CONTACT EQUATIONS, LIPSCHITZ EXTENSIONS AND ISOPERIMETRIC INEQUALITIES

VALENTINO MAGNANI

ABSTRACT. We characterize locally Lipschitz mappings and existence of Lipschitz extensions through a first order nonlinear system of PDEs. We extend this study to graded group-valued Lipschitz mappings defined on compact Riemannian manifolds. Through a simple application, we emphasize the connection between these PDEs and the Rumin complex. We introduce a class of 2-step groups, satisfying some abstract geometric conditions and we show that Lipschitz mappings taking values in these groups and defined on subsets of the plane admit Lipschitz extensions. We present several examples of these groups, called Allcock groups, observing that their horizontal distribution may have any codimension. Finally, we show how these Lipschitz extensions theorems lead us to quadratic isoperimetric inequalities in all Allcock groups.

CONTENTS

1. Introduction 2
2. Preliminaries and known results 7
  2.1. Some elementary facts on graded groups 7
  2.2. Haar measure and Fubini’s theorem 8
  2.3. Differentiability 10
  2.4. Some auxiliary results 11
3. Technical lemmata 12
4. Group-valued mappings on Riemannian manifolds 15
  5. The differential characterization 18
  5.1. Smooth functions that are not locally Lipschitz 20
  5.2. Contact equations and Rumin complex 21
  5.3. Lipschitz mappings as boundary value problems 23
6. Allcock groups 25
7. Lipschitz extensions from the plane to Allcock groups 30
8. Quadratic isoperimetric inequalities 38
References 44
1. Introduction

The local Lipschitz property for mappings between Euclidean spaces is characterized by local $L^\infty$-bounds on distributional derivatives. If we replace the target space even with the simplest sub-Riemannian manifold, as the Heisenberg group, then the previous statement does not hold. This elementary fact is the starting point of our study.

An important instance is the case of mappings defined on a Riemannian manifold, that is related to the study of horizontal submanifolds. In this connection, M. Gromov among other results has treated various Lipschitz approximation theorems along with Lipschitz extensions problems, see Section 3.5 of [18].

It is well known that smooth Lipschitz mappings taking values into a contact Carnot-Carathéodory manifold are horizontal and the converse also holds, [18]. The main concern of this work is to understand this fact under Lipschitz regularity, where both source and target spaces have the sub-Riemannian structure of stratified groups. The horizontality condition for mappings of stratified groups yields the system of nonlinear first order PDEs (1), that only depends on the Lie algebra and that it can be written in terms of exponential coordinates.

For general sub-Riemannian structures, these equations can be replaced by the pull-back of forms defining the horizontal distribution, according to Section 4.2 of [18]. Recall that in the case the target of our mappings is Euclidean, then they are all horizontal and their Lipschitz property is characterized in the usual way by bounds on distributional horizontal derivatives. This fact holds for general Carnot-Carathéodory source spaces, [15].

Our point of view is that of considering mappings of stratified groups as solutions to the system of equations (1). Although the main results of this work are Theorem 1.1 and Theorem 1.2, a substantial part of the paper is devoted to applications, regarding construction of both Lipschitz and non-Lipschitz mappings and the relationship between Lipschitz extension theorems and isoperimetric inequalities.

The interest in geometric properties of mappings in the sub-Riemannian setting has recently proved useful in connection with bi-Lipschitz embeddability of sub-Riemannian metric spaces into different classes of infinite dimensional Banach spaces, according to the remarkable work by J. Cheeger and B. Kleiner, [10].

Our framework is that of graded groups, that are real, finite dimensional, connected, simply connected and nilpotent Lie groups, with graded Lie algebra. Notice that they might not be connected by rectifiable curves, according to Example 2.1. When the Lie subalgebra spanned by the first layer of the grading coincides with the whole algebra, we say that the group is stratified, [14]. The so-called horizontal directions of the group are spanned by the left invariant vector fields belonging to the first layer of the algebra, see Section 2 for more details. Stratified groups, also called Carnot groups, represent the foremost models of nilpotent and simply connected sub-Riemannian geometries. As we will see below, our techniques will allow for studying mappings from a stratified group to a graded group.

The above mentioned horizontality for mappings of graded groups is well known as contact property. This property corresponds to preserving horizontal directions.
In Remark 2.16, we derive equations (1) that correspond to this constraint. We will refer to these equations as contact equations. For instance, every parametrization of a horizontal curve has the contact property by definition, hence the property of being horizontal can be equivalently stated in terms of contact equations, that in this case become ODEs (16). More generally, parametrizations of either Legendrian submanifolds in Heisenberg groups or horizontal submanifolds in stratified groups are characterized by being solutions to (1), where horizontal derivatives $X_j$ are replaced by usual partial derivatives. Incidentally, these manifolds coincide with $(\mathbb{R}^k,\mathcal{M})$-regular sets, according to [25]. Quasiconformal mappings of a stratified groups have the contact property, [28], hence they solve the system (1) a.e. For instance, contact equations have been implicitly used in [7], in relation with smoothness of 1-quasiconformal mappings between Carnot groups. In connection with mappings with bounded distortion in two step groups, contact equations explicitly appear in (3.2) of [11]. Certainly, many other interesting cases could be added from the existing literature.

In the present paper, we focus our attention on mappings defined on an open subset $\Omega$ of a stratified group $G$, with graded group target $\mathcal{M}$. Lie algebras of $G$ and $\mathcal{M}$ will be denoted by $\mathfrak{g}$ and $\mathcal{M}$, respectively. Turning to the initial question of characterizing the local Lipschitz property, we state the following

**Theorem 1.1.** Let $\Omega \subset G$ be an open set and let $f : \Omega \to \mathcal{M}$. Consider the mappings $F = \exp^{-1} \circ f$ and $F_j = \pi_j \circ F$, hence $f = \exp \sum_{j=1}^v F_j$ and $F_j$ take values in $W_j$. If $F \in L^1_{\text{loc}}(\Omega, \mathcal{M})$, then the following statements are equivalent:

- $f$ is locally Lipschitz,

- all the distributional derivatives $\nabla_{X_i} F$ belong to $L^\infty_{\text{loc}}(\Omega, \mathcal{M})$ and the system

$$
\nabla_H F_j = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j ([F, \nabla_H F]_{n-1})
$$

$\mu$-a.e. holds in $\Omega$ for every $j = 2, \ldots, v$,

- $f$ is a.e. $P$-differentiable and $\nabla_H F = (\nabla_{X_1} F, \ldots, \nabla_{X_m} F) \in L^\infty_{\text{loc}}(\Omega, \mathcal{M})^m$,

where $(X_1, \ldots, X_m)$ is a basis of $V_1$ and $\mu$ is the Haar measure of $G$. The mappings $\pi_j : \mathcal{M} \to W_j$ indicate the canonical projection onto the $j$-th layer of the graded algebra $\mathcal{M}$ of $\mathcal{M}$.

Under the stronger assumption that $\mathcal{M}$ is stratified, results of [28] imply that the local Lipschitz property yields both the second and the third condition of Theorem 1.1. Since a graded group need not be connected by rectifiable curves, then the previously mentioned implications are obtained employing some technical tools developed in [25].

However, in this work we are interested in the converse to these implications, since we aim to achieve the Lipschitz property starting from the contact property. Clearly, the standard smoothing argument used in Euclidean spaces to get the Lipschitz property cannot be applied here. In fact, the mapping $F$ is a weak solution to (1), but its mollification might no longer be a solution of these equations, due to their nonlinearity. The direct use of Pansu differentiability does not seem to be of help, since this notion has no corresponding distributional version. Here our point is that Theorem 1.1 allows
us to think of (1) somehow as a distributional counterpart of Pansu differentiability. Our approach boils down to making the problem a.e. one dimensional. In fact, due to the distributional validity of (1) and a Fubini’s decomposition, one shows that for a.e. horizontal line in the domain, the restriction of the mapping to this line is horizontal. Then these restrictions are Lipschitz continuous with uniform Lipschitz constant with respect to the homogeneous distance of the target. This argument can be iterated for a basis of horizontal directions. Since points are locally connected by piecewise horizontal lines and any horizontal curve can be approximated by piecewise horizontal lines, the local Lipschitz property follows. Constructing globally Lipschitz functions on a domain $\Omega$ of a stratified group clearly depends on the geometry of $\Omega$. In this respect, every mapping that satisfies one of the equivalent conditions in Theorem 1.1 and also $\nabla_H F \in L^\infty(\Omega, \mathcal{M})^n$ is Lipschitz in $\Omega$, if we assume that this set is a John domain.

The approach adopted to get Theorem 1.1 also works replacing the open set of a stratified group with a complete Riemannian manifold, since one replaces horizontal lines with Riemannian geodesics, that are more manageable than sub-Riemannian’s. By Theorem 4.5 and taking into account Remark 4.6, we are lead to the following

**Theorem 1.2.** Let $N$ be a compact connected Riemannian manifold and let $\nu$ be the canonical Riemannian measure of $N$. Let $F : N \to \mathcal{M}$ be a Lipschitz mapping that $\nu$-a.e. satisfies

$$
\nabla F_j = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j ([F, \nabla F]_{n-1})\quad \text{for every } j = 2, \ldots, v.
$$

Then $f : N \to \mathcal{M}$, where $f = \exp \circ F$, is Lipschitz and there exists a geometric constant $C > 0$ such that $\text{Lip}(f) \leq C \text{Lip}(F_1)$.

A special instance of this theorem will be used in the proof of Theorem 1.3, about the existence of Lipschitz extension. Next, we discuss some applications of Theorem 1.1. In fact, this theorem provides a general PDEs approach to construct either Lipschitz or non-Lipschitz mappings. For instance, it is easy to construct smooth mappings in the Heisenberg group that are nowhere locally Lipschitz: it suffices to consider the parametrization of the vertical line in the first Heisenberg group. In Subsection 5.1, through contact equations, we provide another case for a mapping of Heisenberg groups. To construct examples of Lipschitz mappings, we use the natural relationship between contact equations and the Rumin complex. In fact, an elementary computation only relying on contact equations and using the complex property of Rumin differential allows us to determine all smooth Lipschitz mappings of the Heisenberg group, whose horizontal components are affine functions, see Subsection 5.2. Our point here clearly is in the method, that rests on the solution of the simple contact equations specialized to this case. It is also clear how our approach could be extended to either higher dimensional Heisenberg groups or other variants. Recall that affine Lipschitz mappings of the first Heisenberg group have been characterized in [4] through the explicit use of the Lipschitz condition.
Another application of Theorem 1.1 concerns existence of Lipschitz extensions that can be interpreted as existence of solutions to (1) with assigned boundary datum. In general, this PDEs problem is very difficult, already for Euclidean spaces and Heisenberg groups. We will state and prove this characterization in Subsection 5.3. This leads us to two possible methods to approach Lipschitz extension problems in sub-Riemannian geometry: the geometric approach and the PDEs approach.

These two demanding projects cannot be treated here in depth, then we will limit ourselves to show first elementary examples of how one can apply both Theorem 1.1 and Theorem 1.2 to obtain Lipschitz extensions, showing how this issue is essentially equivalent to finding isoperimetric inequalities. We concentrate our attention on the geometric approach, following the method by D. Allcock to get quadratic isoperimetric inequalities in higher dimensional Heisenberg groups. Essentially, contact equations permit us to rephrase Allcock construction in a larger class of two step groups that must satisfy some abstract geometric properties. We call these groups Allcock groups. This leads us to a disk extension theorem in this class of groups, see Theorem 1.3. Notice that M. Gromov pointed out how his disk extension theorem for contact simply connected compact Carnot-Carathéodory target gives a quadratic isoperimetric inequality, see p.218 of [18]. As a final application of the tools developed in this work, we will also prove this implication for all Allcock groups, see Theorem 1.5.

In Section 6, we introduce Allcock groups and present several examples, that include higher dimensional quaternionic H-type groups, the complexified Heisenberg group and other classes of two step groups. It is easy to show that Allcock groups also include Heisenberg groups, see Remark 6.7. Next, we state the following

**Theorem 1.3 (Disk extension theorem).** Let $A^n$ be an Allcock group, with $n \geq 2$. Then there exists a geometric constant $c > 0$ such that for every Lipschitz mapping $f : S^1 \to A^n$ there exists a Lipschitz extension on the closed disk $\tilde{f} : D \to A^n$ such that $\text{Lip}(\tilde{f}) \leq c \text{Lip}(f)$.

It is a rather general fact that Lipschitz extension theorems from spheres to the corresponding higher dimensional disks imply full Lipschitz extension theorems, from an Euclidean source space. The argument of the proof essentially relies on Whitney cube decomposition, according for instance to Theorem 1.2 of [3]. A different statement of this fact can be found for mappings taking values in a compact Carnot-Carathéodory manifold, see p.219 of [18]. In the recent work by Lang and Schlichenmaier, [21], using Nagata dimension, more general source spaces are considered. In Section 7 we use both terminology and results of this paper to show how Theorem 1.3 yields the following

**Corollary 1.4.** Let $A^n$ be an Allcock group, with $n \geq 2$. Then the couple $(\mathbb{R}^2, A^n)$ has the Lipschitz extension property.

