Theorem 12. *Let γ be a minimizer. A nonzero covector λ_1 satisfies (8) or (9) if and only if there exists a Lipschitz curve $\lambda(t) \in T_{\gamma(t)}^*M$, $t \in [0, 1]$, such that $\lambda(1) = \lambda_1$ and*

- if (8) holds, then $\lambda(t)$ is a solution of $\dot{\lambda}(t) = \overrightarrow{H}(\lambda(t))$ for a.e. $t \in [0, 1]$,
- if (9) holds, then $\lambda(t)$ satisfies $\sigma(\dot{\lambda}(t), T_{\lambda(t)}\Delta^\perp) = 0$ for a.e. $t \in [0, 1]$.

The curve $\lambda(t)$ is said to be an extremal associated to $\gamma(t)$. In the first case $\lambda(t)$ is called a normal extremal while in the second one an abnormal extremal.

Remark 13. It is possible to give a unified characterization of normal and abnormal extremals in terms of the symplectic form. Indeed the Hamiltonian H is always constant on extremals, hence $\lambda(t) \subset H^{-1}(c)$ for some $c \geq 0$. Theorem 12 can be rephrased as follows: any extremal $\lambda(t)$ such that $H(\lambda(t)) = c$ is a reparametrization of a characteristic curve of the differential form $\sigma|_{H^{-1}(c)}$, where $c = 0$ for abnormal extremals, and $c > 0$ for normal ones.

Also notice that, if $\lambda(t)$ is a normal extremal, then, for every $\alpha > 0$, $\lambda_\alpha(t) := \alpha \lambda(\alpha t)$ is also a normal extremal. If the curve is parametrized in such a way that $H = \frac{1}{2}$ then we say that the extremal is arclength parameterized. Trajectories parametrized by arclength corresponds to initial covectors λ_0 belonging to the hypercylinder $\Lambda_{q_0} := T_{q_0}^*M \cap H^{-1}(\frac{1}{2}) \simeq S^{k-1} \times \mathbb{R}^{n-k}$ in $T_{q_0}^*M$.

Remark 14. From Theorem 12 it follows that $\lambda(t) = e^{t\vec{H}}(\lambda_0)$ is the normal extremal with initial covector $\lambda_0 \in \Lambda_{q_0}$. If $\pi : T^*M \rightarrow M$ denotes the canonical projection, then it is well known that $\gamma(t) = \pi(\lambda(t))$ is a geodesic (starting from q_0). On the other hand, in every 2-step sub-Riemannian manifold all geodesics are projection of normal extremals, since there is no strict abnormal minimizer (see Goh conditions, [6]).

The following proposition resumes some basic properties of small sub-Riemannian balls

Proposition 15. *Let \mathbf{S} be a sub-Riemannian manifold and $B_{q_0}(\varepsilon)$ the sub-Riemannian ball of radius ε at fixed point $q_0 \in M$. For $\varepsilon > 0$ small enough we have:*

- (i) $\forall q \in B_{q_0}(\varepsilon)$ there exists a minimizer that join q and q_0 ,
- (ii) $\text{diam}(B_{q_0}(\varepsilon)) = 2\varepsilon$.

Claim (i) is a consequence of Filippov theorem (see [6, 16]). To prove (ii) it is sufficient to show that, for ε small enough, there exists two points in $q_1, q_2 \in \partial B_{q_0}(\varepsilon)$ such that $d(q_1, q_2) = 2\varepsilon$.

To this purpose, consider the projection $\gamma(t) = \pi(\lambda(t))$ of a normal extremal starting from $\gamma(0) = q_0$, and defined in a small neighborhood of zero $t \in]-\delta, \delta[$. Using arguments of Chapter 17 of [6] one can prove that $\gamma(t)$ is globally minimizer. Hence if we consider $0 < \varepsilon < \delta$ we have that $q_1 = \gamma(-\varepsilon)$ and $q_2 = \gamma(\varepsilon)$ satisfy the property required, which proves claim (ii).

Definition 16. Fix $q_0 \in M$. We define the *Exponential map* starting from q_0 as

$$\text{Exp}_{q_0} : T_{q_0}^*M \rightarrow M, \quad \text{Exp}_{q_0}(\lambda_0) = \pi(e^{\vec{H}}(\lambda_0)).$$

Using the homogeneity property $H(c\lambda) = c^2H(\lambda)$, $\forall c > 0$, we have that

$$e^{\vec{H}}(s\lambda) = e^{s\vec{H}}(\lambda), \quad \forall s > 0.$$

In other words we can recover the geodesic on the manifold with initial covector λ_0 as the image under Exp_{q_0} of the ray $\{t\lambda_0, t \in [0, 1]\} \subset T_{q_0}^*M$ that join the origin to λ_0 .

$$\text{Exp}_{q_0}(t\lambda_0) = \pi(e^{\vec{H}}(t\lambda_0)) = \pi(e^{t\vec{H}}(\lambda_0)) = \pi(\lambda(t)) = \gamma(t).$$

Next, we recall the definition of cut and conjugate time.

Definition 17. Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a sub-Riemannian manifold. Let $q_0 \in M$ and $\lambda_0 \in \Lambda_{q_0}$. Assume that the geodesic $\gamma(t) = \text{Exp}_{q_0}(t\lambda_0)$ for $t > 0$, is not abnormal.

(i) The *first conjugate time* is $t(\lambda_0) = \min\{t > 0, t\lambda_0 \text{ is a critical point of } \text{Exp}_{q_0}\}$.

(ii) The *cut time* is $t_c(\lambda_0) = \min\{t > 0, \exists \lambda_1 \in \Lambda_{q_0}, \lambda_1 \neq \lambda_0 \text{ s.t. } \text{Exp}_{q_0}(t_c(\lambda_0)\lambda_0) = \text{Exp}_{q_0}(t_c(\lambda_0)\lambda_1)\}$.

It is well known that if a geodesic is not abnormal then it loses optimality either at the cut or at the conjugate locus (see for instance [4]).

2.3 Hausdorff measures

In this section we recall definitions of Hausdorff measure and spherical Hausdorff measure. We start with the definition of smooth volume.

Definition 18. Let M be a n -dimensional smooth manifold, which is connected and orientable. By a *smooth volume* on M we mean a measure μ on M associated to a smooth non-vanishing n -form $\omega_\mu \in \Lambda^n M$, i.e. for every subset $A \subset M$ we set

$$\mu(A) = \int_A \omega_\mu.$$

The Popp volume \mathcal{P} , which is a smooth volume in the sense of Definition 18, is the volume associated to a n -form $\omega_{\mathcal{P}}$ that can be intrinsically defined via the sub-Riemannian structure (see [2, 35]).

Let (M, d) be a metric space and denote with \mathcal{B} the set of balls in M .

Definition 19. Let A be a subset of M and $\alpha > 0$.

The α -dimensional Hausdorff measure of A is

$$\mathcal{H}^\alpha(A) := \lim_{\delta \rightarrow 0} \mathcal{H}_\delta^\alpha(A),$$

where

$$\mathcal{H}_\delta^\alpha(A) := \inf\left\{\sum_{i=1}^{\infty} \text{diam}(A_i)^\alpha, A \subset \bigcup_{i=1}^{\infty} A_i, \text{diam}(A_i) < \delta\right\}.$$

The α -dimensional spherical Hausdorff measure of A is

$$\mathcal{S}^\alpha(A) := \lim_{\delta \rightarrow 0} \mathcal{S}_\delta^\alpha(A),$$

where

$$\mathcal{S}_\delta^\alpha(A) := \inf\left\{\sum_{i=1}^{\infty} \text{diam}(B_i)^\alpha, A \subset \bigcup_{i=1}^{\infty} B_i, B_i \in \mathcal{B}, \text{diam}(B_i) < \delta\right\}.$$

These two measures are commensurable since it holds (see [21])

$$2^{-\alpha} \mathcal{S}^\alpha(A) \leq \mathcal{H}^\alpha(A) \leq \mathcal{S}^\alpha(A), \quad \forall A \subset M. \quad (10)$$

The *Hausdorff dimension* of A is defined as

$$\inf\{\alpha > 0, \mathcal{H}^\alpha(A) = 0\} = \sup\{\alpha > 0, \mathcal{H}^\alpha(A) = +\infty\}. \quad (11)$$

Formula (10) guarantees that Hausdorff dimension of A does not change if we replace \mathcal{H}^α with \mathcal{S}^α in formula (11).

It is a standard fact that the Hausdorff dimension of a Riemannian manifold, considered as a metric space, coincides with its topological dimension. On the other side, we have the following

Theorem 20. *Let (M, Δ, \mathbf{g}) be a regular sub-Riemannian manifold. Its Hausdorff dimension as a metric space is*

$$Q = \sum_{i=1}^m ik_i, \quad k_i := \dim \Delta^i - \dim \Delta^{i-1}.$$

Moreover \mathcal{S}^Q is commensurable to a smooth volume μ on M , i.e. for every compact $K \subset M$ there exists $\alpha_1, \alpha_2 > 0$ such that

$$\alpha_1 \mathcal{S}^Q \leq \mu \leq \alpha_2 \mathcal{S}^Q. \quad (12)$$

This theorem was proved by Mitchell in [34]. In its original version it was stated for the Lebesgue measure and the standard Hausdorff measure.

3 The nilpotent approximation

In this section we briefly recall the concept of nilpotent approximation. For details see [5, 12].

3.1 Privileged coordinates

Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a sub-Riemannian manifold and (X_1, \dots, X_k) an orthonormal frame. Fix a point $q \in M$ and consider the flag of the distribution $\Delta_q^1 \subset \Delta_q^2 \subset \dots \subset \Delta_q^m$. Recall that $k_i = \dim \Delta_q^i - \dim \Delta_q^{i-1}$ for $i = 1, \dots, m$, and that $k_1 + \dots + k_m = n$.