Section 8 is devoted to quadratic isoperimetric inequalities for Allcock groups. We wish to stress that here we certainly do not exhibit novel techniques to get isoperimetric inequalities in nilpotent groups, but rather aim to clarify known results, presenting some possibly new cases and eventually support future progress to study this vast and interesting issue by methods of sub-Riemannian Geometric Measure Theory.
Theorem 1.5 (Quadratic isoperimetric inequality). Let $\mathbb{A}_n$ be an Allcock group, with $n \geq 2$. Then there exists a geometric constant $K > 0$ such that for every Lipschitz loop $\Gamma : S^1 \to \mathbb{A}_n$ there exists a Lipschitz map $f : D \to \mathbb{A}_n$ such that $f|_{S^1} = \Gamma$ and the “spanning disk” $(D, f)$ satisfies the quadratic isoperimetric inequality

$$H_{\rho_0}^2(f(D)) \leq K \text{length}_{\rho_0}(\Gamma)^2,$$

where $\rho_0$ is a fixed Carnot-Carathéodory distance.

Recall that quadratic isoperimetric inequalities impose quadratic behaviour to the corresponding “geometric Dehn functions”. It is interesting to point out that there is also a combinatorial notion of Dehn function, that is shown to be equivalent to the geometric one, [5], [6]. This represents a fascinating connection between Combinatorial Group Theory and Geometric Group Theory. To have a glimpse of this vast research area, we mention just a few references, [5], [8], [16], [17].

In fact, quadratic isoperimetric inequalities in Heisenberg groups have also a combinatorial proof, [26]. More recently, R. Young have found various estimates for Dehn functions in some specials classes of stratified groups, [31], [32]. In particular, for some central powers of two step groups he establishes quadratic isoperimetric inequalities using combinatorial methods, [31]. It is interesting to compare his results with our Theorem 1.5 since the construction of central powers is somehow similar to the construction of an Allcock group $\mathbb{A}_n$, starting from its model algebra $n$, see Section 6.

Notice that the core of Allcock isoperimetric inequality is establishing this result in the “symplectic part” $(\mathbb{R}^{2n}, \omega)$ of the Heisenberg group $\mathbb{H}^n$ using the Euclidean Hausdorff measure in $\mathbb{R}^{2n}$. In this respect, he raises the question on the proper notion of area to read isoperimetric inequalities in the Heisenberg group equipped with its Carnot-Carathéodory distance, see at p.230 of [1]. Theorem 1.5 answers this question showing that the 2-dimensional Hausdorff measure with respect to the Carnot-Carathéodory distance works. This is a simple consequence of the sub-Riemannian area formula, [22], that shows how the Euclidean surface measure of the projected surface in $\mathbb{R}^{2n}$ corresponds to its Hausdorff measure with respect to the Carnot-Carathéodory distance in $\mathbb{H}^n$. This works more generally for Allcock groups, see Proposition 8.5.

In this case, we first obtain an Euclidean isoperimetric inequality in the “multi-symplectic space” $(\mathbb{R}^{mn}, \omega)$, where the multi-symplectic form $\omega$ is defined in [33]. In fact, $\mathbb{A}_n$ can be identified with $\mathbb{R}^{mn} \times \mathbb{R}^s$ with respect to suitable graded coordinates and the horizontal subspace given by $\mathbb{R}^{mn}$ inherits $\omega$ from the Lie algebra of $\mathbb{A}_n$. The quadratic isoperimetric inequality in $(\mathbb{R}^{mn}, \omega)$ is then the main point, corresponding to Theorem 8.6 Then the sub-Riemannian area formula leads us to Theorem 1.5.

Since the Riemannian Heisenberg group is quasi-isometric to the sub-Riemannian one, it is natural to expect the same isoperimetric inequality with respect to both distances, according to the case of finitely presented groups, [2]. In fact, under the same assumptions of Theorem 1.5 we have

$$H_{\bar{\rho}}^2(f(D)) \leq K \text{length}_{\bar{\rho}}(\Gamma)^2,$$
where $\varrho$ is the Riemannian distance obtained by the fixed left invariant Riemannian metric defining the Carnot-Carathéodory distance $\rho_0$, see Section 8. In fact, since $\rho_0$ is greater than or equal to $\varrho$ and Proposition 8.2 show that Riemannian and sub-Riemannian lengths of a horizontal curve coincide, it follows that (3) implies (4). Finally, we wish to remark that using the Euclidean Hausdorff measure to study quadratic isoperimetric inequalities for instance in $\mathbb{H}^2$ does not seem promising, according to Example 8.8.

2. Preliminaries and known results

2.1. Some elementary facts on graded groups. A graded group is a real, finite dimensional, connected and simply connected Lie group $M$, whose Lie algebra $\mathcal{M}$ can be written as the direct sum of subspaces $W_i$, called layers, such that

$$[W_i, W_j] \subset W_{i+j}$$

and $\mathcal{M} = W_1 \oplus \cdots \oplus W_v$. The integer $v$ is the step of nilpotence of $\mathbb{M}$. A graded group $\mathbb{M}$ is stratified if its layers satisfy the stronger condition $[W_i, W_j] = W_{i+j}$, see for instance [14]. The horizontal tangent spaces

$$H_x \mathbb{M} = \{ Z(x) \mid Z \in W_1 \} \subset T_x \mathbb{M}, \quad x \in \mathbb{M}$$

define all horizontal directions of the group, that are collected into the so-called horizontal subbundle $H_x \mathbb{M}$. We also define subbundles of higher order $H^j \mathbb{M}$, setting

$$H^j_x \mathbb{M} = \{ Z(x) \mid Z \in W_j \} \subset T_x \mathbb{M}, \quad x \in \mathbb{M}$$

We fix a norm $\| \cdot \|$ in $\mathcal{M}$, then the bilinearity of Lie brackets gives

$$\|[X,Y]\| \leq \beta \|X\| \|Y\| \quad \text{for every } X, Y \in \mathcal{M}$$

for some constant $\beta > 0$, depending on the norm and on the algebra. The grading of $\mathcal{M}$ allows us to introduce a one-parameter group of Lie algebra automorphisms $\delta_r : \mathcal{M} \longrightarrow \mathcal{M}$, defined as $\delta_r(X) = r^i$ if $X \in V_i$, where $r > 0$. These mappings are called dilations. Taking into account that the exponential mapping $\exp : \mathcal{M} \longrightarrow \mathbb{M}$ is a diffeomorphism for simply connected nilpotent Lie groups, we can read dilations in the group $\mathbb{M}$ through the mapping $\exp$ and maintain the same notation. We fix a homogeneous distance $\rho$ on $\mathbb{M}$, namely, a left invariant continuous distance that is 1-homogeneous with respect to dilations $\delta_r$. We will use the convention $\rho(x) = \rho(x,e)$, where $e$ denotes here the unit element of $\mathbb{M}$.

Example 2.1. Let $\mathbb{M} = \mathbb{R}^n \times \mathbb{R}$ be equipped with the sum of vectors as commutative group operation and define the parabolic distance $|x,t| = |x| + \sqrt{|t|}$ and dilations $\delta_r(x,t) = (rx, r^2t)$. Here we have the grading $V_1 \oplus V_2$, where $V_1$ and $V_2$ can be identified with $\mathbb{R}^n \times \{0\}$ and $\{0\} \times \mathbb{R}$, respectively. Clearly, $\mathbb{M}$ is a 2-step graded group, but it is not stratified.

Following the notation of [28], Section 4.5, in the next definition we introduce the iterated Lie bracket.
Definition 2.2. Let $X,Y \in \mathcal{M}$. The $k$-th bracket is defined by
\[
[X,Y]_k = \underbrace{[X,[X,\cdots,[X,Y],],\ldots,]}_{k \text{ times}} \quad \text{and} \quad [X,Y]_0 = Y.
\]
Notice that $[X,Y]_k = \text{ad}^k(Y)$.

We can express the group operation in the Lie algebra by the Baker-Campbell-Hausdorff formula, that we present in the following form
\[
X \circ Y = \sum_{j=1}^\nu c_n(X,Y),
\]
where $c_1(X,Y) = X + Y$, $c_2(X,Y) = [X,Y]/2$ and the subsequent terms can be defined inductively. Notice that the sum has $\nu$ addends, since $X,Y$ belong to $\mathcal{M}$ that has step of nilpotence $\nu$. The addends $c_n$ in general are given by induction through the following Baker-Campbell-Hausdorff-Dynkin formula
\[
(n+1)c_{n+1}(X,Y) = \frac{1}{2} [X - Y, c_n(X,Y)] + \sum_{\substack{p \geq 1 \\ 2p \leq n}} K_{2p} \sum_{k_1,\ldots,k_{2p}>0 \\ k_1+\cdots+k_{2p}=n} [c_{k_1}(X,Y),[\cdots,[c_{k_{2p}}(X,Y),X+Y],\ldots,],
\]
see Lemma 2.15.3 of [33].

Lemma 2.3. Let $\nu > 0$ and let $n = 2,\ldots,\nu$. Then there exists a constant $\alpha_n(\nu)$ only depending on $n$ and $\nu$ such that
\[
\|c_n(X,Y)\| \leq \alpha_n(\nu) \|[X,Y]\|
\]
whenever $\|X\|,\|Y\| \leq \nu$.

Proof. Our statement is trivial for $n = 2$, being $c_2(X,Y) = [X,Y]/2$. Assume that it is true for every $j = 2,\ldots,n$, with $n \geq 2$. We observe that $[c_{k_{2p}}(X,Y),X+Y] \neq 0$ in (9) implies $k_{2p} > 1$, then inductive hypothesis yields
\[
\|c_{k_{2p}}(X,Y)\| \leq \alpha_{k_{2p}}(\nu) \|[X,Y]\|.
\]
Using this estimate in (9) and observing that $\|c_{k_i}(X,Y)\| \leq 2\nu$, whenever $k_i = 1$, our claim follows. \qed

Definition 2.4. A homogeneous subgroup $H$ of $\mathbb{G}$ is a Lie subgroup that is closed under dilations. Analogously, for subalgebras the same terminology is adopted.

Definition 2.5. Let $N$ and $H$ be homogeneous subgroups of $\mathbb{G}$, where $N$ is normal, $N \cap H = \{e\}$ and $NH = \mathbb{G}$. Then $\mathbb{G}$ is an inner semidirect product of $N$ and $H$ and we write $\mathbb{G} = N \rtimes H$.

We denote by $H_X$ the one-dimensional subgroup of $\mathbb{G}$, spanned by $\exp X$, where $X \in \mathcal{G}$. Next, we recall a standard fact concerning direct sums of homogeneous subalgebras, whose proof can be found for instance in [25]. Recall that a homogeneous algebra $\mathfrak{p}$ satisfies $\delta_r \mathfrak{p} \subset \mathfrak{p}$ for every $r > 0$. 
Proposition 2.6. Let $p$ and $h$ be homogeneous subalgebras of $G$ and let $P$ and $H$ denote their corresponding homogeneous subgroups, respectively. Then the condition $p \oplus h = G$ is equivalent to require that $P \cap H = \{e\}$ and $PH = G$. Furthermore, if one of these conditions hold, then the mapping
\begin{equation}
\phi : p \times h \longrightarrow G, \quad \phi(W, Y) = \exp W \exp Y
\end{equation}
is a diffeomorphism.

Lemma 2.7. Let $X \in V_1 \setminus \{0\}$. Then there exists a normal homogeneous subgroup $N \subset G$ such that $G = N \rtimes H_X$.

Proof. Let $U_1 \subset V_1$ be a subspace of dimension $\dim V_1 - 1$ that does not contain $X$ and set $\mathcal{N} = U_1 \oplus V_2 \oplus \cdots \oplus V_i$. We notice that $\mathcal{N}$ is a homogeneous ideal, hence $N = \exp \mathcal{N}$ is a normal homogeneous subgroup of $G$. By definition of $\mathcal{N}$, we have $\mathcal{N} \cap \text{span}\{X\} = \{0\}$ and this implies that $N \cap H_X = \{e\}$ and $NH_X = G$, see for instance Proposition 2.6. In other words $G = N \rtimes H_X$. \qed

2.2. Haar measure and Fubini’s theorem. We consider a graded group $G$ with grading $G = V_1 \oplus \cdots \oplus V_\infty$ and we fix a left invariant Riemannian metric $g$ on $G$. Then the associated volume measure $\text{vol}_g$ is clearly left invariant and defines the Haar measure of $G$. We will denote by $\mu$ this measure. Let $X \in V_1$ and $N$ be a homogeneous normal subgroup $N$ such that $G = N \rtimes H_X$. Then we have the following Fubini’s theorem with respect to this factorization.