Let O_q be an open neighborhood of the point $q \in M$. We say that a system of coordinates $\psi : O_q \rightarrow \mathbb{R}^n$ is *linearly adapted* to the flag if, in these coordinates, we have $\psi(q) = 0$ and

$$\psi_*(\Delta_q^i) = \mathbb{R}^{k_1} \oplus \dots \oplus \mathbb{R}^{k_i}, \quad \forall i = 1, \dots, m.$$

Consider now the splitting $\mathbb{R}^n = \mathbb{R}^{k_1} \oplus \dots \oplus \mathbb{R}^{k_m}$ and denote its elements $x = (x_1, \dots, x_m)$ where $x_i = (x_i^1, \dots, x_i^{k_i}) \in \mathbb{R}^{k_i}$. The space of all differential operators in \mathbb{R}^n with smooth coefficients forms an associative algebra with composition of operators as multiplication. The differential operators with polynomial coefficients form a subalgebra of this algebra with generators $1, x_i^j, \frac{\partial}{\partial x_i^j}$, where $i = 1, \dots, m; j = 1, \dots, k_i$. We define weights of generators as

$$\nu(1) = 0, \quad \nu(x_i^j) = i, \quad \nu\left(\frac{\partial}{\partial x_i^j}\right) = -i,$$

and the weight of monomials

$$\nu\left(y_1 \cdots y_\alpha \frac{\partial^\beta}{\partial z_1 \cdots \partial z_\beta}\right) = \sum_{i=1}^\alpha \nu(y_i) - \sum_{j=1}^\beta \nu(z_j).$$

Notice that a polynomial differential operator homogeneous with respect to ν (i.e. whose monomials are all of same weight) is homogeneous with respect to dilations $\delta_t : \mathbb{R}^n \rightarrow \mathbb{R}^n$ defined by

$$\delta_t(x_1, \dots, x_m) = (tx_1, t^2x_2, \dots, t^m x_m), \quad t > 0. \quad (13)$$

In particular for a homogeneous vector field X of weight h it holds $\delta_{t*}X = t^{-h}X$. A smooth vector field $X \in \text{Vec}(\mathbb{R}^n)$, as a first order differential operator, can be written as

$$X = \sum_{i,j} a_i^j(x) \frac{\partial}{\partial x_i^j}$$

and considering its Taylor expansion at the origin we can write the formal expansion

$$X \approx \sum_{h=-m}^{\infty} X^{(h)}$$

where $X^{(h)}$ is the homogeneous part of degree h of X (notice that every monomial of a first order differential operator has weight not smaller than $-m$). Define the filtration of $\text{Vec}(\mathbb{R}^n)$

$$\mathcal{D}^{(h)} = \{X \in \text{Vec}(\mathbb{R}^n) : X^{(i)} = 0, \forall i < h\}, \quad \ell \in \mathbb{Z}.$$

Definition 21. A system of coordinates $\psi : O_q \rightarrow \mathbb{R}^n$ defined near the point q is said *privileged* for a sub-Riemannian structure \mathbf{S} if these coordinates are linearly adapted to the flag and such that $\psi_*X_i \in \mathcal{D}^{(-1)}$ for every $i = 1, \dots, k$.

Theorem 22. *Privileged coordinates always exists. Moreover there exist $c_1, c_2 > 0$ such that in these coordinates, for all $\varepsilon > 0$ small enough, we have*

$$c_1 \text{Box}(\varepsilon) \subset B(q, \varepsilon) \subset c_2 \text{Box}(\varepsilon), \quad (14)$$

where $\text{Box}(\varepsilon) = \{x \in \mathbb{R}^n, |x_i| \leq \varepsilon^i\}$.

Existence of privileged coordinates is proved in [5, 7, 12, 13]. In the regular case the construction of privileged coordinates was also done in the context of hypoelliptic operators (see [39]). The second statement is known as *Ball-Box theorem* and a proof can be found in [12]. Notice however that privileged coordinates are not unique.

Definition 23. Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular sub-Riemannian manifold and (X_1, \dots, X_k) a local orthonormal frame near a point q . Fixed a system of privileged coordinates, we define the *nilpotent approximation of \mathbf{S} near q* , denoted by $\widehat{\mathbf{S}}_q$, the sub-Riemannian structure on \mathbb{R}^n having $(\widehat{X}_1, \dots, \widehat{X}_k)$ as an orthonormal frame, where $\widehat{X}_i := (\psi_*X_i)^{(-1)}$.

Remark 24. It is well known that under the regularity hypothesis, $\widehat{\mathbf{S}}_q$ is naturally endowed with a Lie group structure whose Lie algebra is generated by left-invariant vector fields $\widehat{X}_1, \dots, \widehat{X}_k$. Moreover the sub-Riemannian distance \widehat{d} in $\widehat{\mathbf{S}}_q$ is homogeneous with respect to dilations δ_t , i.e. $\widehat{d}(\delta_t(x), \delta_t(y)) = t\widehat{d}(x, y)$. In particular, if $\widehat{B}_q(r)$ denotes the ball of radius r in $\widehat{\mathbf{S}}_q$, this implies $\delta_t(\widehat{B}_q(1)) = \widehat{B}_q(t)$.

Theorem 25. *The nilpotent approximation $\widehat{\mathbf{S}}_q$ of a sub-Riemannian structure \mathbf{S} near a point q is the metric tangent space to M at point q in the sense of Gromov, that means*

$$\delta_{1/\varepsilon}B(q, \varepsilon) \longrightarrow \widehat{B}_q, \quad (15)$$

where \widehat{B}_q denotes the sub-Riemannian unit ball of the nilpotent approximation $\widehat{\mathbf{S}}_q$.

Remark 26. Convergence of sets in (15) is intended in the Gromov-Hausdorff topology [12, 27]. In the regular case this theorem was proved by Mitchell in [33]. A proof in the general case can be found in [12].

Definition 27. Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular sub-Riemannian manifold and $\widehat{\mathbf{S}}_q$ its nilpotent approximation near q . If μ a smooth volume on M , associated to the smooth non-vanishing n -form ω_μ , we define the *induced volume* $\widehat{\mu}_q$ at the point q as the left-invariant volume on $\widehat{\mathbf{S}}_q$ canonically associated with $\omega_\mu(q) \in \wedge^n(T_q^*M)$ (cf. isomorphism (1)).

From Theorem 25 and the relation³ $\mu(\delta_\varepsilon A) = \varepsilon^Q \widehat{\mu}_q(A) + o(\varepsilon^Q)$ when $\varepsilon \rightarrow 0$, one gets

Corollary 28. *Let μ be a smooth volume on M and $\widehat{\mu}_q$ the induced volume on the nilpotent approximation at point q . Then, for $\varepsilon \rightarrow 0$, we have*

$$\mu(B(q, \varepsilon)) = \varepsilon^Q \widehat{\mu}_q(\widehat{B}_q) + o(\varepsilon^Q).$$

3.2 Normal forms for nilpotent approximation in dimension ≤ 5

In this section we provide normal forms for the nilpotent approximation of regular sub-Riemannian structures in dimension less or equal than 5. One can easily shows that in this case the only possibilities for growth vectors are:

- $\dim(M) = 3$: $\mathcal{G}(\mathbf{S}) = (2, 3)$,
- $\dim(M) = 4$: $\mathcal{G}(\mathbf{S}) = (2, 3, 4)$ or $\mathcal{G}(\mathbf{S}) = (3, 4)$,
- $\dim(M) = 5$: $\mathcal{G}(\mathbf{S}) \in \{(2, 3, 4, 5), (2, 3, 5), (3, 5), (3, 4, 5), (4, 5)\}$.

We have the following.

Theorem 29. *Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular sub-Riemannian manifold and $\widehat{\mathbf{S}}_q$ its nilpotent approximation near q . Up to a change of coordinates and a rotation of the orthonormal frame we have the following expression for the orthonormal frame of $\widehat{\mathbf{S}}_q$:*

Case $n = 3$ • $\mathcal{G}(\mathbf{S}) = (2, 3)$. (*Heisenberg*)

$$\begin{aligned}\widehat{X}_1 &= \partial_1, \\ \widehat{X}_2 &= \partial_2 + x_1 \partial_3.\end{aligned}$$

Case $n = 4$ • $\mathcal{G}(\mathbf{S}) = (2, 3, 4)$. (*Engel*)

$$\begin{aligned}\widehat{X}_1 &= \partial_1, \\ \widehat{X}_2 &= \partial_2 + x_1 \partial_3 + x_1 x_2 \partial_4.\end{aligned}$$

• $\mathcal{G}(\mathbf{S}) = (3, 4)$. (*Quasi-Heisenberg*)

$$\begin{aligned}\widehat{X}_1 &= \partial_1, \\ \widehat{X}_2 &= \partial_2 + x_1 \partial_4, \\ \widehat{X}_3 &= \partial_3.\end{aligned}$$

³Notice that this formula is meaningful in privileged coordinates near q .

Case $n = 5$ • $\mathcal{G}(\mathbf{S}) = (2, 3, 5)$. (*Cartan*)

$$\begin{aligned}\widehat{X}_1 &= \partial_1, \\ \widehat{X}_2 &= \partial_2 + x_1\partial_3 + \frac{1}{2}x_1^2\partial_4 + x_1x_2\partial_5.\end{aligned}$$

• $\mathcal{G}(\mathbf{S}) = (2, 3, 4, 5)$. (*Goursat rank 2*)

$$\begin{aligned}\widehat{X}_1 &= \partial_1, \\ \widehat{X}_2 &= \partial_2 + x_1\partial_3 + \frac{1}{2}x_1^2\partial_4 + \frac{1}{6}x_1^3\partial_5.\end{aligned}$$

• $\mathcal{G}(\mathbf{S}) = (3, 5)$. (*Corank 2*)

$$\begin{aligned}\widehat{X}_1 &= \partial_1 - \frac{1}{2}x_2\partial_4, \\ \widehat{X}_2 &= \partial_2 + \frac{1}{2}x_1\partial_4 - \frac{1}{2}x_3\partial_5, \\ \widehat{X}_3 &= \partial_3 + \frac{1}{2}x_2\partial_5.\end{aligned}$$

• $\mathcal{G}(\mathbf{S}) = (3, 4, 5)$. (*Goursat rank 3*)

$$\begin{aligned}\widehat{X}_1 &= \partial_1 - \frac{1}{2}x_2\partial_4 - \frac{1}{3}x_1x_2\partial_5, \\ \widehat{X}_2 &= \partial_2 + \frac{1}{2}x_1\partial_4 + \frac{1}{3}x_1^2\partial_5, \\ \widehat{X}_3 &= \partial_3.\end{aligned}$$

• $\mathcal{G}(\mathbf{S}) = (4, 5)$. (*Bi-Heisenberg*)

$$\begin{aligned}\widehat{X}_1 &= \partial_1 - \frac{1}{2}x_2\partial_5, \\ \widehat{X}_2 &= \partial_2 + \frac{1}{2}x_1\partial_5, \\ \widehat{X}_3 &= \partial_3 - \frac{\alpha}{2}x_4\partial_5, \quad \alpha \in \mathbb{R}, \\ \widehat{X}_4 &= \partial_4 + \frac{\alpha}{2}x_3\partial_5.\end{aligned} \tag{16}$$

Proof. It is sufficient to find, for every such a structure, a basis of the Lie algebra such that the structural constants⁴ are uniquely determined by the sub-Riemannian structure. We give a sketch of the proof for the $(2, 3, 4, 5)$ and $(3, 4, 5)$ and $(4, 5)$ cases. The other cases can be treated in a similar way.

(i). Let $\widehat{\mathbf{S}} = (G, \Delta, \mathfrak{g})$ be a nilpotent $(3, 4, 5)$ sub-Riemannian structure. Since we deal with a left-invariant sub-Riemannian structure, we can identify the distribution Δ with its value at the

⁴Let X_1, \dots, X_k be a basis of a Lie algebra \mathfrak{g} . The coefficients c_{ij}^ℓ that satisfy $[X_i, X_j] = \sum_\ell c_{ij}^\ell X_\ell$ are called structural constant of \mathfrak{g} .

identity of the group Δ_{id} . Let $\{e_1, e_2, e_3\}$ be a basis for Δ_{id} , as a vector subspace of the Lie algebra. By our assumption on the growth vector we know that

$$\dim \text{span}\{[e_1, e_2], [e_1, e_3], [e_2, e_3]\}/\Delta_{id} = 1. \quad (17)$$

In other words, we can consider the skew-symmetric mapping

$$\Phi(\cdot, \cdot) := [\cdot, \cdot]/\Delta_{id} : \Delta_{id} \times \Delta_{id} \rightarrow T_{id}G/\Delta_{id}, \quad (18)$$

and condition (17) implies that there exists a one dimensional subspace in the kernel of this map. Let \widehat{X}_3 be a normalized vector in the kernel and consider its orthogonal subspace $D \subset \Delta_{id}$ with respect to the Euclidean product on Δ_{id} . Fix an arbitrary orthonormal basis $\{X_1, X_2\}$ of D and set $\widehat{X}_4 := [X_1, X_2]$. It is easy to see that \widehat{X}_4 does not change if we rotate the base $\{X_1, X_2\}$ and there exists a choice of this frame, denoted $\{\widehat{X}_1, \widehat{X}_2\}$, such that $[\widehat{X}_2, \widehat{X}_4] = 0$. Then set $\widehat{X}_5 := [\widehat{X}_1, \widehat{X}_4]$. Therefore we found a canonical basis for the Lie algebra that satisfies the following commutator relations:

$$[\widehat{X}_1, \widehat{X}_2] = \widehat{X}_4, \quad [\widehat{X}_1, \widehat{X}_4] = \widehat{X}_5,$$

and all other commutators vanish. A standard application of the Campbell-Hausdorff formula gives the coordinate expression above.

(ii). Let us assume now that $\widehat{\mathbf{S}}$ is a nilpotent (2, 3, 4, 5) sub-Riemannian structure. As before we identify the distribution Δ with its value at the identity and consider any orthonormal basis $\{e_1, e_2\}$ for the 2-dimensional subspace Δ_{id} . By our assumption on $\mathcal{G}(\mathbf{S})$

$$\begin{aligned} \dim \text{span}\{e_1, e_2, [e_1, e_2]\} &= 3 \\ \dim \text{span}\{e_1, e_2, [e_1, e_2], [e_1, [e_1, e_2]], [e_2, [e_1, e_2]]\} &= 4. \end{aligned} \quad (19)$$

As in (i), it is easy to see that there exists a choice of the orthonormal basis on Δ_{id} , which we denote $\{\widehat{X}_1, \widehat{X}_2\}$, such that $[\widehat{X}_2, [\widehat{X}_1, \widehat{X}_2]] = 0$. From this property and the Jacobi identity it follows $[\widehat{X}_2, [\widehat{X}_1, [\widehat{X}_1, \widehat{X}_2]]] = 0$. Then we set $\widehat{X}_3 := [\widehat{X}_1, \widehat{X}_2]$, $\widehat{X}_4 = [\widehat{X}_1, [\widehat{X}_1, \widehat{X}_2]]$ and $\widehat{X}_5 := [\widehat{X}_1, [\widehat{X}_1, [\widehat{X}_1, \widehat{X}_2]]]$. It is easily seen that (19) implies that these vectors are linearly independent and give a canonical basis for the Lie algebra, with the only nontrivial commutator relations:

$$[\widehat{X}_1, \widehat{X}_2] = \widehat{X}_3, \quad [\widehat{X}_1, \widehat{X}_3] = \widehat{X}_4, \quad [\widehat{X}_1, \widehat{X}_4] = \widehat{X}_5.$$

(iii). In the case (4, 5) since $\dim T_{id}G/\Delta_{id} = 1$, the map (18) is represented by a single 4×4 skew-symmetric matrix L . By skew-symmetry its eigenvalues are purely imaginary $\pm ib_1, \pm ib_2$, one of which is different from zero. Assuming $b_1 \neq 0$ we have that $\alpha = b_2/b_1$. Notice that the structure is contact if and only if $\alpha \neq 0$ (see also Section 5.1 for more details on the normal form). \square

Remark 30. Notice that in all other cases but the last the nilpotent approximation does not depend on any parameter, while in the (4, 5) case (and starting from dimension 5 in general) the sub-Riemannian structure induced on the tangent space depends on the point.

As a consequence, the Popp's measure \mathcal{P} also depends on the point. In particular, being \mathcal{P} smooth by definition, in the (4, 5) case it must depend smoothly on b_1, b_2 , the eigenvalues of the skew-symmetric matrix that represent the Lie bracket map. Indeed it can be computed that

$$\mathcal{P} = \frac{1}{\sqrt{b_1^2 + b_2^2}} dx_1 \wedge \dots \wedge dx_5$$

The uniqueness of the normal form in Theorem 29, when $\mathcal{G}(\mathbf{S}) \neq (4, 5)$, implies the following corollary

Corollary 31. *Let $\mathbf{S} = (M, \Delta, \mathbf{g})$ be a regular sub-Riemannian manifold such that $\dim(M) \leq 5$ and $\mathcal{G}(\mathbf{S}) \neq (4, 5)$. Then if $q_1, q_2 \in M$ we have that $\widehat{\mathbf{S}}_{q_1}$ is isometric to $\widehat{\mathbf{S}}_{q_2}$ as sub-Riemannian manifolds.*

4 Proof of Theorem 1: the density is the volume of nilpotent balls

In this section we prove Theorem 1, i.e.

$$f_{\mu\mathcal{S}}(q) = \frac{1}{2^Q} \widehat{\mu}_q(\widehat{B}_q). \quad (20)$$

It is well known that, being μ absolutely continuous with respect to \mathcal{S}^Q (see Theorem 20), the Radon-Nikodym derivative of μ with respect to \mathcal{S}^Q , namely $f_{\mu\mathcal{S}}$, can be computed almost everywhere as

$$\lim_{r \rightarrow 0} \frac{\mu(B(q, r))}{\mathcal{S}^Q(B(q, r))}.$$

By Corollary 28 we get

$$\frac{\mu(B(q, r))}{\mathcal{S}^Q(B(q, r))} = \frac{r^Q \widehat{\mu}_q(\widehat{B}_q) + o(r^Q)}{\mathcal{S}^Q(B(q, r))} = \frac{\widehat{\mu}_q(\widehat{B}_q)}{2^Q} \frac{2^Q r^Q}{\mathcal{S}^Q(B(q, r))} + \frac{o(r^Q)}{\mathcal{S}^Q(B(q, r))}.$$

Then we are left to prove the following

Lemma 32. *Let A be an open subset of M . For \mathcal{S}^Q -a.e. $q \in A$ we have*

$$\lim_{r \rightarrow 0} \frac{\mathcal{S}^Q(A \cap B(q, r))}{(2r)^Q} = 1. \quad (21)$$

Proof. In the following proof we make use of Vitali covering lemma⁵ and we always assume that balls of our covering are small enough to satisfy property (ii) of Proposition 15.

We prove that the set where (21) exists and is different from 1 has \mathcal{S}^Q -null measure.

(i). First we show

$$\mathcal{S}^Q(E_\delta) = 0, \quad \forall 0 < \delta \leq 1,$$

where

$$E_\delta := \{q \in A : \mathcal{S}^Q(A \cap B(q, r)) < (1 - \delta)(2r)^Q, \forall 0 < r < \delta\}$$

Let $\{B_i\}$ a ball covering of E_δ with $\text{diam}(B_i) < \delta$ and such that

$$\sum_i \text{diam}(B_i)^Q \leq \mathcal{S}_\delta^Q(E_\delta) + \varepsilon \leq \mathcal{S}^Q(E_\delta) + \varepsilon.$$

⁵**Theorem.**(Vitali covering lemma, [21, 9]) Let E be a metric space, $B \subset E$ and $\alpha > 0$ such that $\mathcal{H}^\alpha(B) < \infty$, and let \mathcal{F} be a fine covering of B . Then there exist a countable disjoint subfamily $\{V_i\} \subset \mathcal{F}$ such that

$$\mathcal{H}^\alpha(B \setminus \bigcup V_i) = 0.$$

We recall that \mathcal{F} is a *fine covering* of B if for every $x \in B$ and $\varepsilon > 0$ there exists $V \in \mathcal{F}$ such that $x \in V$ and $\text{diam}(V) < \varepsilon$.