Proposition 2.8. Let $\mu$ be the Haar measure of $G$. Then for every measurable set $A \subset G$, we have
\begin{equation}
\mu(A) = \int_N \nu_X(A_n) \, d\mu_N(n),
\end{equation}
where $A_n = \{h \in H_X \mid nh \in A\}$. We have denoted by $\mu_N$ and $\nu_X$ the Haar measure of $N$ and of $H_X$, respectively.

Proof. We fix an orthonormal basis $(X_1, \ldots, X_q)$ of $G$ with respect to the metric $g$. In addition we assume that this basis is adapted to the grading of $G$, such that $X_1$ is proportional to $X$. By Proposition 2.6, the mapping
\begin{equation}
\psi : \mathbb{R} \times \mathbb{R}^{q-1} \longrightarrow M, \quad (\xi, t) \longrightarrow \exp \left( \sum_{j=2}^q x_j X_j \right) \exp (x_1 X_1)
\end{equation}
is a diffeomorphism. Our fixed basis also introduces the special system of coordinates $F : \mathbb{R}^q \longrightarrow G$, $F(x) = \exp \left( \sum_{i=1}^q x_i X_i \right)$ on $G$. Now we observe that
\begin{equation}
F_i \mathcal{L}^q = \mu, \quad (F_{[0]} \times \mathbb{R}^{q-1})_t \mathcal{L}^{q-1} = \mu_N, \quad \text{and} \quad (F_{\mathbb{R} \times [0]})_t \mathcal{L}^1 = \nu_X,
\end{equation}
see for instance Proposition 2.3.47 of [24]. Thus, the fact that $X_1 \in V_1$ implies that the mapping $F^{-1} \circ \psi$ has jacobian equal to one. Combining these facts with classical Fubini’s theorem, we get our claim. \qed
2.3. Differentiability. In the present subsection, we recall the notion of Pansu differentiability. \(G\) and \(M\) denote two stratified groups and \(\Omega\) is an open subset of \(G\).

**Definition 2.9** (h-homomorphism). A group homomorphism \(L : G \rightarrow M\) such that 
\(L(\delta^r_x) = \delta^r_M L(x)\) for every \(x \in G\) and \(r > 0\) is called homogeneous homomorphism, in short h-homomorphism.

Analogous terminology will be used for the corresponding Lie algebra homomorphisms of graded algebras that commute with dilations.

**Definition 2.10** (P-differentiability). Let \(d\) and \(\rho\) be homogeneous distances of \(G\) and \(M\), respectively. We consider the mapping \(f : \Omega \rightarrow M\). We say that \(f\) is P-differentiable at \(x \in \Omega\) if there exists an h-homomorphism \(L : G \rightarrow M\) such that
\[
\rho\left(f(x)^{-1}f(xh), L(h)\right) \rightarrow 0 \quad \text{as} \quad h \rightarrow e.
\]
The h-homomorphism \(L\) satisfying this limit is unique and it is called P-differential of \(f\) at \(x\). We denote \(L\) by \(Df(x)\), when we read the P-differential between the corresponding Lie algebras, we will denote it by \(df(x)\).

**Theorem 2.11.** Every Lipschitz mapping \(f : \Omega \rightarrow M\) is \(\mu\)-a.e. P-differentiable.

This theorem is an important result due to Pansu, [28]. Here have presented a slightly more general version where \(M\) is graded, but it might not be stratified, [25].

**Definition 2.12** (Distributional derivatives). Let \(\Omega\) be an open subset of a stratified group \(G\), let \(X\) be a left invariant vector field of \(G\) and let \(E\) be a finite dimensional normed space. Then for every \(F \in L^1_{\text{loc}}(\Omega, E)\) we say that \(G \in L^1_{\text{loc}}(\Omega, E)\) is the distributional derivative of \(F\) with respect to \(X\) if
\[
\int_{\Omega} F X \varphi \, d\mu = - \int_{\Omega} G \varphi \, d\mu
\]
for every \(\varphi \in C^\infty_c(\Omega)\). Uniqueness of \(G\) allows us to use the notation \(\nabla_X F\). In the case \(\Omega\) is an open subset of \(\mathbb{R}^n\) and \(X = \partial_x\), we will use the notation \(\nabla_j\), where \(e_j\) belongs to the canonical basis of \(\mathbb{R}^n\).

**Remark 2.13.** To avoid confusion, we stress that the symbol \(\nabla_X\) will always denote a distributional derivative, since we will never consider connections in this work.

**Definition 2.14** (Horizontal gradient). Under the conditions of Definition 2.12 if we have equipped \(V_1\) with a scalar product and \((X_1, \ldots, X_m)\) is an orthonormal basis of \(V_1\), then we introduce the notation
\[
\nabla_H F = (\nabla_{X_1} F, \ldots, \nabla_{X_m} F)
\]
to denote the distributional horizontal gradient of \(F : \Omega \rightarrow E\).
2.4. Some auxiliary results. In this subsection $\mathbb{M}$ denotes a graded group equipped with Lie algebra $\mathcal{M}$. The next lemma easily follows from the scaling property of the homogeneous distance in a graded group.

**Lemma 2.15.** Let $\mathcal{M} = W_1 \oplus \cdots \oplus W_\upsilon$ and let $\pi^i : \mathcal{M} \longrightarrow W_i \oplus \cdots \oplus W_\upsilon$ be the canonical projection. Let $U$ be a bounded open neighbourhood of the unit element $e \in \mathbb{M}$. Then there exists a constant $K_U > 0$, depending on $U$, such that

$$\|\pi^i(\exp^{-1}(x))\| \leq K_U d(x)^i \quad (13)$$

holds for every $x \in U$ and every $i = 1, \ldots, \upsilon$.

**Remark 2.16** (Contact property). We wish to point out how the contact property of a mapping $f : \Omega \longrightarrow \mathbb{M}$ at some point $x \in \Omega$, namely,

$$df(x)(H_x G) \subset H_{f(x)}\mathbb{M}$$

is equivalent to the differential constraint

$$(14) \quad X_i F_j = \sum_{n=2}^\upsilon \frac{(-1)^n}{n!} \pi_j([F_{j-1}(x), X_i F_{j-1}(x)]_{n-1}) = 0,$$

for every $i = 1, \ldots, m$ and every $j = 2, \ldots, \upsilon$, where $(X_1, \ldots, X_m)$ is a basis of $V_1$, $f = \exp \circ F$ and $F_j = \pi_j \circ F$. The exponential mapping $\exp : \mathcal{M} \longrightarrow \mathbb{M}$ satisfies

$$(15) \quad d\exp(X) = \text{Id} - \sum_{n=2}^\upsilon \frac{(-1)^n}{n!} \text{ad}(X)^{n-1},$$

see Theorem 2.14.3 of [33]. Identifying for every $y \in \mathbb{M}$ the tangent space $T_y\mathbb{M}$ with $\mathcal{M}$ and applying formula (15), we have

$$X_i F(x) = X_i F(x) - \sum_{n=2}^\upsilon \frac{(-1)^n}{n!} \text{ad}(F(x))^{n-1}(X_i F(x)).$$

Then $X_i F(x) \in H_{f(x)}\mathbb{M}$ if and only if

$$\pi_j \left( X_i F(x) - \sum_{n=2}^\upsilon \frac{(-1)^n}{n!} [F(x), X_i F(x)]_{n-1} \right) = 0 \quad \text{for all} \quad j \geq 2.$$ 

This proves the characterizing property of equations (14). We also notice that this characterization holds even if the mapping is differentiable along horizontal directions, see [25] for the precise definition of horizontal differentiability.

The following result corresponds to Corollary 5.4 of [25].

**Theorem 2.17.** Let $\Gamma : [a, b] \longrightarrow \mathbb{M}$ be a curve and define $\gamma = \exp^{-1} \circ \Gamma = \sum_{i=1}^\upsilon \gamma_i$, where $\gamma_i$ takes values in $W_i$. Then the following statements are equivalent:

1. $\Gamma$ is Lipschitz continuous,
(2) $\gamma$ is Lipschitz continuous and the differential equation

$$
\dot{\gamma}_i(t) = \sum_{n=2}^{v} \frac{(-1)^n}{n!} \pi_i ([\gamma(t), \dot{\gamma}(t)]_{n-1})
$$

is a.e. satisfied for every $i \geq 2$.

If one of the previous conditions holds, then there exists a constant $C > 0$ only depending on $\rho$ and $\| \cdot \|$, such that for any $\tau_1 < \tau_2$, we have

$$
\rho(\Gamma(\tau_1), \Gamma(\tau_2)) \leq C \int_{\tau_1}^{\tau_2} \| \dot{\gamma}_i(t) \| dt.
$$

**Remark 2.18.** If $\mathbb{M}$ is a stratified group, namely, its horizontal distribution satisfies the Lie bracket generating condition, then the previous theorem is well known. In fact, in this case equivalence of (1) and (2) would follow by Proposition 11.4 of [19] joined with Remark 2.16. This proposition shows that Lipschitz curves can be characterized as horizontal curves in the more general Carnot-Carathéodory spaces, that clearly include stratified groups. Moreover, estimate (17) can be obtained for instance by the sub-Riemannian area formula [22] joined with Theorem 2.10.13 of [13]. However, graded groups need not satisfy the Lie bracket generating condition, then they are not included in the family of Carnot-Carathéodory spaces.

### 3. Technical lemmata

In this section we study the properties of mappings $f : A \rightarrow \mathbb{M}$ where $\mathbb{M}$ is a graded group and $A$ may vary. $F$ denotes the mapping $\exp^{-1} \circ f : A \rightarrow \mathcal{M}$, where $\mathcal{M}$ is the Lie algebra of $\mathbb{M}$. The canonical projections onto the layers $W_i$ of the algebra $\mathcal{M}$ are the mappings $\pi : \mathcal{M} \rightarrow W_i$. We use the notation $F_i = \pi \circ F$, hence

$$
F = F_1 + \cdots + F_v \quad \text{and} \quad f = \exp (F_1 + \cdots + F_v).
$$

**Lemma 3.1.** Let $X \in V_1$ and let $N \subset G$ be a normal subgroup such that $G = N \times H_X$. Let $O \subset N$ and $J \subset H_X$ be open subsets, where $J$ is connected, and consider $u \in L^1_{\text{loc}}(\Omega)$, where $\Omega = OJ$ is an open set. If the distributional derivative $D_X u$ belongs to $L^\infty(\Omega)$, then up to redefinition of $u$ an a $\mu$-negligible set, for $\mu_N$-a.e. $n \in O$, we have

$$
|u(n \exp(tX)) - u(n \exp(\tau X))| \leq \|D_X u\|_{L^\infty(\Omega)} |t - \tau|.
$$

**Proof.** We adopt the same notation used in the proof of Proposition 2.8. Then we recall the mapping $\psi : \mathbb{R} \times \mathbb{R}^{q-1} \rightarrow \mathbb{M}$, observing that $\tilde{u} = u \circ \psi$ belongs to $L^1_{\text{loc}}(\psi^{-1}(\Omega))$. Taking into account that $F^{-1} \circ \psi$ has jacobian equal to one, by definition of distributional derivative one easily gets the following equality of distributional derivatives

$$
\nabla_1 \tilde{u} = (D_X u) \circ \psi.
$$

Then $\tilde{u} \in L^1_{\text{loc}}(\psi^{-1}(\Omega))$ and $\nabla_1 \tilde{u} \in L^\infty(\psi^{-1}(\Omega))$. The set $\psi^{-1}(O) = \tilde{O}$ is an open subset of $\{0\} \times \mathbb{R}^{q-1}$ and $\tilde{J} = \psi^{-1}(J)$ is an open interval of $\mathbb{R} \times \{0\} \subset \mathbb{R}^q$, hence $\psi^{-1}(O) = \tilde{J} \times \tilde{O}$. Thus, by a standard mollification argument, see for instance [12], Theorem 2 of Section 4.9.2, our claim follows. $\square$
Lemma 3.2. Let $X \in V_1$ and let $G = N \rtimes H_X$, where $N$ is a homogeneous normal subgroup. Let $O$ and $J$ be open subsets of $N$ and of $H_X$, respectively, where $J$ is connected, and let $z \in G$. We consider the open set $\Omega = zOJ$ along with the continuous mapping $f : \Omega \rightarrow \mathbb{M}$. Let $(X_1, \ldots, X_m)$ be a basis of $V_1$ and assume that there exist

$$\nabla X_i F_j \in L_{loc}^\infty(\Omega, W_j) \quad \text{and} \quad \nabla X_i F_1 \in L^\infty(\Omega, W_1)$$

for every $i = 1, \ldots, m$ and $j = 2, \ldots, v$ and the contact equations

$$X_i F_j = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j ([F, X_i F]_{n-1})$$

$\mu$-a.e. hold in $\Omega$, then there exists $C > 0$, depending on $\rho$ and $\| \cdot \|$, such that

$$\rho(f(\gamma_n), f(\gamma_m)) \leq C \|X F_1\|_{L^\infty(\Omega)} |t - \tau|$$

for every $\exp(tX), \exp(\tau X) \in J$ and every $n \in O$.

Proof. Up to a left translation, taking into account the left invariance of $X$, it is not restrictive to assume that $z$ is the unit element. Then we consider the curve

$$\Gamma_n(s) = n \exp(sX) = \exp(\gamma_{n,1}(s) + \cdots + \gamma_{n,v}(s))$$

where $\gamma_{n,j} = \pi_j \circ \gamma_n$ and $\gamma_n = \exp^{-1} \circ \Gamma_n$. We choose $t, \tau \in \mathbb{R}$ with $t < \tau$ such that $\exp(tX), \exp(\tau X) \in J$. By linearity, we have $D_X F$ exists and belongs to $L_{loc}^\infty(\Omega, \mathcal{M})$. Taking into account continuity of all $F_j$'s, Lemma 3.1 gives

$$\|F_j(\Gamma_n(t)) - F_j(\Gamma_n(\tau))\| \leq c \|X F_j\|_{L^\infty(\Omega_n)} (\tau - t)$$

for all $n \in O$, where $c > 0$ depends on $\| \cdot \|$ and $\Omega_n$ is an open neighbourhood of $\Gamma_n([t, \tau])$, that is compactly contained in $\Omega$. In particular, $t \rightarrow F \circ \Gamma_n(t)$ is absolutely continuous on compact intervals for all $n \in O$. By Proposition 2.8 from the $\mu$-a.e. validity of (19) it follows that for $\mu_N$-a.e. $n \in O$, we have

$$\frac{d}{ds} (F_j \circ \Gamma_n)(s) = X F_j \circ \Gamma_n(s) = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j ([F \circ \Gamma_n(s), (X F) \circ \Gamma_n(s)]_{n-1})$$

for a.e. $s \in \{l \in \mathbb{R} \mid \exp(lX) \in J\}$. Thus, we can apply Theorem 2.17 to the curve

$$f \circ \Gamma_n = \exp(F_1 \circ \Gamma_n + \cdots + F_v \circ \Gamma_n)$$

getting

$$\rho(f \circ \Gamma_n(t), f \circ \Gamma_n(\tau)) \leq C \int_{\tau}^t \|X F_1\circ \Gamma_n(s)\| \, ds.$$ 

Then the continuity of $f$ and the hypothesis $X F_1 \in L^\infty(\Omega)$ lead us to the conclusion.

Lemma 3.3. Let $f : X \rightarrow \mathbb{M}$ be a Lipschitz mapping, where $X$ is a metric space. Then there exists a constant $C > 0$, depending on the norm $\| \cdot \|$ of $\mathcal{M}$ and the distance $\rho$ of $\mathbb{M}$ such that

$$\|F_1(x) - F_1(y)\| \leq C \text{Lip}(f) \, d(x, y) \quad \text{for every} \quad x, y \in X.$$
Proof. We define
\[ |\xi| = \sum_{j=0}^{\nu} \|\xi_j\|^{1/j} \]
where \( \xi = \xi_1 + \xi_2 + \cdots + \xi_\nu \) and \( \xi_j \in W_j \). By continuity and homogeneity one can check that there exists a constant \( C > 0 \) such that
\[ \frac{1}{C} \rho(\exp \xi) \leq |\xi| \leq C \rho(\exp \xi) \]
Then we have
\[ \|\xi_1 - \eta_1\| \leq |\xi \ominus \eta| \quad \text{for every} \quad \xi, \eta \in \mathcal{M}, \]
where \( \xi_1 = \pi_1(\xi) \) and \( \eta_1 = \pi_1(\eta) \). It follows that
\[ \|F_1(x) - F_1(y)\| \leq C \rho(f(x), f(y)) \leq C \text{Lip}(f) d(x, y). \]
This concludes the proof. \( \square \)

Definition 3.4 (Piecewise horizontal line). A continuous curve \( \Gamma : [a, b] \to \mathbb{G} \) is a piecewise horizontal line if there exist \( n \in \mathbb{N} \), numbers \( a \leq t_0 < \cdots < t_n \leq b \) and \( X_k \in V_1 \) with \( k = 1, \ldots, n \) such that \( \Gamma|_{[t_k-1, t_k]}(t) = \Gamma(t_{k-1}) \exp \left((t - t_{k-1})X_k\right) \).

Lemma 3.5. Let \( \Omega \) be an open subset of \( \mathbb{G} \) and let \( \Gamma : [a, b] \to \Omega \) be a Lipschitz curve. Then there exists a sequence of piecewise horizontal lines uniformly converging to \( \Gamma \), whose lengths are also converging to the length of \( \Gamma \).

Proof. By Lemma 3.3 it follows that the curve \( \gamma_1 = \pi_1 \circ \exp^{-1} \circ \Gamma : [a, b] \to W_1 \) is Lipschitz. The derivative \( \dot{\gamma}_1 \) is essentially bounded. Let \( \varphi_k : [a, b] \to W_1 \) be a sequence of piecewise constant functions that a.e. converge to \( \dot{\gamma}_1 \) and \( |\varphi_k(t)| \leq |\dot{\gamma}_1(t)| \) for every \( k \) and a.e. \( t \in [a, b] \). We define the Lipschitz functions
\[ \gamma_{1k}(t) = \gamma_1(0) + \int_0^t \varphi_k(s) \, ds \]
that uniformly converge to \( \gamma_1 \). Exploiting (16), we can define
\[ \gamma_{2k}(t) = \gamma_2(0) + \frac{1}{2} \int_0^t [\gamma_{1k}(s), \dot{\gamma}_{1k}(s)] \, ds \]
and more generally by iteration
\[ \gamma_{ik}(t) = \gamma_{ik}(0) + \sum_{n=2}^{\nu} \frac{(-1)^n}{n!} \int_0^t \pi_i\left(\left[ \sum_{j=1}^{i-1} \gamma_{jk}(s), \sum_{j=1}^{i-1} \dot{\gamma}_{jk}(s) \right]_{n-1}\right) \, ds. \]
The curve \( \Gamma_k = \exp \left( \sum_{j=1}^{\nu} \gamma_{jk} \right) \) is the unique horizontal lifting of \( \gamma_k \). We wish to show that it is piecewise horizontal and uniformly converges to \( \Gamma \). To see this, we observe that the projection on the first layer of any curve
\[ t \to \exp \xi \exp tX \]
with \( X \in V_1 \) and \( \xi \in \mathbb{G} \) has the form \( t \to \xi_1 + tX \), where \( \pi_1(\xi) = \xi_1 \in V_1 \). We notice that \( \gamma_{1k} \) has exactly this form, since \( \varphi_k \) is piecewise constant, then the uniqueness of
the horizontal lifting implies that $\Gamma_k$ is piecewise horizontal. By construction, the fact that $\gamma_{1k}$ uniformly converges to $\gamma_1$ implies the uniform convergence of $\Gamma_k$ to $\Gamma$. From Corollary 5.5 of [25], we get the formula

$$\text{length}(\Gamma_k) = \text{Var}_a^b \rho = \int_a^b \rho \left( \exp \left( \dot{\gamma}_{1k}(t) \right) \right) dt.$$ 

It follows that $\text{length}(\Gamma_k) \to \text{length}(\Gamma)$ as $k \to \infty$. □

4. Group-valued mappings on Riemannian manifolds

In this section we extend contact equations to graded group-valued mappings on Riemannian manifolds and find the corresponding differential characterization of the Lipschitz property. In this case, contact equations can be written using the standard differential for differentiable manifolds.

Throughout this section, we have the following assumptions. We denote by $\mathbb{M}$ a graded group with graded algebra $\mathcal{M}$ and $N$ indicates a Riemannian manifold. The symbol $\rho$ denotes a homogeneous distance of $\mathbb{M}$. Any function $f : N \to \mathbb{M}$ is also written as the composition $f = \exp \circ F$, where $F : N \to \mathcal{M}$ and $\exp : \mathcal{M} \to \mathbb{M}$ denotes the exponential of Lie groups. We will also use the notation $F_i = \pi \circ F$, where $F = F_1 + \cdots + F_v$, $f = \exp (F_1 + \cdots + F_v)$ and $\pi_i : \mathcal{M} \to W_i$ is the canonical projection onto the $i$-th layer.

**Proposition 4.1.** Let $N$ be a one dimensional connected Riemannian manifold and let $F : N \to \mathcal{M}$ be a Lipschitz mapping such that

$$\dot{F}_j = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j \left( [F, \dot{F}]_{n-1} \right)$$

a.e. holds in $N$ for every $j = 2, \ldots, v$. Then the associated mapping $f : N \to \mathbb{M}$ is Lipschitz continuous and there exists a geometric constant $C > 0$ depending on the norm of $\mathcal{M}$ and the distance of $\mathbb{M}$ such that $\text{Lip}(f) \leq C \text{Lip}(F_1)$.

**Proof.** Let $x, y$ be two points of $N$ and let $c : [0, L] \to N$ be the length minimizing geodesic such that $c(0) = x$, $c(L) = y$ and $L = d(x, y)$. The curve $\gamma = F \circ c$ is Lipschitz and in view of (23) it satisfies (16) for a.e. $t$ in $[0, L]$, then by Theorem 2.17 we get

$$\rho(f(x), f(y)) = \rho(\Gamma(0), \Gamma(L)) \leq C \int_0^L \| (F_1 \circ c)'(t) \| dt$$

for a suitable geometric constant $C > 0$ depending on $\mathbb{M}$ and $\mathcal{M}$, where we have set $\Gamma = f \circ c$. Since

$$\left| (F_1 \circ c)(t + h) - (F_1 \circ c)(t) \right| \leq \text{Lip}(F_1) d(c(t + h), c(t)) = \text{Lip}(F_1) |h|,$$

where $d$ is the Riemannian distance of $N$. Then (24) leads us to our claim. □

**Remark 4.2.** As the linearity of $X \to [Z, X]_{n-1} = \text{ad}^{n-1}(X)$ makes (23) intrinsic, the same holds for higher dimensional manifolds. In fact, if $N$ has dimension $k \geq 1$ and
f is differentiable at a point $\bar{x}$ of $N$, then the contact condition $df(\bar{x})(T_{\bar{x}}N) \subset H_{f(\bar{x})}\mathbb{M}$ is equivalent to the validity of the system

$$\partial x_l F_j(\bar{x}) = \sum_{n=2}^{v} \frac{(-1)^n}{n!} \pi_j ([F(\bar{x}), \partial x_l F(\bar{x})]_{n-1}),$$

for every $j = 2, \ldots, v$ and $l = 1, \ldots, k$. This follows by Remark 2.16. Since the previous formula is independent of the local coordinates $(x_i)$ chosen around $\bar{x}$ in $N$, then an equivalent intrinsic version is the following one

$$dF_j(\bar{x}) = \sum_{n=2}^{v} \frac{(-1)^n}{n!} \pi_j ([F(\bar{x}), dF(\bar{x})]_{n-1}),$$

for every $j = 2, \ldots, v$. Here $d$ denotes the standard differential for mappings on differentiable manifolds and

$$dF_j(\bar{x}) : T_{\bar{x}}N \rightarrow H_{f(\bar{x})}^j \mathbb{M}.$$ 

Formula (25) can be also written in a shorter and intrinsic form adopting the Riemannian gradient as follows

$$\nabla F_j(\bar{x}) = \sum_{n=2}^{v} \frac{(-1)^n}{n!} \pi_j ([F(\bar{x}), \nabla F(\bar{x})]_{n-1}).$$

**Definition 4.3.** Let $N$ be a Riemannian manifold and let $d$ denote the Riemannian distance. We say that a subset $O \subset N$ is geodetically convex if for every $x, y \in O$, there exists a length minimizing geodesic $c : [0, L] \rightarrow O$ such that $c(0) = x, c(L) = y$.

**Remark 4.4.** Notice that notion of geodetic convexity we are adopting is not standard, since in Riemannian Geometry different forms of uniqueness of the connecting and minimizing geodesic are also required, see for instance [9]. For instance, according to our definition, the intersection of the 2-dimensional sphere $S^2$ embedded in $\mathbb{R}^3$ with any closed half space, where $S^2$ is equipped with the canonical Riemannian metric, is an example of geodetically convex set.

**Theorem 4.5.** Let $N$ be a connected Riemannian manifold of dimension higher than one and let $O \subset N$ be an open and geodetically convex set. Let $\nu$ be the canonical Riemannian measure on $N$ and let $F : O \rightarrow \mathcal{M}$ be a Lipschitz mapping such that

$$dF_j = \sum_{n=2}^{v} \frac{(-1)^n}{n!} \pi_j ([F, dF]_{n-1})$$

$\nu$-a.e. holds in $O$ for every $j = 2, \ldots, v$. Then the associated mapping $f : O \rightarrow \mathbb{M}$ is Lipschitz continuous and there exists a geometric constant $C > 0$ depending on the norm of $\mathcal{M}$ and the distance of $\mathbb{M}$ such that $\text{Lip}(f) \leq C \text{Lip}(F_1)$.