Then we have

$$\begin{aligned}
\mathcal{S}^Q(E_\delta) &\leq \mathcal{S}^Q(A \cap \bigcup B_i) \\
&\leq \sum \mathcal{S}^Q(A \cap B_i) \\
&\leq (1 - \delta) \sum \text{diam}(B_i)^Q \\
&\leq (1 - \delta)(\mathcal{S}^Q(E_\delta) + \varepsilon).
\end{aligned}$$

then $\varepsilon \rightarrow 0$ and $1 - \delta < 1$ implies $\mathcal{S}^Q(E_\delta) = 0$.

(ii). Next we prove that

$$\mathcal{S}^Q(E_t) = 0, \quad \forall t > 1,$$

where

$$E_t := \{q \in A : \mathcal{S}^Q(A \cap B(q, r)) > t(2r)^Q, \forall r \text{ small enough}\}.$$

Now let U be an open set such that $E_t \subset U$ and $\mathcal{S}^Q(A \cap U) < \mathcal{S}^Q(E_t) + \varepsilon$. We define

$$\mathcal{F} := \{B(q, r) : q \in E_t, B(q, r) \subset U, \text{diam } B(q, r) \leq \delta\}.$$

Now we can apply Vitali covering lemma to \mathcal{F} and get a family $\{B_i\}$ of disjoint balls such that $\mathcal{S}^Q(E_t \setminus \bigcup_i B_i) = 0$. Then we get

$$\begin{aligned}
\mathcal{S}^Q(E_t) + \varepsilon &> \mathcal{S}^Q(A \cap U) \\
&\geq \mathcal{S}^Q(A \cap \bigcup B_i) \\
&\geq t \sum \text{diam}(B_i)^Q \\
&\geq t \mathcal{S}_\delta^Q(E_t \cap \bigcup B_i) \\
&\geq t \mathcal{S}_\delta^Q(E_t).
\end{aligned}$$

Letting $\varepsilon, \delta \rightarrow 0$ we have an absurd because $t > 1$. □

Since A is open, from this lemma follows formula (20).

Remark 33. Notice that, for a n -dimensional Riemannian manifold, the tangent spaces at different points are isometric. As a consequence the Riemannian volume of the unit ball in the tangent space is constant and equal to $C_n = \pi^{\frac{n}{2}} / \Gamma(\frac{n}{2} + 1)$. The above formula, where $\mu = \text{Vol}$ is the Riemannian volume, recover the well-known relation between Vol and the (spherical) Hausdorff measure

$$d\text{Vol} = \frac{C_n}{2^n} d\mathcal{S}^n = \frac{C_n}{2^n} d\mathcal{H}^n$$

4.1 Proof of Corollary 2: continuity of the density

In this section we prove Theorem 2. More precisely we study the continuity of the map

$$f_{\mu\mathcal{S}} : q \mapsto \hat{\mu}_q(\hat{B}_q). \tag{22}$$

To this purpose, it is sufficient to study the regularity under the hypothesis that $\hat{\mu}_q$ does not depend on the point. Indeed it is easily seen that the smooth measure μ , which is defined on the manifold,

induces on the nilpotent approximations a smooth family of measures $\{\widehat{\mu}_q\}_{q \in M}$. If μ is the Popp measure is a consequence of equality (2). In other words we can identify all tangent spaces in coordinates with \mathbb{R}^n and fix a measure $\widehat{\mu}$ on it.

We are then reduced to study the regularity of the volume of the unit ball of a smooth family of nilpotent structures in \mathbb{R}^n , with respect to a fixed smooth measure. Notice that this family depends on an n -dimensional parameter.

To sum up, we have left to study the regularity of the map

$$q \mapsto \mathcal{L}(\widehat{B}_q), \quad q \in M, \quad (23)$$

where \widehat{B}_q is the unit ball of a family of nilpotent structures $\widehat{\mathbf{S}}_q$ in \mathbb{R}^n and \mathcal{L} is the standard Lebesgue measure.

Let us denote \widehat{d}_q the sub-Riemannian distance in $\widehat{\mathbf{S}}_q$ and $\rho_q := \widehat{d}_q(0, \cdot)$. Following this notation $\widehat{B}_q = \{x \in \mathbb{R}^n \mid \rho_q(x) \leq 1\}$ and the coordinate expression (13) implies that

$$\mathcal{L}(\delta_\alpha(\widehat{B}_q)) = \alpha^Q \mathcal{L}(\widehat{B}_q), \quad \forall \alpha > 0. \quad (24)$$

Notice that, since our sub-Riemannian structure is regular, we can choose privileged coordinates $\psi_q : O_q \rightarrow \mathbb{R}^n$ smoothly with respect to q . Let now $q' \neq q$, there exists $\alpha = \alpha(q, q')$ such that (see Remark 24)

$$\delta_{\frac{1}{\alpha}} \widehat{B}_{q'} \subset \widehat{B}_q \subset \delta_\alpha \widehat{B}_{q'}. \quad (25)$$

Using (24), (25) and monotonicity of the volume we get

$$\left(\frac{1}{\alpha^Q} - 1\right) \mathcal{L}(\widehat{B}_{q'}) \leq \mathcal{L}(\widehat{B}_q) - \mathcal{L}(\widehat{B}_{q'}) \leq (\alpha^Q - 1) \mathcal{L}(\widehat{B}_{q'}).$$

It is sufficient then to show that $\alpha(q, q') \rightarrow 1$ when $q' \rightarrow q$. This property follows from the next

Lemma 34. *The family of functions $\rho_q|_K$ is equicontinuous for every compact $K \subset \mathbb{R}^n$. Moreover $\rho_{q'} \rightarrow \rho_q$ uniformly on compacts in \mathbb{R}^n , as $q' \rightarrow q$.*

In the case in which $\{\rho_t\}_{t>0}$ is the approximating family of the nilpotent distance $\widehat{\rho}$, this result is proved in [3]. See also [1] for a more detailed proof, using chronological calculus. With the same arguments one can extend this result to any smooth family of regular sub-Riemannian structures. The key point is that we can construct a basis for the tangent space to the structure with bracket polynomials of the orthonormal frame where the structure of the brackets does not depend on the parameter.

5 Proof of Theorem 5: differentiability of the density in the corank 1 case

In this section we prove Theorem 5, We start by studying the contact case. Then we complete our analysis by reducing the quasi-contact case and the general case to the contact one.

5.1 Normal form of the nilpotent contact case

Consider a 2-step nilpotent sub-Riemannian manifold in \mathbb{R}^n of rank k .

Select a basis $\{X_1, \dots, X_k, Z_1, \dots, Z_{n-k}\}$ such that

$$\begin{cases} \Delta = \text{span}\{X_1, \dots, X_k\}, \\ [X_i, X_j] = \sum_{h=1}^{n-k} b_{ij}^h Z_h, & i, j = 1, \dots, k, \quad \text{where } b_{ij}^h = -b_{ji}^h, \\ [X_i, Z_j] = [Z_j, Z_h] = 0, & i = 1, \dots, k, \quad j, h = 1, \dots, n-k. \end{cases} \quad (26)$$

Hence the Lie bracket can be considered as a map

$$[\cdot, \cdot] : \Delta \times \Delta \longrightarrow TM/\Delta \quad (27)$$

and is represented by the $n - k$ skew-symmetric matrices $L^h = (b_{ij}^h), h = 1, \dots, n - k$.

In the contact case we have $(k, n) = (2\ell, 2\ell + 1)$ and our structure is represented by one non degenerate skew-symmetric matrix L . Take coordinates in such a way that L is normalized in the following block-diagonal form

$$L = \begin{pmatrix} B_1 & & \\ & \ddots & \\ & & B_\ell \end{pmatrix}, \quad \text{where } B_i := \begin{pmatrix} 0 & -b_i \\ b_i & 0 \end{pmatrix}, \quad b_i > 0.$$

with eigenvalues $\pm ib_1, \dots, \pm ib_\ell$. Hence we can find vector fields $\{X_1, \dots, X_\ell, Y_1, \dots, Y_\ell, Z\}$ such that relations (26) reads

$$\begin{cases} \Delta = \text{span}\{X_1, \dots, X_\ell, Y_1, \dots, Y_\ell\}, \\ [X_i, Y_i] = -b_i Z, & i = 1, \dots, \ell \\ [X_i, Y_j] = 0, & i \neq j \\ [X_i, Z] = [Y_i, Z] = 0, & i = 1, \dots, \ell \end{cases} \quad (28)$$

In the following we call b_1, \dots, b_ℓ *frequencies* of the contact structure.

We can recover the product on the group by the Campbell-Hausdorff formula. If we denote points $q = (x, y, z)$, where

$$x = (x_1, \dots, x_\ell) \in \mathbb{R}^\ell, \quad y = (y_1, \dots, y_\ell) \in \mathbb{R}^\ell, \quad z \in \mathbb{R},$$

we can write the group law in coordinates

$$q \cdot q' = \left(x + x', y + y', z + z' - \frac{1}{2} \sum_{i=1}^{\ell} b_i (x_i x'_i - y_i y'_i) \right). \quad (29)$$

Finally, from (29), we get the coordinate expression of the left-invariant vector fields of the Lie algebra, namely

$$\begin{aligned} X_i &= \partial_{x_i} + \frac{1}{2} b_i y_i \partial_z, & i = 1, \dots, \ell, \\ Y_i &= \partial_{y_i} - \frac{1}{2} b_i x_i \partial_z, & i = 1, \dots, \ell, \\ Z &= \partial_z. \end{aligned} \quad (30)$$

In this expression one of frequencies b_i can be normalized to 1.