**Proof.** We choose two arbitrary points $p, q \in O$ and consider the smooth length minimizing geodesic $c : [0, L] \rightarrow N$, where $L = d(p, q)$, $c(0) = p$ and $c(L) = q$. Here $d$ denotes the Riemannian distance on $N$. Let $\tau \in [0, L]$ and let $(c(\tau), \dot{c}(\tau)) \in TN$. Let $TN$ be the open subset of $TN$ corresponding to the domain of the Riemannian
exponential mapping $\exp : TN \longrightarrow N$, see Theorem I.3.2 of [9]. We wish to construct a tubular neighbourhood of geodesics containing also the image of the restriction of $c$ to a neighbourhood of $\tau$.

To do this, we fix $z = c(\tau)$, $\xi = c(\tau)$ and select an arbitrary embedded smooth one codimensional submanifold $\Sigma \subset N$ passing through $z$, such that $T_z \Sigma$ is orthogonal to $\xi \in T_z N$. We will select all geodesics tangent to a normal field of $\Sigma$ in a neighbourhood of $z$. To make this argument rigorous, according to Chapter 5, p.132, of [29], we consider the vector bundle $T_{\Sigma}^\perp$ of fibers

$$T_s \Sigma^\perp = \{ v \in T_s N \mid v \in (T_s \Sigma)^\perp \subset T_s N \} \quad \text{for every } s \in \Sigma.$$

By definition, we have the orthogonal decomposition $T_s N = T_s \Sigma \oplus T_s \Sigma^\perp$. Now, we consider the normal exponential mapping $\exp^\perp$ as the following restriction

$$\exp^\perp : TN \cap T\Sigma^\perp \longrightarrow N.$$

Since the differential of $\exp^\perp$ is nonsingular at every point $(s, 0)$, in particular there exist open neighbourhoods $U$ of $(z, 0)$ in $TN \cap T\Sigma^\perp$ and $U$ of $z$ in $N$ such that $\exp^\perp : U \longrightarrow U$ is a smooth diffeomorphism. This provides us with a local system of coordinates around $z$ made by the local geodesic flow. Let $k$ be the dimension of $N$, hence up to shrinking both $U$ and $U$, we can select local coordinates $(y_1, \ldots, y_{k-1})$ of $\Sigma$ centered at zero, around $z$ and fix the local unit normal field $n$ of $\Sigma$ around $z$ such that $n(z) = \xi$. Then we define

$$H(y, t) = \exp^\perp \left( \zeta(y), t n(\zeta(y)) \right)$$

where $(A, \zeta)$ is a local chart of $\Sigma$, $A$ is an open subset of $\mathbb{R}^{k-1}$ containing the origin, with $\zeta(0) = z$, and $\left( \zeta(y), t n(\zeta(y)) \right) \in U$ if and only if $(y, t) \in A \times I$ for a suitable open interval $I$ of $\mathbb{R}$. By local uniqueness of geodesics, we get

$$I \ni t \longrightarrow c(\tau + t) = H(0, t) = \exp^\perp(z, t n(z)) = \exp(z, t \xi).$$

Since $H : A \times I \longrightarrow U$ is bilipschitz, then in view of Fubini’s theorem for a.e. $y \in A$ we have that $F$ is differentiable at $H(y, t)$ for a.e. $t \in I$ and there satisfies (28) for every $j = 2, \ldots, v$. Then in particular, we have the partial derivatives

$$\partial_j (F_j \circ H)(y, \cdot) = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j \left( \left[ (F \circ H)(y, \cdot), \partial_j (F_j \circ H)(y, \cdot) \right]_{n-1} \right)$$

a.e. in $I$ for a.e. $y \in A$. Since $F \circ H(y, \cdot)$ is also Lipschitz, then for a.e. $y \in A$ Proposition 4.1 yields a geometric constant $C_1 > 0$ such that

$$\rho(f \circ H(y, t), f \circ H(y, t')) \leq \text{Lip}(F_1 \circ H(y, \cdot)) C |t - t'|.$$

Continuity extends the previous estimate to all $y \in A$ and $t \in I$. In particular, for $y = 0$, it follows that

$$\rho((f \circ c)(\tau + t), (f \circ c)(\tau + t')) \leq C \text{Lip}(F_1) |t - t'|$$
for every \( t, t' \in I \). The arbitrary choice of \( \tau \in [0, L] \) gives a finite open covering of \([0, L]\) made of intervals satisfying (30). Since the constants of this estimate are independent of \( \tau \), this leads us to the end of the proof. \( \square \)

**Remark 4.6.** Under the hypotheses of the previous theorem, if we assume in addition that \( N \) is a complete Riemannian manifold, then the mapping \( f \) extends to a Lipschitz mapping \( f : \mathcal{O} \rightarrow \mathbb{M} \), with \( \text{Lip}(f) \leq C \text{Lip}(F_1) \). Since compact Riemannian manifolds are complete and then geodetically convex, then Theorem 4.5 obviously implies Theorem 1.2.

5. The differential characterization

In this section we give a proof of the differential characterization of locally Lipschitz mappings and show a simple application.

**Proof of Theorem 1.1.** Let \( f \) be locally Lipschitz. By Lemma 3.3, there exists \( C > 0 \) only depending on the norm \( \| \cdot \| \) and the distance \( \rho \), such that
\[
\| F_1(x) - F_1(y) \| \leq C \rho(f(x), f(y))
\]
for every \( x, y \in \Omega \). Then \( F_1 \) is also locally Lipschitz. We wish to show that all \( F_j \)'s are locally Lipschitz. By the Baker-Campbell-Hausdorff formula (8), we have
\[
\| - \xi_i + \eta_i + \sum_{n=2}^{\nu} \pi_i(c_n(-\xi, \eta)) \| = \| \pi_i(-\xi \odot \eta) \|^i \leq \left[ C \rho(\exp \xi, \exp \eta) \right]^i,
\]
then estimate (10) yields \( \tilde{\alpha}_n(\nu) > 0 \), with \( \nu = \max\{\|\xi\|, \|\eta\|\} \), such that
\[
\| - \xi_i + \eta_i \| \leq \sum_{n=2}^{\nu} \tilde{\alpha}_n(\nu) \|\xi, \eta\| \left[ C \rho(\exp \xi, \exp \eta) \right]^i,
\]
then taking into account (6) and (13), we get
\[
\| - \xi_i + \eta_i \| \leq \left( \sum_{n=2}^{\nu} \tilde{\alpha}_n(\nu) \beta \nu K_U + C^i \rho(\exp \xi, \exp \eta)^{i-1} \right) \rho(\exp \xi, \exp \eta)
\]
where \( U \) is a compact set containing \( \exp \xi \) and \( \exp \eta \), hence depending on \( \nu \). As a consequence, replacing both \( \xi_i \) and \( \eta_i \) with \( F_i(x) \) and \( F_i(y) \), we have shown that all \( F_i \)'s are locally Lipschitz. As a standard fact, it follows that all \( \nabla X_iF_j \) belong to \( L^\infty_{\text{loc}}(\Omega) \).

To prove the a.e. validity of (11), we choose an arbitrary left invariant vector field \( X \in V_1 \) and by Lemma 2.7 we get a normal subgroup \( N \) of \( G \) such that \( G = N \rtimes H_X \).

We can cover \( \Omega \) by a countable union of open sets of the form \( zOJ \), where \( z \in G \), \( O \) and \( J \) are connected open neighbourhoods of the unit element in \( N \) and in \( H_X \), respectively. Thus, we can reduce ourselves to prove our claim in the case \( \Omega = zOJ \).

Due to Theorem 2.11 in the case \( \mathbb{M} \) is a vector space, we consider the full measure set \( A \subset \Omega \) of points where all \( F_j \)'s are P-differentiable. Then formula (12) implies
\[
\nu_X(J \setminus p^{-1}z^{-1}A) = 0
\]
for $\mu_N$-a.e. $p \in O$. Let us pick one of these $p$'s. For a.e. $t$ we have $zp \exp(tX) \in A$, then all $F_j$'s are $P$-differentiable at this point. Definition of $P$-differentiability yields 
\[ DF_j(c_{zp,X}(t))(X) = XF_j(c_{zp,X}(t)), \]
for every $j = 2, \ldots, v$, where we have defined 
\[ \tau \mapsto c_{zp,X}(\tau) = zp \exp(\tau X). \]
As a result, setting 
\[ \tau \mapsto \Gamma_{zp,X}(\tau) = f \circ c_{zp,X}(\tau) = \exp \gamma_{zp,X}(\tau) = \exp \sum_{j=1}^v \gamma_{zp,X,j}(\tau), \]
where we have set $\gamma_{zp,X,j} = \pi_j \circ \gamma_{zp,X}$. We have then proved that 
\[ \gamma_{zp,X,j}(t) = XF_j \circ c_{zp,X}(t). \]
By Theorem 2.17 the local Lipschitz property of $f$ implies the a.e. validity of (16) for the curve $\gamma_{zp,X}$. Thus, taking into account (32), we obtain 
\[ XF_j(c_{zp,X}(\tau)) = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j ([F(c_{zp,X}(\tau)), XF(c_{zp,X}(\tau))]_{n-1}), \]
for a.e. $\tau$. We have proved the a.e. validity of the previous equation for $\mu_N$-a.e. $p \in O$, therefore Proposition 2.8 implies the validity of 
\[ (33) \quad XF_j = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j ([F, XF]_{n-1}) \]
$\mu$-a.e. in $\Omega$ for every $j = 2, \ldots, v$. The arbitrary choice of $X$ gives the validity of (16).

Conversely, we assume the validity of the second condition. By linearity of distributional derivative and of (16) for every $X \in V_1$, we have $D_X F \in L_{\infty}(\Omega, \mathcal{M})$ and (33) holds for every $j = 2, \ldots, v$. We fix $\|X\| = 1$ and choose $z \in \Omega$. We select connected open neighbourhoods $O \subset N$ and $J \subset H_X$ of the unit element $e \in G$ such that $\Omega = zOJ$ is compactly contained in $\Omega$. Thus, we can apply Lemma 3.2 getting 
\[ (34) \quad \rho(f(n \exp(tX)), f(n \exp(\tau X))) \leq C' \|XF_1\|_{L_{\infty}(\Omega)} |t - \tau| \]
\[ \leq C' \|\nabla_H F_1\|_{L_{\infty}(\Omega)} |t - \tau| \]
for every $n \in O$ and every $\exp(tX), \exp(\tau X) \in J$. Now, we consider the Carnot-Carathéodory distance $\delta$ generated by a left invariant metric fixed on $G$, see for instance [18]. This distance is homogeneous and then it is equivalent to $d$. It is well known that $(G, \delta)$ is a length space, namely any couple of points $x, y \in G$ is connected by a geodesic whose length equals their distance $\delta(x, y)$. Then from (34), we get 
\[ (35) \quad \rho(f(p), f(p')) \leq C \|\nabla_H F_1\|_{L_{\infty}(\Omega)} \delta(p, p') \quad \text{for every} \quad p, p' \in zJ, \]
where $C = C'/\inf_{|Y| = 1} \delta(\exp Y)$. Now we arbitrarily choose $r > 0$ and $p \in \Omega$ such that $B_{p, 6r} \subset \Omega$, where we have denoted by $B_{p, r}$ the open ball of center $x$ and radius $r$ with respect to $\delta$. Let $x, y \in B_{p, 2r}$ and let $\Gamma : [0, \delta(x, y)] \longrightarrow G$ be the geodesic connecting $x$ with $y$. By triangle inequality, it follows that the image of $\Gamma$ is contained
in $B_{p,6r}$. By Lemma 3.5 we can find a sequence $(\Gamma_k)$ of piecewise horizontal lines defined in $[0, \delta(x,y)]$ and contained in $B_{p,6r}$ that uniformly converge to $\Gamma$ and their lengths converge to $\delta$. On any horizontal segment of $\Gamma_k$ with horizontal direction $X_i$, that is also a geodesic, we apply the estimate (35) where $J$ is considered contained in the subgroup $H_{X_i}$. Thus, triangle inequality yields

$$\rho(f(\Gamma_k(0)), f(\Gamma_k(\delta(x,y)))) \leq C \|\nabla_H F_1\|_{L^\infty(\Omega)} l(\Gamma_k),$$

where $K_{p,r} = \overline{B_{p,6r}}$ and $l(\Gamma_k)$ is the length of $\Gamma_k$ with respect to $\delta$. Passing to the limit as $k \to \infty$, we have shown that

$$\rho(f(x), f(y)) \leq C \|\nabla_H F_1\|_{L^\infty(\Omega)} \delta(x,y)$$

for every $x, y \in B_{p,2r}$. Adopting the same argument of Theorem 3.18 of [23], it follows that $f$ is Lipschitz continuous on compact sets of $\Omega$, namely, $f$ is locally Lipschitz. Now, we show that the third condition is equivalent to the previous ones. We first assume that the first condition holds, namely, $f$ is locally Lipschitz. Due to Theorem 2.11, $f$ is a.e. P-differentiable and the equivalence of the first two conditions clearly yields $\nabla_{X_i} F \in L^\infty_{loc}(\Omega, \mathcal{M})$ for every $i = 1, \ldots, m$. If we know that $f$ is a.e. P-differentiable and $\nabla_{X_i} F \in L^\infty_{loc}(\Omega, \mathcal{M})$ for every $i = 1, \ldots, m$, then we apply Theorem 4.8 of [25], according to which the pointwise P-differentiability of $f$ implies the pointwise P-differentiability of all $F_j : \Omega \to \mathcal{M}$ with the validity of formulae

$$\pi_1 \circ df(x) = dF_1(x)$$

$$dF_i(x)(h) = \sum_{n=2}^\nu \frac{(-1)^n}{n!} \pi_i([F(x), dF(x)(h)]_{n-1}).$$

at P-differentiability points $x$. Taking into account that $dF_i$ denotes the P-differential read in the Lie algebras and that $dF_i(x)(\exp X) = X F_i$, these formulae implies the a.e. validity of contact equations (1). We have then shown that that the third conditions implies the second one. This concludes the proof. \(\square\)