5.2 Exponential map in the nilpotent contact case

Now we apply the PMP to find the exponential map Exp_{q_0} where q_0 is the origin. Define the hamiltonians (linear on fibers)

$$h_{u_i}(\lambda) = \langle \lambda, X_i(q) \rangle, \quad h_{v_i}(\lambda) = \langle \lambda, Y_i(q) \rangle, \quad h_w(\lambda) = \langle \lambda, Z(q) \rangle.$$

Recall from Section 2.2 that $q(t)$ is a normal extremal if and only if there exists $\lambda(t)$ such that

$$\begin{cases} \dot{u}_i = -b_i w v_i \\ \dot{v}_i = b_i w u_i \\ \dot{w} = 0 \end{cases} \quad \begin{cases} \dot{x}_i = u_i \\ \dot{y}_i = v_i \\ \dot{z} = \frac{1}{2} \sum_i b_i (u_i y_i - v_i x_i) \end{cases} \quad (31)$$

where

$$u_i(t) := h_{u_i}(\lambda(t)), \quad v_i(t) := h_{v_i}(\lambda(t)), \quad w(t) := h_w(\lambda(t)).$$

Remark 35. Notice that from (31) it follows that the sub-Riemannian length of a geodesic coincide with the Euclidean length of its projection on the horizontal subspace $(x_1, \dots, x_n, y_1, \dots, y_n)$.

$$l(\gamma) = \int_0^T \left(\sum_i u_i^2(t) + v_i^2(t) \right)^{\frac{1}{2}} dt$$

Now we solve (31) with initial conditions (see also Remark 13)

$$(x^0, y^0, z^0) = (0, 0, 0), \quad (32)$$

$$(u^0, v^0, w^0) = (u_1^0, \dots, u_\ell^0, v_1^0, \dots, v_\ell^0, w^0) \in S^{2\ell-1} \times \mathbb{R}. \quad (33)$$

Notice that $w \equiv w^0$ is constant on geodesics. We consider separately the two cases:

(i) If $w \neq 0$, we have (denoting $a_i := b_i w$)

$$\begin{aligned} u_i(t) &= u_i^0 \cos a_i t - v_i^0 \sin a_i t, \\ v_i(t) &= u_i^0 \sin a_i t + v_i^0 \cos a_i t, \\ w(t) &= w. \end{aligned} \quad (34)$$

From (31) one easily get

$$\begin{aligned} x_i(t) &= \frac{1}{a_i} (u_i^0 \sin a_i t + v_i^0 \cos a_i t - v_i^0), \\ y_i(t) &= \frac{1}{a_i} (-u_i^0 \cos a_i t + v_i^0 \sin a_i t + u_i^0), \\ z(t) &= \frac{1}{2w^2} (wt - \sum_i \frac{1}{b_i} ((u_i^0)^2 + (v_i^0)^2)) \sin a_i t. \end{aligned} \quad (35)$$

(ii) If $w = 0$, we find equations of straight lines on the horizontal plane in direction of the vector (u^0, v^0) :

$$x_i(t) = u_i^0 t \quad y_i(t) = v_i^0 t \quad z(t) = 0.$$

Remark 36. To recover symmetry properties of the exponential map it is useful to rewrite (35) in polar coordinates, using the following change of variables

$$u_i = r_i \cos \theta_i, \quad v_i = r_i \sin \theta_i, \quad i = 1, \dots, \ell. \quad (36)$$

In these new coordinates (35) becomes

$$\begin{aligned} x_i(t) &= \frac{r_i}{a_i} (\cos(a_i t + \theta_i) - \cos \theta_i), \\ y_i(t) &= \frac{r_i}{a_i} (\sin(a_i t + \theta_i) - \sin \theta_i), \\ z(t) &= \frac{1}{2w^2} (wt - \sum_i \frac{r_i^2}{b_i} \sin a_i t). \end{aligned} \quad (37)$$

and the condition $(u^0, v^0) \in S^{2\ell-1}$ implies that $r = (r_1, \dots, r_\ell) \in S^\ell$.

From equations (37) we easily see that the projection of a geodesic on every 2-plane (x_i, y_i) is a circle, with period T_i , radius ρ_i and center C_i where

$$T_i = \frac{2\pi}{b_i w}, \quad \rho_i = \frac{r_i}{b_i w} \quad C_i = -\frac{r_i}{b_i w} (\cos \theta_i, \sin \theta_i), \quad \forall i = 1, \dots, \ell \quad (38)$$

Moreover (31) shows that z component of the geodesic is the weighted sum (with coefficients b_i) of areas of these circular segments.

Lemma 37. *Let $\gamma(t)$ be a geodesic starting from the origin and corresponding to the parameters (r_i, θ_i, w) . The cut time t_c along γ is equal to first conjugate time and satisfies*

$$t_c = \frac{2\pi}{w \max_i b_i}, \quad (39)$$

with the understanding that $t_c = +\infty$, if $w = 0$.

Proof. We divide the proof into two steps. Recall that a geodesic lose optimality either at a cut time or at a conjugate time, see Definition 17. First we prove that (39) is a conjugate time and then that for every $t < t_c$ our geodesic is optimal.

The case $w = 0$ is trivial. Indeed the geodesic is a straight line and, by Remark 35, we have neither cut nor conjugate points. Then it is not restrictive to assume $w \neq 0$. Moreover, up to relabeling indices, we can also assume that $b_1 \geq b_2 \geq \dots \geq b_\ell \geq 0$.

(i) Since, by assumption, $b_1 = \max_i b_i$, from (37) it is easily seen that projection on the (x_1, y_1) -plane satisfies

$$x_1(t_c) = y_1(t_c) = 0.$$

Next consider the one parametric family of geodesic with initial condition

$$(r_1, r_2, \dots, r_\ell, \theta_1 + \phi, \theta_2, \dots, \theta_\ell, w), \quad \phi \in [0, 2\pi].$$

It is easily seen from equation (37) that all these curves have the same endpoints. Indeed neither (x_i, y_i) , for $i > 1$, nor z depends on this variable. Then it follows that t_c is a critical time for exponential map, hence a conjugate time.

(ii) Since $w \neq 0$, our geodesic is non horizontal (i.e. $z(t) \not\equiv 0$). We know that, for every i , the

projection of every non horizontal geodesic on on the plane (x_i, y_i) is a circle. Moreover for the i -th projected curve, the distance from the origin is easily computed

$$\eta_i(t) = \sqrt{x_i(t)^2 + y_i(t)^2} = r_i t \operatorname{sin}_c\left(\frac{b_i w t}{2}\right), \quad \text{where} \quad \operatorname{sin}_c(x) = \frac{\sin x}{x}$$

Let now $\bar{t} < t_c$, we want to show that there are no others geodesics $\tilde{\gamma}(t)$, starting from the origin, that reach optimally the point $\gamma(\bar{t})$ at the same time \bar{t} . Assume that $\tilde{\gamma}(t)$ is associated to the parameters $(\tilde{r}_i, \tilde{\theta}_i, \tilde{w})$, where $(\tilde{r}_1, \dots, \tilde{r}_\ell) \in S^\ell$, and let us argue by contradiction. If $\gamma(\bar{t}) = \tilde{\gamma}(\bar{t})$ it follows that $\eta_i(\bar{t}) = \tilde{\eta}_i(\bar{t})$ for every i , that means

$$r_i \operatorname{sin}_c\left(\frac{b_i w \bar{t}}{2}\right) = \tilde{r}_i \bar{t} \operatorname{sin}_c\left(\frac{b_i \tilde{w} \bar{t}}{2}\right), \quad i = 1, \dots, \ell \quad (40)$$

Notice that, once \tilde{w} is fixed, \tilde{r}_i are uniquely determined by (40) (recall that \bar{t} is fixed). Moreover, $\tilde{\theta}_i$ also are uniquely determined by relations (38). Finally, from the assumption that $\tilde{\gamma}$ also reach optimally the point $\tilde{\gamma}(\bar{t})$, it follows that

$$\bar{t} < t_c = \frac{2\pi}{b_1 \tilde{w}} \implies \frac{b_i \tilde{w} \bar{t}}{2} < \pi \quad \forall i = 1, \dots, \ell. \quad (41)$$

Since $\operatorname{sin}_c(x)$ is a monotone increasing function on $[-\pi, \pi]$ it follows from (40) that, if $\tilde{w} \neq w$ we have

$$\begin{aligned} \tilde{w} > w &\implies \tilde{r}_i < r_i \quad \forall i = 1, \dots, \ell \implies \sum_i \tilde{r}_i^2 < \sum_i r_i^2 = 1 \\ \tilde{w} < w &\implies \tilde{r}_i > r_i \quad \forall i = 1, \dots, \ell \implies \sum_i \tilde{r}_i^2 > \sum_i r_i^2 = 1 \end{aligned}$$

which contradicts the fact that $(\tilde{r}_1, \dots, \tilde{r}_\ell) \in S^\ell$. □

Consider now the exponential map from the origin

$$(r, \theta, w) \mapsto \operatorname{Exp}(r, \theta, w) = \begin{cases} x_i = \frac{r_i}{b_i w} (\cos(b_i w + \theta_i) - \cos \theta_i) \\ y_i = \frac{r_i}{b_i w} (\sin(b_i w + \theta_i) - \sin \theta_i) \\ z = \frac{1}{2w^2} (w|r|^2 - \sum_i \frac{r_i^2}{b_i} \sin b_i w) = \sum_i \frac{r_i^2}{2b_i w^2} (b_i w - \sin b_i w) \end{cases} \quad (42)$$

By Lemma 37, the set D where geodesics are optimal and their length is less or equal to 1, is characterized as follows

$$D := \{(r, \theta, w), |r| \leq 1, |w| \leq 2\pi / \max b_i\},$$

Thus the restriction of the exponential map to the interior of D gives a regular parametrization of the nilpotent unit ball \widehat{B} , and we can compute its volume with the change of variables formula

$$\mathcal{L}(\widehat{B}) = \int_{\widehat{B}} d\mathcal{L} = \int_D |\det J| R dr d\theta dw. \quad (43)$$

where J denotes the Jacobian matrix of the exponential map (42). Notice that we have to integrate with respect to the measure $dudvdw = R dr d\theta dw$, where $R = \prod_i r_i$ because of the change of variables (36).

Lemma 38. *The Jacobian of Exponential map (42) is given by the formula*

$$\det J(r, \theta, w) = \frac{4^\ell R}{B^2 w^{2\ell+2}} \sum_{i=1}^{\ell} \left(\prod_{j \neq i} \sin^2\left(\frac{b_j w}{2}\right) \right) \sin\left(\frac{b_i w}{2}\right) \left(\frac{b_i w}{2} \cos\left(\frac{b_i w}{2}\right) - \sin\left(\frac{b_i w}{2}\right) \right) r_i^2, \quad (44)$$

where we denote with $B = \prod_i b_i$.

Proof. We reorder variables in the following way

$$(r_1, \theta_1, \dots, r_\ell, \theta_\ell, w), \quad (x_1, y_1, \dots, x_\ell, y_\ell, z),$$

in such a way that the Jacobian matrix J of the exponential map (42) is

$$J = \begin{pmatrix} Q_1 & & & W_1 \\ & Q_2 & & W_2 \\ & & \ddots & \vdots \\ & & & Q_\ell & W_\ell \\ Z_1 & Z_2 & \dots & Z_\ell & \partial_w z \end{pmatrix} \quad (45)$$

where we denote

$$Q_i = (Q_i^r \quad Q_i^\theta) := \begin{pmatrix} \partial_{r_i} x_i & \partial_{\theta_i} x_i \\ \partial_{r_i} y_i & \partial_{\theta_i} y_i \end{pmatrix}, \quad i = 1, \dots, \ell, \quad (46)$$

$$W_i := \begin{pmatrix} \partial_w x_i \\ \partial_w y_i \end{pmatrix}, \quad Z_i := (\partial_{r_i} z \quad \partial_{\theta_i} z).$$

Notice that x_i, y_i depend only on r_i, θ_i .

To compute the determinant of J , we write

$$z = z(r, w) = \sum_i z_i(r_i, w), \quad z_i := \frac{r_i^2}{2b_i w^2} (b_i w - \sin b_i w),$$

and we split the last column of J as a sum

$$\begin{pmatrix} W_1 \\ W_2 \\ \vdots \\ W_\ell \\ \partial_w z \end{pmatrix} = \begin{pmatrix} W_1 \\ 0 \\ \vdots \\ 0 \\ \partial_w z_1 \end{pmatrix} + \begin{pmatrix} 0 \\ W_2 \\ \vdots \\ 0 \\ \partial_w z_2 \end{pmatrix} + \dots + \begin{pmatrix} 0 \\ 0 \\ \vdots \\ W_\ell \\ \partial_w z_\ell \end{pmatrix}. \quad (47)$$

Notice that in the i -th column only the i -th variables appear.

By multilinearity of determinant, $\det J$ is the sum of the determinants of ℓ matrices, obtained by replacing each time last column with one of vectors appearing in the sum (47). If we replace it, for instance, with the first term, we get

$$J_1 = \begin{pmatrix} Q_1 & & & W_1 \\ & Q_2 & & 0 \\ & & \ddots & \vdots \\ & & & Q_\ell & 0 \\ Z_1 & Z_2 & \dots & Z_\ell & \partial_w z_1 \end{pmatrix}$$

Now with straightforward computations (notice that $\partial_{\theta} z_1 = 0$ in Z_1), we get

$$\begin{aligned}\det J_1 &= \det Q_2 \cdots \det Q_\ell \cdot \det \begin{pmatrix} Q_1 & W_1 \\ Z_1 & \partial_w z_1 \end{pmatrix} \\ &= \det Q_2 \cdots \det Q_\ell \cdot \det \begin{pmatrix} Q_1^r & Q_1^\theta & W_1 \\ \partial_{r_1} z_1 & 0 & \partial_w z_1 \end{pmatrix}.\end{aligned}$$

Setting

$$A_i = (Q_i^\theta \quad W_i) \quad i = 1, \dots, \ell,$$

we find

$$\det J_1 = \det Q_2 \cdots \det Q_\ell \cdot (\partial_w z_1 \det Q_1 + \partial_{r_1} z_1 \det A_1).$$

Similarly, we find analogous expressions for J_2, \dots, J_ℓ . Then

$$\det J = \sum_{i=1}^{\ell} \det J_i = \sum_{i=1}^{\ell} \left(\prod_{j \neq i} \det Q_j \right) (\partial_w z_i \det Q_i + \partial_{r_i} z_i \det A_i).$$

From (42), by direct computations, it follows

$$\det Q_j = \frac{4r_j}{b_j^2 w^2} \sin^2\left(\frac{b_j w}{2}\right),$$

Moreover, with some computations, one can get

$$\begin{aligned}\partial_w z_i &= \frac{r_i^2}{w^2} \sin^2\left(\frac{b_i w}{2}\right) - \frac{r_i^2}{b_i w^3} (b_i w - \sin b_i w) \\ \partial_{r_i} z_i &= -\frac{r_i}{b_i w^2} (b_i w - \sin b_i w) \\ \det A_i &= \frac{4r_i^2}{b_i^2 w^3} \sin\left(\frac{b_i w}{2}\right) \left(\frac{b_i w}{2} \cos\left(\frac{b_i w}{2}\right) - \sin\left(\frac{b_i w}{2}\right) \right)\end{aligned}$$

and finally, after some simplifications

$$\partial_w z_i \det Q_i + \partial_{r_i} z_i \det A_i = -\frac{4r_i^3}{b_i^2 w^4} \sin\left(\frac{b_i w}{2}\right) \left(\frac{b_i w}{2} \cos\left(\frac{b_i w}{2}\right) - \sin\left(\frac{b_i w}{2}\right) \right).$$

from which we get (44). □

From the explicit expression of the Jacobian (44) we see that integration with respect to horizontal variables (r_i, θ_i) does not involve frequencies, providing a constant C_ℓ . Hence, the computation of the volume reduces to a one dimensional integral in the vertical variable w :

$$V = \int_{-\frac{2\pi}{\max b_i}}^{\frac{2\pi}{\max b_i}} \frac{C_\ell}{B^2 w^{2\ell+2}} \sum_{i=1}^{\ell} \left(\prod_{j \neq i} \sin^2\left(\frac{b_j w}{2}\right) \right) \sin\left(\frac{b_i w}{2}\right) \left(\frac{b_i w}{2} \cos\left(\frac{b_i w}{2}\right) - \sin\left(\frac{b_i w}{2}\right) \right) dw.$$

Using symmetry property of the jacobian with respect to w and making the change of variable $2s = w$, the volume become (reabsorbing all constants in C_ℓ)

$$V = \int_0^{\frac{\pi}{\max b_i}} \frac{C_\ell}{B^2 s^{2\ell+2}} \sum_{i=1}^{\ell} \left(\prod_{j \neq i} \sin^2(b_j s) \right) \sin(b_i s) (b_i s \cos(b_i s) - \sin(b_i s)) ds. \quad (48)$$

Remark 39. In the case $\ell = 1$ (Heisenberg group) we get

$$V = \frac{1}{12}(1 + 2\pi \operatorname{Si}(2\pi)) \simeq 0.8258, \quad \operatorname{Si}(x) := \int_0^x \frac{\sin t}{t} dt.$$

Notice that this is precisely the value of the constant $f_{\mathcal{P}\mathcal{S}}$, since in this case Popp's measure coincide with the Lebesgue measure in our coordinates.

5.3 Differentiability properties: contact case

Let us come back to the differentiability of the map

$$q \mapsto \mathcal{L}(\widehat{B}_q). \quad (49)$$

A smooth family of sub-Riemannian structures is represented by a smooth family of skew symmetric matrices $L(q)$ (see Section 5.1). Recall that, for a smooth family of skew-symmetric matrices that depend on a n -dimensional parameter, the eigenvalue functions $q \mapsto b_i(q)$ exists and are Lipschitz continuous with respect to q (see [30]).

Thus, if we denote with $V(q)$ the volume of the nilpotent unit ball corresponding to frequencies $b_1(q), \dots, b_\ell(q)$, formula (48) implies that

$$V(q) = \int_0^{a(q)} G(q, s) ds, \quad (50)$$

where

$$a(q) := \frac{\pi}{\max b_i(q)},$$

$$G(q, s) := \frac{1}{s^{2\ell+2} B^2} \sum_{i=1}^{\ell} \left(\prod_{j \neq i} \sin^2(b_j(q)s) \right) \sin(b_i(q)s) (b_i(q)s \cos(b_i(q)s) - \sin(b_i(q)s)). \quad (51)$$

Note. In the following we drop the constant C_ℓ that appear in (48) since it does not affect differentiability of the volume.

Remark 40. Since the family of sub-Riemannian structures $q \mapsto L(q)$ is smooth, the exponential map smoothly depends on the point q . As a consequence the integrand $G(q, s)$, being the Jacobian of the exponential map, is a smooth function of its variables.

In addition, although $q \rightarrow b_i(q)$ is only Lipschitz, it is easy to see that the function $a(q) = \max b_i(q)$ is semiconvex with respect to q (see also [18]). In particular $a(q)$ admits second derivative almost everywhere.

If for all q the moduli of the eigenvalues of $L(q)$ are simple (i.e. there are no resonance points), then the eigenvalue functions can be chosen in a smooth way. As a consequence the volume $V(q)$ is smooth, since all functions that appear in (50) are smooth. This argument provides a proof of Proposition 6 in the case of contact structures.

On the other hand, we prove that, along a curve where $a(q)$ is not smooth, i.e. when the two bigger frequencies cross, $V(q)$ is no longer smooth at that point.

Proof of Theorem 5. In this proof by a resonance point we mean a point q_0 where a is not smooth, i.e. where the two (or more) bigger eigenvalues of $L(q_0)$ coincide. As we noticed, V is smooth at non-resonance points.

We divide the proof into two steps: first we prove that $V \in \mathcal{C}^3$ and then we show that, in general, it is not smooth (but \mathcal{C}^4 on curves).

(i). We have to show that V is \mathcal{C}^3 at every resonance point q_0 . First split the volume as follows

$$V(q) = \int_0^{a(q_0)} G(q, s) ds + \int_{a(q_0)}^{a(q)} G(q, s) ds \quad (52)$$

The first term in the sum is smooth with respect to q , since it is the integral of a smooth function on a domain of integration that does not depend on q . We are then reduced to the regularity of the function

$$W(q) := \int_{a(q_0)}^{a(q)} G(q, s) ds \quad (53)$$

This lemma shows that the first three derivatives of W exist and vanish at q_0 .

Lemma 41. *Let $q_0 \in M$ be a resonance point. Then*

$$\left| \int_{a(q_0)}^{a(q)} G(q, s) ds \right| \leq C|q - q_0|^4 \quad (54)$$

Proof. It is sufficient to prove that every derivative up to second order of G vanish at $(q_0, a(q_0))$. Indeed, being G smooth, computing its Taylor polynomial at $(q_0, a(q_0))$ only terms with order greater or equal than three appear (both in $q - q_0$ and $s - a(q_0)$). Thus, integrating with respect to s and using that $|a(q) - a(q_0)| = O(|q - q_0|)$, we have the desired result.