5.1. Smooth functions that are not locally Lipschitz. We consider exponential coordinates $(x_1, x_2, x_3)$ of the Heisenberg group $\mathbb{H}^1$, with $p = \exp \left( \sum_{j=1}^3 x_j X_j \right) \in \mathbb{H}^1$ and the basis of left invariant vector fields given by

$$X_1 = \partial_{x_1} - x_2 \partial_{x_3}, \quad X_2 = \partial_{x_2} + x_1 \partial_{x_3}, \quad \text{and} \quad X_3 = \partial_{x_3}.$$

As mentioned in the introduction, the simplest case of smooth mapping that is not Lipschitz is the smooth curve $\Gamma : \mathbb{R} \to \mathbb{H}^1$ defined as $\Gamma(t) = (0,0,t)$. In this case, using any explicit homogeneous norm, it is straightforward to see the failure of the Lipschitz property.

On the other hand, if we consider the smooth mapping $I : \mathbb{H}^1 \to \mathbb{H}^1$ defined as

$$I(x_1, x_2, x_3) = (f_1, f_2, f_3) = (x_3, x_1, x_2),$$

then showing that any of its restrictions to open sets is not locally Lipschitz may require slightly more computations using distances and the use of contact equations.
could prove a convenient tool. We have
\[
\begin{align*}
&\begin{cases}
  f_1 X_1 f_2 - f_2 X_1 f_1 = x_3 + x_1 x_2 \\
  f_1 X_2 f_2 - f_2 X_2 f_1 = -x_1^2
\end{cases}
\quad\text{and}\quad
\begin{cases}
  X_1 f_3 = 0 \\
  X_2 f_3 = 1
\end{cases}
\end{align*}
\]
Thus, defining \( F_1 : \mathbb{H}^1 \to \mathfrak{h} \) as \( F_1(x_1, x_2, x_3) = x_3 X_1 + x_1 X_2 \in \mathfrak{h} \) and taking into account both (1) and Theorem 1.1, the mapping \( I \) is smooth, all components of the mapping \( I \) are continuously \( P \)-differentiable.

5.2. **Contact equations and Rumin complex.** In view of our study of Lipschitz mappings in the three dimensional Heisenberg group \( \mathbb{H}^1 \), we limit ourselves to recall the Rumin complex on \( \mathbb{H}^1 \), see [27] for the case of general contact manifolds. We denote by \( \Omega^k(\mathbb{H}^1) \) the module of \( k \)-forms on \( \mathbb{H}^1 \) and also
\[
J^2 = \{ \alpha \in \Omega^2(M) \mid \theta \wedge \alpha = 0 \}, \quad I^1 = \{ \varphi \theta \mid \varphi \in C^\infty(\mathbb{H}^1) \},
\]
where \( \theta = dt + (x_2 dx_1 - x_1 dx_2)/2 \) is the contact form. In this coordinates, we fix the left invariant vector fields
\[
X_1 = \partial_1 - \frac{x_2}{2} \partial_t, \quad X_2 = \partial_2 + \frac{x_1}{2} \partial_t \quad \text{and} \quad X_3 = \partial_t.
\]
We also set \( \Omega^1(\mathbb{H}^1)/I^1 = \{[\alpha_1 dx_1 + \alpha_2 dx_2 + \alpha_3 \theta]_R \mid \alpha_j \in C^\infty(\mathbb{H}^1) \} \). Clearly,
\[
[\alpha_1 dx_1 + \alpha_2 dx_2 + \alpha_3 \theta]_R = [\alpha_1 dx_1 + \alpha_2 dx_2]_R
\]
and we have the following

**Theorem 5.1** (Rumin, 1990). *There exists \( D : \Omega^1(\mathbb{H}^1)/I^1 \to J^2 \) such that*
\[
(37) \quad 0 \longrightarrow \mathbb{R} \longrightarrow C^\infty(\mathbb{H}^1) \longrightarrow \Omega^1(\mathbb{H}^1)/I^1 \overset{D}{\longrightarrow} J^2 \longrightarrow 0
\]
defines a complex whose cohomology coincides with the De Rham cohomology, where
\[
D[\alpha_1 dx_1 + \alpha_2 dx_2]_R = d(\alpha_1 dx_1 + \alpha_2 dx_2 + \alpha_3 \theta) \in J^2 \quad \text{defining} \quad \alpha_3 = X_1 \alpha_2 - X_2 \alpha_1.
\]

We denote by \( d_R \) the differential of this complex. For more information, we address the reader to [27]. In the next proposition, we will also use the following notation
\[
A = \begin{pmatrix}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{pmatrix}, \quad b = \begin{pmatrix}
b_1 \\
b_2
\end{pmatrix}, \quad A_j = \begin{pmatrix}
a_{1j} \\
a_{2j}
\end{pmatrix} \quad \text{and} \quad a = \begin{pmatrix}
a_1 \\
a_2
\end{pmatrix}.
\]

**Proposition 5.2.** *Let \( f = (F_1, f_3) : U \subset \mathbb{H}^1 \to \mathbb{H}^1 \) be a smooth mapping, where \( U \) is an open neighbourhood of the origin and \( F_1(x) = Ax + at + b \) and \( f_3 = f_3(x, t) \). Then \( f \) is locally Lipschitz if and only if for some \( \tau \in \mathbb{R} \) the following conditions hold*
\[
\begin{align*}
\text{(38)} & \quad \det (A, a) = \det (A, a) = 0, \\
\text{(39)} & \quad f_3(x, t) = \tau + \frac{x_2 \det (b, A_2) + x_1 \det (b, A_1)}{2} + t \left( \frac{\det (b, a)}{2} + \det A \right).
\end{align*}
\]
Proof. Since $f$ is Lipschitz and smooth, we have
\begin{equation}
\begin{cases}
X_1 f_3 = \alpha_1 = \frac{1}{2} (f_1 X_1 f_2 - f_2 X_1 f_1) \\
X_2 f_3 = \alpha_2 = \frac{1}{2} (f_1 X_2 f_2 - f_2 X_2 f_1)
\end{cases}
\end{equation}
everywhere in $U$. Since $d_R$ defines a complex, the Rumin’s complex, see [27], then
\[d_R ([\alpha_1 dx_1 + \alpha_2 dx_2]) = d_R (d_R f_3) = 0.\]
Taking into account that
\[d_R [\alpha_1 dx_1 + \alpha_2 dx_2] = d(\alpha_1 dx_1 + \alpha_2 dx_2 + \alpha_3 \theta) = (X_1 \alpha_3 - X_3 \alpha_1) dx_1 \wedge \theta + (X_2 \alpha_3 - X_3 \alpha_2) dx_2 \wedge \theta\]
where $\alpha_3 = X_1 \alpha_2 - X_2 \alpha_1$ and $\theta = dt + (x_2 dx_1 - x_1 dx_2)/2$. Then the system
\begin{equation}
\begin{cases}
X_1 \alpha_3 - \partial_t \alpha_1 = 0 \\
X_2 \alpha_3 - \partial_t \alpha_2 = 0
\end{cases}
\end{equation}
must hold. One can check that
\[2 \alpha_1 = f_1 X_1 f_2 - f_2 X_1 f_1\]
\[X_1 f_2 = a_{2} - \frac{a_2 x_2}{2} \quad X_1 f_1 = a_{1} - \frac{a_1 x_2}{2}\]
\[X_2 f_1 = a_{12} + \frac{a_1 x_1}{2} \quad X_2 f_2 = a_{2} + \frac{a_2 x_1}{2}.
\]
Thus, a direct calculation shows that
\[2 \alpha_1 = f_1 X_1 f_2 - f_2 X_1 f_1\]
\[= \left( \frac{\det (a b)}{2} - \det (A) \right) x_2 + \frac{\det (a A_1)}{2} x_1 x_2 + \frac{\det (a A_2)}{2} x_2^2\]
\[+ \det (a A_1) t + \det (b A_1)\]
and analogously
\[2 \alpha_2 = f_1 X_2 f_2 - f_2 X_2 f_1\]
\[= \left( \frac{\det (b a)}{2} + \det (A) \right) x_1 + \frac{\det (A_2 a)}{2} x_1 x_2 + \frac{\det (A_1 a)}{2} x_1^2\]
\[+ \det (a A_2) t + \det (b A_2)\].
Taking into account that $2 \alpha_3 = X_1 (2 \alpha_2) - X_2 (2 \alpha_1)$ and also
\[X_1 (2 \alpha_2) = - \frac{\det (a b)}{2} + \det (A) + \det (A_1 a) x_1 + \det (A_2 a) x_2,\]
\[X_2 (2 \alpha_1) = - \det (A) + \frac{\det (a b)}{2} + \det (A_2 a) x_2 + \det (A_1 a) x_1\]
we obtain
\[\alpha_3 = \det (A) + x_1 \det (A_1 a) + x_2 \det (A_2 a) - \frac{\det (a b)}{2}.\]
In view of (41), we get
\[ \det (A_1 a) = \det (A_2 a) = 0. \]
Notice that by Rumin complex, this is a necessary condition in order that the system (40) admits solutions. In order to solve the system (40), taking into account (42), direct computations yield
\[
\begin{cases}
2(f_3 \circ c_{p,X_1})'(s) = \left( \frac{\det (a b)}{2} - \det A \right) p_2 + \det (b A_1) \\
2(f_3 \circ c_{p,X_2})'(s) = \left( \frac{\det (b a)}{2} - \det A \right) p_1 + \det (b A_2)
\end{cases}
\]
where \( c_{p,X_j}(s) = p \exp(sX_j) \). Set \( p = (0, p_2, p_3) \), hence
\[
2f_3(p_1, p_2, p_3 - \frac{p_1 p_2}{2}) = 2f_3(0, p_2, p_3) + p_1 p_2 \left( \frac{\det (a b)}{2} - \det A \right) + p_1 \det (b A_1)
\]
and \( f(0, p_2, p_3) = f_3(0, 0, p_3) + p_2 \det (b A_2) \). It follows that
\[
f_3(x) = g \left( t + \frac{x_1 x_2}{2} \right) + \frac{x_2 \det (b A_2) + x_1 \det (b A_1)}{2} + \frac{x_1 x_2}{2} \left( \frac{\det (a b)}{2} - \det A \right),
\]
where \( g : \mathbb{R} \to \mathbb{R} \) is defined as \( g(t) = f_3(0, 0, t) \). By previous computations, we have
\[
\frac{f_1 X_2 f_2 - f_2 X_2 f_1}{2} = \frac{x_1}{2} \left( \frac{\det (b a)}{2} + \det A \right) + \frac{\det (b A_2)}{2}
\]
and our formula for \( f_3 \) gives
\[
X_2f_3 = x_1 g' \left( t + \frac{x_1 x_2}{2} \right) + \frac{\det (b A_2)}{2} - \frac{x_1}{2} \left( \frac{\det (b a)}{2} + \det A \right).
\]
Imposing the validity of the second equation of (40), it follows that
\[
g(t) = \tau + \left( \frac{\det (b a)}{2} + \det A \right) t.
\]
If \( f_3 \) satisfies formula (39), then \( f \) satisfies contact equations, then it is locally Lipschitz. This concludes the proof. \( \square \)

5.3. **Lipschitz mappings as boundary value problems.** This subsection is devoted to the characterization of existence of Lipschitz extensions as solutions of contact equations with assigned boundary datum.