From the explicit formula (51) it is easy to see that $G(q, a(q)) \equiv 0$ for every $q \in M$. In particular $G(q_0, a(q_0)) = 0$. Moreover, since at a resonance point q_0 at least the two bigger eigenvalues coincide, say $b_1(q_0) = b_2(q_0) = \beta$, we have $a(q_0) = \pi/\beta$ and in a neighborhood of $(q, s) = (q_0, a(q_0))$

$$\sin^2(b_1(q) \frac{\pi}{a(q)}) \sin(b_2(q) \frac{\pi}{a(q)}) = O(|b_1 - b_2|^3) = O(|q - q_0|^3). \quad (55)$$

due to the Lipschitz property of $b_j(q)$, for $j = 1, 2$.

From (51) and (55) one can easily get that every derivative of G up to second order (in both variables q and s) vanish at $(q_0, a(q_0))$. □

To complete the proof it is sufficient to show that the first three derivatives of W are small at points close to q_0 . For simplicity of the notation, we will denote by $\frac{\partial W}{\partial q}$ the partial derivative with respect to some coordinate function on M (recall that $q \in M$, where $\dim M = n$). At non resonance points q near q_0 we can compute

$$\begin{aligned} \frac{\partial W}{\partial q}(q) &= \underbrace{G(q, a(q))}_{=0} \frac{\partial a}{\partial q}(q) + \int_{a(q_0)}^{a(q)} \frac{\partial G}{\partial q}(q, s) ds \\ &= \int_{a(q_0)}^{a(q)} \frac{\partial G}{\partial q}(q, s) ds \end{aligned}$$

and

$$\frac{\partial^2 W}{\partial q^2}(q) = \frac{\partial G}{\partial q}(q, a(q)) \frac{\partial a}{\partial q}(q) + \int_{a(q_0)}^{a(q)} \frac{\partial^2 G}{\partial q^2}(q, s) ds \quad (56)$$

which shows that first and second derivatives are continuous (and vanish) at q_0 since $\frac{\partial G}{\partial q}(q_0, a(q_0)) = 0$ (see proof of Lemma 41) and $\frac{\partial a}{\partial q}$ is bounded. In analogous way one can compute for the third derivative

$$\begin{aligned} \frac{\partial^3 W}{\partial q^3}(q) &= \frac{\partial^2 G}{\partial q \partial s}(q, a(q)) \left(\frac{\partial a}{\partial q}(q) \right)^2 + 2 \frac{\partial^2 G}{\partial q^2}(q, a(q)) \frac{\partial a}{\partial q}(q) + \\ &+ \frac{\partial G}{\partial q}(q, a(q)) \frac{\partial^2 a}{\partial q^2}(q) + \int_{a(q_0)}^{a(q)} \frac{\partial^3 G}{\partial q^3}(q, s) ds \end{aligned} \quad (57)$$

Using again that every second derivative of G vanish at $(q_0, a(q_0))$ and that $\frac{\partial a}{\partial q}$ is bounded, it remains to check only the continuity of the term where the second derivative $\frac{\partial^2 a}{\partial q^2}$ appear. From the proof of Lemma 41 one can see that $\frac{\partial G}{\partial q} = O(|b_j - b_k|^2)$, where b_j and b_k denotes the maximum eigenvalues (see also (55)). Hence it is sufficient to prove that $\frac{\partial^2 a}{\partial q^2} = O(1/|b_j - b_k|)$. This is a consequence of the following lemma, written for Hermitian matrices, that can be easily readapted to estimate the growth of the modulus of the eigenvalues of a skew-symmetric family. Indeed, for every skew-symmetric matrix A , with eigenvalues $\pm i\lambda_j$, with $j = 1, \dots, n$, the matrix iA is an Hermitian matrix which has eigenvalues $\pm\lambda_j$, $j = 1, \dots, n$.

Lemma 42. *Let A, B be two $n \times n$ Hermitian matrices and assume that for every t the matrix $A + tB$ has a simple eigenvalue $\lambda_j(t)$. Then the the following equation is satisfied*

$$\ddot{\lambda}_j = 2 \sum_{k \neq j} \frac{|\langle Bx_j, x_k \rangle|^2}{\lambda_j - \lambda_k} \quad (58)$$

where $\{x_k(t)\}_{k=1, \dots, n}$ is an orthonormal basis of eigenvectors and $x_j(t)$ is the eigenvector associated to $\lambda_j(t)$.

Proof. In this proof we endow \mathbb{C}^n with the standard scalar product $\langle z, w \rangle = \sum_{k=1}^n z_k \overline{w_k}$. Since $\lambda_j(t)$ is simple for every t , both $\lambda_j(t)$ and the associated eigenvector $x_j(t)$ can be chosen smoothly with respect to t . By definition

$$(A + tB)x_j(t) = \lambda_j(t)x_j(t), \quad |x_j(t)| = 1.$$

Then we compute the derivative with respect to t of both sides

$$(A + tB)\dot{x}_j(t) + Bx_j(t) = \dot{\lambda}_j(t)x_j(t) + \lambda_j(t)\dot{x}_j(t), \quad (59)$$

and computing the scalar product with $x_j(t)$ we get

$$\dot{\lambda}_j(t) = \langle Bx_j(t), x_j(t) \rangle, \quad \text{hence} \quad \ddot{\lambda}_j(t) = 2 \operatorname{Re} \langle Bx_j(t), \dot{x}_j(t) \rangle. \quad (60)$$

using that $A + tB$ is Hermitian and $\langle \dot{x}_k(t), x_k(t) \rangle = 0$. On the other hand, the scalar product of (59) with x_k , with $k \neq j$ gives

$$\langle \dot{x}_j, x_k \rangle = \frac{\langle Bx_j, x_k \rangle}{\lambda_j - \lambda_k}.$$

Substituting $\dot{x}_j = \sum_{k=1}^n \langle \dot{x}_j, x_k \rangle x_k$ in (60) we have (58). \square

(ii). Now we study the restriction of the map (49) along any curve on the manifold, and we see that, due to a simmetry property, V is \mathcal{C}^4 on every curve but in general is not \mathcal{C}^5 . Notice that (54) gives no information on the fourth derivative but the fact that it is bounded. Indeed it happens that it is continuous but its value depend on the curve we choose.

From now on, we are left to consider a smooth *one-parametric* family of sub-Riemannian structure, i.e. of skew symmetric matrices.

Remark 43. An analytic family of skew-symmetric matrices $t \mapsto L(t)$ depending on one parameter, can be simultaneously diagonalized (see again [30]), in the sense that there exists an analytic (with respect to t) family of orthogonal changes of coordinates and analytic functions $b_i(t) > 0$ such that

$$L = \begin{pmatrix} B_1(t) & & \\ & \ddots & \\ & & B_\ell(t) \end{pmatrix}, \quad \text{where} \quad B_i(t) := \begin{pmatrix} 0 & -b_i(t) \\ b_i(t) & 0 \end{pmatrix}. \quad (61)$$

In the case of a \mathcal{C}^∞ family $t \mapsto L(t)$, we can apply the previous result to the Taylor polynomial of this family. As a consequence we get an approximate diagonalization for $L(t)$, i.e. for every $N > 0$ there exists a smooth family of orthogonal changes of coordinates and smooth functions $b_i(t) > 0$ such that every entry out of the diagonal in $L(t)$ is $o(t^N)$. Namely

$$L(t) = \begin{pmatrix} B_1(t) & & o(t^N) \\ & \ddots & \\ o(t^N) & & B_\ell(t) \end{pmatrix}, \quad \text{where} \quad B_i(t) := \begin{pmatrix} o(t^N) & -b_i(t) \\ b_i(t) & o(t^N) \end{pmatrix}. \quad (62)$$

Since we are interested, in the study the \mathcal{C}^k regularity of (49), for k finite, in what follows we can ignore higher order terms and assume that $L(t)$ can be diagonalized as in the analytic case (61).

From the general analysis we know that V is \mathcal{C}^3 . To prove that $t \mapsto V(t)$ is actually \mathcal{C}^4 we discuss first the easiest case $\ell = 2$ (i.e. the contact (4,5) case) and then generalize to any ℓ .

(i) Case $\ell = 2$. Assume that $b_1(t), b_2(t)$ cross transversally at $t = 0$ (from the proof it will follow that this assumption is not restrictive). This means that for the volume $V(t)$ we have the expression

$$V(t) = \begin{cases} \int_0^{\frac{\pi}{b_1(t)}} G(t, s) ds, & \text{if } t > 0, \\ \int_0^{\frac{\pi}{b_2(t)}} G(t, s) ds, & \text{if } t < 0, \end{cases} \quad \text{and} \quad \begin{cases} b_1(0) = b_2(0) \\ b'_1(0) \neq b'_2(0) \end{cases} \quad (63)$$

Since the regularity of the volume does not depend on the value $b_1(0) = b_2(0)$, we can make the additional assumption

$$b_i(t) = 1 + t c_i(t), \quad i = 1, 2$$

for some suitable functions $c_1(t), c_2(t)$. Notice that $a'(t)$ is discontinuous at $t = 0$ and the left and right limits are

$$a'_+ := \lim_{t \rightarrow 0^+} a'(t) = -\pi c_1(0), \quad a'_- := \lim_{t \rightarrow 0^-} a'(t) = -\pi c_2(0),$$