**Theorem 5.3.** Let \( E \) be a closed set of a stratified group \( G \), let \( \Omega = G \setminus E \) and let \( f : E \to \mathbb{M} \) be a Lipschitz mapping. Existence of a Lipschitz extension \( \tilde{f} : G \to \mathbb{M} \) of \( f \) with possible larger Lipschitz constant is equivalent to the existence of a mapping \( g : \overline{\Omega} \to \mathbb{M} \), with \( G = \exp^{-1} \circ g \), \( G_j = \pi_j \circ G \) and \( g = \exp (G_1 + \cdots + G_n) \), such that

1. \( g_{\partial\Omega} = f_{\partial\Omega} \),
2. \( G_1 : \Omega \to W_1 \) is Lipschitz,
all the distributional derivatives $\nabla_H G \in L^\infty_{\text{loc}}(\Omega, \mathcal{M})^m$ and the system

$$\nabla_H G_j = \sum_{n=2}^v \frac{(-1)^n}{n!} \pi_j ([G, \nabla_H G]_{n-1})$$

$\mu$-a.e. holds in $\Omega$, for each $j = 2, \ldots v$.

Whenever a function $g$ satisfying these three conditions exists, then the Lipschitz extension of $f$ is given by $\tilde{f} = f 1_E + g 1_{\Omega}$ and we have the estimate

$$\text{Lip}(\tilde{f}) \leq C \left(\|\nabla_H G_1\|_{L^\infty(\Omega)} + \text{Lip}(f)\right),$$

for some geometric constant $C > 0$.

**Proof.** If $f$ admits a Lipschitz extension $\tilde{f}$, in view of Theorem 1.1 taking the restriction $\tilde{f}|_{\Omega} = g$, it follows that $g$ satisfies conditions (1), (2) and (3). Conversely, let us assume the existence of a mapping $g : \overline{\Omega} \rightarrow \mathbb{M}$ satisfying these three conditions. Let $p \in G$ and let $X \in V_1$ be arbitrarily fixed, with $\|X\| = 1$. Let us consider the curve

$$\mathbb{R} \ni t \rightarrow c_{p,X}(t) = p \exp(tX) \in G.$$

The open set $c_{p,X}^{-1}(\Omega)$ is the disjoint union $\cup_{j \in \mathbb{N}} I_j$ of open intervals of $\mathbb{R}$, where $I_j = [a_j, b_j]$. Let $N$ be a homogeneous normal subgroup such that $G = N \times H_X$ and let $a_j < a_j^+ < b_j < b_j^+$. The element $p$ is written in a unique way as $\pi \exp(\tilde{t}X)$, due to Proposition 2.6. Then we fix a relatively compact, connected open neighbourhood $O \subset N$ of the unit element $e$ of $G$ such that

$$O \overline{\exp \left( (\tilde{t} + s)X \right)} \subset \Omega \quad \text{for every} \quad s \in [a_j^+, b_j^+].$$

By assumptions (2) and (3) on $g$, Lemma 3.2 applies, hence we have in particular

$$\rho(g(c_{p,X}(b_j)), g(c_{p,X}(a_j))) \leq C' \|XG_1\|_{L^\infty(\Omega)} d(c_{p,X}(b_j), c_{p,X}(a_j)),$$

where we have set $C' = C / \inf_{\|Y\| = 1} d(\exp Y)$. Passing to the limit as $a_j^+ \to a_j$ and $b_j^+ \to b_j^-$, we achieve

$$\rho(g(c_{p,X}(b_j)), g(c_{p,X}(a_j))) \leq C' \|XG_1\|_{L^\infty(\Omega)} d(c_{p,X}(b_j), c_{p,X}(a_j)).$$

Now, consider $\tilde{f} = f 1_E + g 1_{\Omega}$ and arbitrary select $t, \tau \in \mathbb{R}$, with $t < \tau$. We have the following cases. If $t, \tau \in I_j$ for some $j$, then (44) yields

$$\rho(\tilde{f}(c_{p,X}(t)), \tilde{f}(c_{p,X}(\tau))) \leq C' \|XG_1\|_{L^\infty(\Omega)} d(c_{p,X}(t), c_{p,X}(\tau)).$$

If $t \in I_j$ and $\tau \notin c_{p,X}^{-1}(\Omega)$, then the triangle inequality and the fact that $c_{p,X}$ is a geodesic yield

$$\rho(\tilde{f}(c_{p,X}(t)), \tilde{f}(c_{p,X}(\tau)))$$

$$\leq C' \|XG_1\|_{L^\infty(\Omega)} d(c_{p,X}(t), c_{p,X}(b_j)) + \text{Lip}(f) d(c_{p,X}(b_j), c_{p,X}(\tau))$$

$$\leq (C' \|XG_1\|_{L^\infty(\Omega)} + \text{Lip}(f)) d(c_{p,X}(t), c_{p,X}(\tau)).$$
The same estimate is obtained in the analogous case $\tau \in I_j$ and $t \notin c_{p,X}^{-1}(\Omega)$. If $t \in I_j$ and $\tau \in I_k$ with $j \neq k$, with analogous argument we get
$$
\rho(\tilde{f}(c_{p,X}(t)), \tilde{f}(c_{p,X}(\tau))) \leq (2C' \|XG_1\|_{L^\infty(\Omega)} + \text{Lip}(f)) \, d(c_{p,X}(t), c_{p,X}(\tau)).
$$
The remaining case $t, \tau \notin c_{p,X}^{-1}(\Omega)$ is trivial, since $c_{p,X}(t), c_{p,X}(\tau) \in E$, where $f$ is Lipschitz. We have shown that for every $r, r' \in c_{p,X}(\mathbb{R})$, we have
$$
(47) \quad \rho(\tilde{f}(r), \tilde{f}(r')) \leq (2C' \|\nabla H G_1\|_{L^\infty(\Omega)} + \text{Lip}(f)) \, d(r, r').
$$
Finally, we adopt the same argument used in the proof of Theorem 1.1, where we connect two arbitrary points $p, p' \in G$ by a geodesic with respect to the length distance $\delta$, then we approximate the geodesic by a sequence of piecewise horizontal lines, according to Lemma 3.5 and we let the estimate (47) pass to the limit. \(\square\)

6. ALLCOCK GROUPS

In this section, we introduce the class of Allcock groups, along with examples. The characterizing geometric property of these groups is related to the notion of isotropic homotopy. Throughout this section, we fix a 2-step graded algebra $n$ with first layer $v$ and second layer $z$. We select a scalar product on $n$ such that $v$ and $z$ are orthonormal.

**Definition 6.1.** Let $a, b : [0, 1] \rightarrow v$ be Lipschitz loops. We say that $\Gamma : [0, 1]^2 \rightarrow v$ is an isotropic homotopy carrying $a$ to $b$ if $(\tau, t) \rightarrow \Gamma(\tau, t)$ is Lipschitz,
$$
[\partial_\tau \Gamma, \partial_t \Gamma] = 0 \quad \text{a.e. in } [0, 1]^2,
$$
$\Gamma(\cdot, 0) = a$, $\Gamma(\cdot, 1) = b$ and $\Gamma(0, \cdot) = \Gamma(1, \cdot)$.

**Remark 6.2.** For our purposes, the points $\Gamma(0, \cdot)$ and $\Gamma(1, \cdot)$ need not coincide with some fixed point.

In the sequel, we will use the following linear space
$$
\text{Av}_0^\infty = \left\{ \sigma \in L^\infty([0,1], \mathcal{J}) \mid \int_0^1 \sigma = 0 \right\}.
$$
In the next definition the same symbol $| \cdot |$ will denote both the norm of vectors and 2-vectors. The individual cases will be clear from the context.

**Definition 6.3.** We say that $n$ is surjective on isotropic loops if there exists a constant $C > 0$ such that for every $\lambda > 0$ and every $\sigma \in \text{Av}_0^\infty$ one can find a loop $a \in \text{Lip}([0, 1], v)$ and an isotropic homotopy $\Gamma : [0, 1]^2 \rightarrow v$ that carries $a$ to a fixed point, such that $[a, \dot{a}] = \sigma$ a.e., $|a(0)| \leq C \lambda$ and the estimates
$$
(48) \quad \text{Lip}(\Gamma) \leq C \left( \lambda + \frac{\|\sigma\|_{L^\infty}}{\lambda} \right) \quad \text{and} \quad |\dot{a}| \leq \frac{C'}{\lambda} |\sigma| \quad \text{a.e.}
$$
along with
$$
(49) \quad \int_0^1 \int_0^1 |\Gamma_\tau \wedge \Gamma_t| \, d\tau \, dt \leq C \left( \int_0^1 |\dot{a}(t)| \, dt \right)^2.
$$
Remark 6.4. The $L^\infty$-norm understood for $\sigma \in L^\infty([0,1], \mathbb{R})$ in the previous definition is given by
\[
\|\sigma\|_{L^\infty} = \max_{j=1,\ldots,s} \|\sigma_j\|_{L^\infty},
\]
where $\sigma = \sum_{j=1}^s \sigma_j Z_j$ and $(Z_1, \ldots, Z_s)$ is an orthonormal basis of $\mathfrak{z}$.

Definition 6.5 (Allcock group). Let $n$ be surjective on isotropic loops and define
\[
V_1 = \bigoplus_{j=1}^n v_j, \quad V_2 = \mathfrak{z}, \quad [v_i, v_j] = \{0\} \quad \text{whenever} \quad i \neq j,
\]
where all the two step algebras $v_j \oplus \mathfrak{z}$ are isomorphic to $n = v \oplus \mathfrak{z}$ for every $j = 1, \ldots, n$. Then the two step algebra $V_1 \oplus V_2$ is denoted by $A\mathfrak{l}_n^0$. This algebra defines a unique stratified group $A\mathfrak{l}_n^0$, that we call Allcock group of model $n$. If the model $n$ is understood, then we denote an Allcock group simply by $A\mathfrak{l}_n^0$.

The following example shows that the three dimensional Heisenberg algebra $h^1$ is surjective on isotropic loops.

Example 6.6. Let $h$ be the Heisenberg algebra, with layers $\text{span}\{X_1, X_2\} = v$ and $\text{span}\{Z\} = \mathfrak{z}$, where $[X_1, X_2] = Z$. Let $\sigma \in A\mathfrak{v}_0^\infty$ and $\lambda > 0$, where $\sigma = \sigma_1 Z$. Then we define the Lipschitz curve
\[
a = a_1 X_1 + a_2 X_2, \quad \text{where} \quad a_1 \equiv \lambda \quad \text{and} \quad a_2(t) = \frac{1}{\lambda} \int_0^t \sigma_1(s) \, ds.
\]
Clearly $[a, \dot{a}] = \sigma$ and the corresponding homotopy is
\[
\Gamma(\tau, t) = \lambda X_1 + (1 - t)a_2(\tau)X_2
\]
is clearly isotropic and carries $a$ to the point $\lambda X_1 \in h$. We have $|a(0)| = \lambda$ and simple calculations yield
\[
\text{Lip}(\Gamma) \leq \frac{2 \|\sigma\|_{L^\infty}}{\lambda} \quad \text{and} \quad |\dot{a}| = \frac{|\sigma|}{\lambda}.
\]
The validity of (49) is trivial, since $|\Gamma_t \wedge \Gamma_{\tau}| = 0$ at any differentiability point of $\Gamma$.

Remark 6.7. The previous example shows that Heisenberg groups $H^n_3$ are Allcock groups, since $H^n_3 = A\mathfrak{l}_n^0$, where $h$ is the 3-dimensional Heisenberg algebra. Let us mention that for Heisenberg groups Lipschitz extensions from the Euclidean plane could be also treated by a different method, according to Fassler Master’s thesis, 2007.

Remark 6.8. Arguing as in the previous example, one can find several examples of 2-step stratified algebras $n$ with one codimensional horizontal distribution that suitably yield Allcock groups.