From the explicit expression of G it is easy to compute that

$$\begin{aligned}\frac{\partial^3 G}{\partial t^3}(0, a(0)) &= \frac{6}{\pi^2} c_1 c_2 (c_1 + c_2), \\ \frac{\partial^3 G}{\partial t^2 \partial s}(0, a(0)) &= \frac{2}{\pi^3} (c_1^2 + 4c_1 c_2 + c_2^2), \quad \frac{\partial^3 G}{\partial t \partial s^2}(0, a(0)) = \frac{6}{\pi^4} (c_1 + c_2),\end{aligned}$$

where we denote for simplicity $c_i := c_i(0)$. From the fact that second derivatives of G vanish at $(t, s) = (0, a(0))$ (see the proof of Lemma 41) we have that the 4-th derivative of W at $t = 0$ is computed as follows

$$\begin{aligned}\lim_{t \rightarrow 0^+} W^{(4)}(t) &= 3 \frac{\partial^3 G}{\partial t^3} a'_+ + 3 \frac{\partial^3 G}{\partial t^2 \partial s} (a'_+)^2 + \frac{\partial^3 G}{\partial t \partial s^2} (a'_+)^3, \\ \lim_{t \rightarrow 0^-} W^{(4)}(t) &= 3 \frac{\partial^3 G}{\partial t^3} a'_- + 3 \frac{\partial^3 G}{\partial t^2 \partial s} (a'_-)^2 + \frac{\partial^3 G}{\partial t \partial s^2} (a'_-)^3\end{aligned}$$

where W is defined in (53). It is easily checked that $W^{(4)}$ is continuous (but does not vanish!)

$$\lim_{t \rightarrow 0^+} W^{(4)}(t) = \lim_{t \rightarrow 0^-} W^{(4)}(t) = -\frac{12}{\pi} c_1^2 c_2^2$$

The same argument produce an example that, in general, $V(t)$ is not C^5 . Assuming

$$b_i(t) = 1 + t c_i, \quad c_1 \neq c_2 \quad \text{constant}, \quad i = 1, 2,$$

a longer computation, but similar to the one above, shows that

$$\begin{aligned}\lim_{t \rightarrow 0^+} W^{(5)}(t) &= -\frac{2}{\pi} c_1^3 (13c_1^2 - 29c_1 c_2 + 22c_2^2) \\ \lim_{t \rightarrow 0^-} W^{(5)}(t) &= -\frac{2}{\pi} c_2^3 (13c_2^2 - 29c_1 c_2 + 22c_1^2)\end{aligned}$$

and the 5-th derivatives do not coincide.

(ii) General case. We reduce to case (i).

We can write $G(t, s) = \sum_{i=1}^{\ell} G_i(t, s)$ and $V(t) = \sum_{i=1}^{\ell} V_i(t)$ where we set

$$\begin{aligned}G_i(t, s) &:= \frac{1}{s^{2\ell+2}} \left(\prod_{j \neq i} \sin^2(b_j(t)s) \right) \sin(b_i(t)s) (b_i(t)s \cos(b_i(t)s) - \sin(b_i(t)s)), \\ V_i(t) &:= \int_0^{a(t)} G_i(t, s) ds, \quad i = 1, \dots, \ell.\end{aligned} \tag{64}$$

Assume that b_1, b_2 are the bigger frequencies and that they cross at $t = 0$, i.e.

$$\begin{aligned}b_i(t) &< b_2(t) < b_1(t), \quad \forall t < 0, \quad \forall i = 3, \dots, n. \\ b_i(t) &< b_1(t) < b_2(t), \quad \forall t > 0, \quad \forall i = 3, \dots, n.\end{aligned}$$

From the explicit expression above it is easy to recognise that for G_1 and G_2 we can repeat the same argument used in (i). Indeed if we denote with $\tilde{G}(t, s)$ the integrand of the (4, 5) case we can write $G_1 + G_2$ as the product of a smooth function and \tilde{G}

$$G_1(t, s) + G_2(t, s) = \left(\frac{1}{s^{2\ell-4}} \prod_{j=3}^{\ell} \sin^2(b_j(t)s) \right) \tilde{G}(t, s),$$

which implies that that $V_1 + V_2$ is a \mathcal{C}^4 function.

Moreover it is also easy to see that V_3, \dots, V_n are \mathcal{C}^4 . Indeed from the fact that $b_1(t)$ and $b_2(t)$ both appear in \sin^2 terms in $G_i(t, s)$ for $i > 3$, it follows that in this case

$$G_i(t, a(t)) \equiv \frac{\partial G_i}{\partial t}(t, a(t)) \equiv 0, \quad \frac{\partial^2 G_i}{\partial t^2}(0, a(0)) = 0, \quad \frac{\partial^3 G_i}{\partial t^3}(0, a(0)) = 0, \quad i = 3, \dots, n,$$

and we can apply the same argument used in (i) to the function $V_i'(t) = \int_0^{a(t)} \frac{\partial G_i}{\partial t}(t, s) ds$, for $i = 3, \dots, n$, showing that it is \mathcal{C}^3 , that means $V_i \in \mathcal{C}^4$. \square

Remark 44. As we said the value of the 4-th derivative depend on the curve we choose, hence we cannot conclude that $V \in \mathcal{C}^4$ in general. Nevertheless we explicitly proved that $V \notin \mathcal{C}^\infty$ since in general is not \mathcal{C}^5 , even when restricted on curves.

Moreover from the proof it also follows that, if more than two frequencies coincide at some point (for instance if we get a triple eigenvalue), we have a higher order regularity for every V_i , and the regularity of V increases.

5.4 Extension to the quasi-contact case

Recall that in the quasi contact case the dimension of the distribution is odd and the kernel of the contact form is one dimensional. Hence, applying the same argument used in Section 5.1, we can always normalize the matrix L in the following form:

$$L = \begin{pmatrix} B_1 & & & \\ & \ddots & & \\ & & B_\ell & \\ & & & 0 \end{pmatrix}, \quad B_i := \begin{pmatrix} 0 & -b_i \\ b_i & 0 \end{pmatrix}, \quad b_i > 0.$$

In other words we can select a basis $\{X_1, \dots, X_\ell, Y_1, \dots, Y_\ell, K, Z\}$ such that

$$\begin{cases} \Delta = \text{span}\{X_1, \dots, X_\ell, Y_1, \dots, Y_\ell, K\}, \\ [X_i, Y_i] = -b_i Z, & i = 1, \dots, \ell \\ [X_i, Y_j] = 0, & i \neq j \\ [X_i, K] = [Y_i, K] = 0, & i = 1, \dots, \ell \\ [X_i, Z] = [Y_i, Z] = 0, & i = 1, \dots, \ell \end{cases} \quad (65)$$

where the new vector field K is in the kernel of the bracket mapping, i.e. it commutes with all others elements. Since abnormal extremals are never optimal in quasi contact case (see Remark 14), we are reduced to compute the exponential map to find geodesics. With analogous computations of contact case we get the following expression for the exponential map from the origin

$$\begin{aligned} x_i(t) &= \frac{r_i}{b_i w} (\cos(b_i w t + \theta_i) - \cos \theta_i), \\ y_i(t) &= \frac{r_i}{b_i w} (\sin(b_i w t + \theta_i) - \sin \theta_i), \\ x_{2\ell+1}(t) &= u_{2\ell+1} t, \\ z(t) &= \frac{1}{2w^2} (|r|^2 w t - \sum_i \frac{r_i^2}{b_i} \sin b_i w t). \end{aligned} \quad (66)$$

From (66) it is easily seen that the jacobian of exponential map has exactly the same expression as in contact case (42). Since zero is always an eigenvalue of L , but is never the maximum one, we can proceed as in the contact case and all the regularity results extend to this case.

6 Proof of Theorem 4 and extension to general corank 1 case

We start this section with the proof of Theorem 4, after that we extend the result to the general corank 1 case.

Proof of Theorem 4. Let \mathbf{S} be a sub-Riemannian structure such that $\dim M \leq 5$.

(i). If $\mathcal{G}(\mathbf{S}) \neq (4, 5)$. From Theorem 29 we know that at every point $q \in M$, the nilpotent approximation $\widehat{\mathbf{S}}_q$ has a unique normal form, hence by Corollary 31 all nilpotent approximations are isometric. From this property it easily follows that $f_{\mathcal{P}\mathcal{S}}$, the Popp volume of the unit ball, is constant (recall that Popp measure is intrinsic for the sub-Riemannian structure). This also implies that for a smooth volume μ the density $f_{\mu\mathcal{S}}$ is smooth.

(ii). If $\mathcal{G}(\mathbf{S}) = (4, 5)$ by Theorem 29 it is sufficient to consider the case when the family of nilpotent structure has the normal form (16), where $\alpha = \alpha(q)$ depends on the point. Notice that the formula (48) for the volume of the unit ball is still valid, where now $b_1 = 1$ and $b_2 = |\alpha|$.

Theorem 5 proves that the density is \mathcal{C}^3 at points where $|\alpha| > 0$, i.e. in the contact case. We are then reduced to the study of the volume near a point where the eigenvalue α crosses zero. In particular we show that the volume is smooth at these points. Since the eigenvalue α is approaching zero, it is not restrictive to assume $0 \leq |\alpha| < 1$. Let us consider then the function defined on the interval $(-1, 1)$

$$W(\alpha) = \int_0^\pi \frac{1}{\alpha^2 s^6} \left(\sin^2(\alpha s) \sin s (s \cos s - \sin s) + \sin^2 s \sin(\alpha s) (\alpha s \cos(\alpha s) - \sin(\alpha s)) \right) ds. \quad (67)$$

Note that $V(\alpha) = W(|\alpha|)$, where $V(\alpha)$ denotes the volume of the nilpotent ball relative to frequencies 1 and $\alpha < 1$. It is easy to see that both

$$\frac{\sin^2 \alpha s}{\alpha^2} \quad \text{and} \quad \frac{1}{\alpha^2} \sin(\alpha s) (\alpha s \cos(\alpha s) - \sin(\alpha s)) \quad (68)$$

are smooth as functions of α (also at $\alpha = 0$). Hence W is a smooth function (for $\alpha \in (-1, 1)$). Moreover it is easy to see that W it is an even smooth function of α . Thus W it is smooth also as a function of $|\alpha|$, that completes our proof. \square

The same argument applies to prove that the \mathcal{C}^3 regularity holds in the general corank 1 case. Indeed, from (51) and the fact that (68) are smooth functions at $\alpha = 0$, it follows that the integrand $G(q, s)$ is smooth as soon as one of the eigenvalue $b_i(q)$ is different from zero (recall that $b_i \geq 0$ by definition). Since the structure is regular (i.e. the dimension of the flag do not depend on the point) and bracket generating, we have that $\max_i b_i(q) > 0$ for every q , hence the conclusion.

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