On the other hand, it is easy to find Allcock groups where the horizontal distribution has codimension higher than one, as in the following
Example 6.9. Let us consider the 2-step algebra \( \mathfrak{f}_s = \mathfrak{v} \oplus \mathfrak{z} \), where \( \mathfrak{v} = \{X_1, \ldots, X_{s+1}\} \), \( \mathfrak{z} = \text{span}\{Z_1, \ldots, Z_s\} \) and the only nontrivial bracket relations are

\[
[X_1, X_j] = Z_{j-1}, \quad \text{and} \quad j = 2, \ldots, s + 1.
\]

Let \( \sigma \in \mathrm{Av}_0^\infty \) and \( \lambda > 0 \), where \( \sigma = \sum_{k=1}^s \sigma_k Z_k \). Then we define the Lipschitz curve

\[
a = \sum_{j=1}^{s+1} a_j X_j, \quad \text{where} \quad a_1 = \lambda \quad \text{and} \quad a_j(t) = \frac{1}{\lambda} \int_0^t \sigma_{j-1}(s) \, ds
\]

for every \( j = 2, \ldots, s + 1 \). Clearly \([a, \dot{a}] = \sigma\), \(|a(0)| = \lambda\) and \(\dot{a} = \lambda^{-1} \sigma\). Moreover, the mapping

\[
\Gamma(\tau, t) = \lambda X_1 + (1 - t) \sum_{j=2}^{s+1} a_j(\tau) X_j
\]

is an isotropic homotopy carrying \( a \) to \( \lambda X_1 \in \mathfrak{v} \) and a simple computation yields

\[
\text{Lip}(\Gamma) \leq \frac{\sqrt{2s}}{\lambda} \|\sigma\|_{L^\infty}.
\]

Finally, to prove (19) we observe that there exists \( \beta_0 > 0 \) such that

\[
\int_{[0,1]^2} |\Gamma_t \wedge \Gamma_\tau| \, d\tau dt \leq \beta_0 \int_{[0,1]^2} \left| \sum_{j=2}^{s+1} a_j X_j \right| \left| \sum_{j=2}^{s+1} \dot{a}_j X_j \right| \, d\tau dt
\]

\[
= \beta_0 \int_{[0,1]^2} \left| \int_0^t \sum_{j=2}^{s+1} \dot{a}_j X_j \right| \left| \int_0^t \sum_{j=2}^{s+1} \dot{a}_j X_j \right| \, d\tau dt \leq \beta_0 \left( \int_0^1 |\dot{a}| \, dt \right)^2.
\]

We have proved that \( \mathfrak{f}_s \) is surjective on isotropic loops.

Remark 6.10. In view of the previous example, we have obtained other Allcock groups, corresponding to \( \mathcal{A}^n_k \). Notice that they have horizontal distribution of codimension \( s \), for every integer \( s \geq 1 \). Clearly, a 2-step algebra \( \mathfrak{n} \) of \( s \)-dimensional second layer, having a subalgebra isomorphic to \( \mathfrak{f}_s \) is surjective on isotropic loops.

Example 6.11. We define the “multi-Heisenberg algebra” \( \mathcal{M}h^s \), as the 2-step stratified algebra, where first and second layers are spanned by the bases \( \{X_1, \ldots, X_{2s}\} \) and \( \{Z_1, \ldots, Z_s\} \), respectively, and the only nontrivial brackets are

\[
[X_j, X_{s+j}] = Z_j \quad \text{for every} \quad j = 1, \ldots, s.
\]

Let \( \sigma \in \mathrm{Av}_0^\infty \), with \( \sigma = \sum_{j=1}^s \sigma_j Z_j \) and let \( \lambda > 0 \). We define

\[
a = \lambda \sum_{j=1}^s X_j + \sum_{j=1}^s a_j X_{s+j}, \quad \text{where} \quad a_j(\tau) = \frac{1}{\lambda} \int_0^\tau \sigma_j(s) ds.
\]

Then the homotopy \( \Gamma(\tau, t) = \lambda \sum_{j=1}^s X_j + (1 - t) \sum_{j=1}^s a_j(\tau) X_{s+j} \) is clearly isotropic and carries \( a \) to the point \( \lambda \sum_{j=1}^s X_j \). Clearly, \( |a(0)| = \lambda \sqrt{s} \) and \( |\dot{a}| = \lambda^{-1} |\sigma| \). We also
have the estimate
\[
\text{Lip}(\Gamma) \leq \text{ess sup}_{[0,1]^2} \left( \sum_{j=1}^{s} \alpha_j^2 + \dot{\alpha}_j^2 \right)^{1/2} \leq \frac{\sqrt{2s}}{\lambda} \|\sigma\|_{L^\infty}
\]

Finally, arguing exactly as in Example 6.9, we get
\[
\int_{[0,1]^2} |\Gamma_t \wedge \Gamma_r| \, dt \leq \beta_0 \int_{[0,1]^2} \left| \sum_{j=1}^{s} a_j X_{s+j} \right| \left| \sum_{j=1}^{s} \dot{a}_j X_j \right| \, dt \leq \beta_0 \left( \int_{0}^{1} |\dot{a}| \, dt \right)^2.
\]

We have shown that \( \mathcal{M}^{\ast}_h \) is surjective on isotropic loops, then \( \mathcal{A}^{[1]}_{\ast} \) are Allcock groups for every \( n, s \in \mathbb{N} \setminus \{0\} \).

**Example 6.12.** The *complexified Heisenberg algebra* is surjective on isotropic loops. Recall that this algebra \( \mathbb{C} \mathfrak{h} = \mathfrak{v} \oplus \mathfrak{z} \) is an H-type algebra, with \( J_Z : \mathfrak{v} \rightarrow \mathfrak{v} \) and \( J_Z^2 = -|Z|^2I \) for every \( Z \in \mathfrak{z} \). We fix an orthonormal basis \( (Z_1, Z_2) \) of \( \mathfrak{z} \) and define the unit vectors \( R_0 = X_0, R_1 = JZ_1X_0, R_2 = JZ_2X_0 \) and \( R_3 = JZ_1JZ_2X_0 \), that form an orthonormal basis of \( \mathfrak{v} \). For more information on the complexified Heisenberg algebra, see [30]. Let us fix \( \lambda > 0 \) and choose a curve \( \sigma = \sigma_1Z_1 + \sigma_2Z_2 \in \mathbb{A}v^\infty_0 \). We define
\[
a = \sum_{j=0}^{s} a_j R_j \in \text{Lip}_0([0,1], \mathfrak{v})
\]
as follows
\[
a_0 \equiv \lambda, \quad a_3 \equiv 0, \quad a_1(t) = \frac{1}{\lambda} \int_{0}^{t} \sigma_1 \quad \text{and} \quad a_2(t) = \frac{1}{\lambda} \int_{0}^{t} \sigma_2.
\]

This loop satisfies \([a, \dot{a}] = \lambda \dot{a}_1 Z_1 + \lambda \dot{a}_2 Z_2 = \sigma, \; |a(0)| = \lambda \) and \(|\dot{a}| = \lambda^{-1} |\sigma|\). We also notice that the isotropic homotopy \( \Gamma(\tau, t) = \lambda R_0 + (1-t)(a_1(\tau)R_1 + a_2(\tau)R_2) \) carries \( a \) to \( \lambda R_0 \in \mathfrak{v} \) and satisfies \( \text{Lip}(\Gamma) \leq 2 \|\sigma\|_{L^\infty} \lambda^{-1} \). As in the previous cases, the validity of (19) is straightforward. This proves that the complexified Heisenberg group \( \mathbb{C} \mathfrak{h}_1 \), corresponding to \( \mathcal{A}^{[1]}_{\mathfrak{c}^0} \), is an Allcock group and more generally \( \mathcal{A}^{[n]}_{\mathfrak{c}^0} = \mathbb{C} \mathfrak{h}_n \) are also Allcock groups. Notice that these groups are all H-type groups.

**Example 6.13.** Let us consider the *quaternionic H-type group* \( \mathfrak{n}_3 \), whose center \( \mathfrak{z} \) is spanned by the orthonormal basis \( (Z_1, Z_2, Z_3) \) and \( (X_0, JZ_1X_0, JZ_2X_0, JZ_3X_0) \) is an orthonormal basis of \( \mathfrak{v} \). \( X_0 \) is a fixed unit vector of \( \mathfrak{v} \) and we have
\[
J_{Z_1}J_{Z_2} = J_{Z_3}, \quad J_{Z_1}J_{Z_3} = -J_{Z_2}, \quad \text{and} \quad J_{Z_2}J_{Z_3} = J_{Z_1}.
\]

We define the direct product algebra \( \mathfrak{n}_3^3 = \mathfrak{v} \oplus \mathfrak{z} \), where \( \mathfrak{v} = \mathfrak{v}_1 \oplus \mathfrak{v}_2 \oplus \mathfrak{v}_3 \) and (\( R_{l0}, R_{l1}, R_{l2}, R_{l3} \)) is the orthonormal basis of \( \mathfrak{v}_l \) for \( l = 1, 2, 3 \). We have defined
\[
R_{l0} = X_l, \quad R_{l1} = JZ_lX_l, \quad R_{l2} = JZ_2X_l, \quad R_{l3} = JZ_3X_l
\]

where \( X_l \) is a unit vector of \( \mathfrak{v}_l \). Furthermore, whenever \( l \neq s \) we set
\[
[R_{li}, R_{lj}] = 0 \quad \text{for every} \; i, j = 1, 2, 3, 4.
\]

Let \( \lambda > 0 \) and let \( \sigma = \sum_{j=1}^{3} \sigma_j Z_j \in \mathbb{A}v^\infty_0 \). We define the curve
\[
a = \lambda \left( \sum_{l=1}^{3} R_{l0} \right) + \frac{1}{\lambda} \sum_{l=1}^{3} \left( \int_{0}^{t} \sigma_l(s) \, ds \right) R_{li} \in \text{Lip}_0([0,1], \mathfrak{v})
\]
Then one can easily check that 
\[ [a, \dot{a}] = \sum_{l=1}^{3} \sigma_l Z_l, \quad |a(0)| = \sqrt{3}\lambda \quad \text{and} \quad |\dot{a}| = \lambda^{-1} |\sigma|. \]

Finally, the isotropic homotopy
\[ \Gamma(\tau, t) = \lambda \left( \sum_{l=1}^{3} R_{l0} \right) + (1-t) \frac{1}{\lambda} \sum_{l=1}^{3} \left( \int_{0}^{\tau} \sigma_l(s) \, ds \right) R_{lt} \]
carries \( a \) to \( \lambda \left( \sum_{l=1}^{3} R_{l0} \right) \in v \) and satisfies \( \operatorname{Lip}(\Gamma) \leq \sqrt{6} \|\sigma\|_{L^\infty} \lambda^{-1} \). Thus, we have proved estimates (48). Estimates (49) are obtained as in the previous examples. We have then proved that the quaternionic H-type group \( N^3_H \) is an Allcock group. Notice that higher dimensional Allcock groups \( A^m_{n+3} \) are H-type groups of dimension \( 12n + 3 \), where \( n \) is a positive integer.

**Remark 6.14.** Similar computations can be adapted to the octonionic H-type group. In general, the principle to find Allcock groups is to add as many copies of the horizontal subspace as possible. This yields the suitable “room” to construct isotropic homotopies.

On the other hand, it is not difficult to find 2-step groups that are not Allcock groups. It suffices to show that some 2-step stratified algebras are not surjective on isotropic loops, as we show in the next example.

**Example 6.15.** The free 2-step free Lie algebra \( g_{5,2} = V_1 \oplus V_2 \) on five generators is not surjective on isotropic loops. Let \( (X_1, \ldots, X_5) \) be a basis of generators of \( V_1 \) and let \( Z_{lp} = [X_l, X_p] \) be the vectors defining a basis of \( V_2 \) where \( 1 \leq l < p \leq 5 \). The curve
\[ \sigma(t) = \sum_{1 \leq l < p \leq 5} \sigma_{lp}(t) Z_{lp} \]
is defined by setting

\[
\sigma_{lp}(t) = \begin{cases} 
0 & \text{if } (l, p) = (1, 2) \\
t - 1/2 & \text{if } (l, p) = (1, 3) \\
t - 1/2 & \text{if } (l, p) = (2, 3) \\
t - 1/2 & \text{if } (l, p) = (1, 4) \\
t - 1/2 & \text{if } (l, p) = (2, 4) \\
t - 1/2 & \text{if } (l, p) = (1, 5) \\
0 & \text{if } (l, p) = (2, 5) \\
0 & \text{otherwise}
\end{cases}
\]

(50)

Clearly belongs to \( \operatorname{Av}^0 \). Now, suppose by contradiction that there exists a Lipschitz function \( a \in \operatorname{Lip}_0([0, 1], V_1) \) such that
\[ [a, \dot{a}] = \sigma \quad \text{a.e. in } [0, 1]. \]

Writing \( a(t) = \sum_{j=1}^{5} a_j(t) X_j \), the previous condition yields
\[ \det \begin{pmatrix} \dot{a}_l & a_p \\ \dot{a}_l & a_p \end{pmatrix} = \sigma_{lp} \quad \text{a.e. in } [0, 1]. \]

(51)
Then $\sigma_{12} \equiv 0$, $\sigma_{13} = \sigma_{23} \neq 0$ a.e. imply $(a_1, \dot{a}_1) = \lambda (a_2, \dot{a}_2)$ and $\lambda(t) \neq 0$ for a.e. $t \in [0, 1]$. Then constraints (51) give $\sigma_{14} = \lambda \sigma_{24}$ that yields $\lambda \equiv 1$. Therefore the condition $\sigma_{15} = \lambda \sigma_{25}$ corresponds to $\sigma_{15} = \sigma_{25}$, that conflicts with (50).

7. Lipschitz extensions from the plane to Allcock groups

This section is devoted to the proof of Theorem 1.3. Let $\mathbb{A}_l^n$ be an Allcock group with $n = v \oplus z$ and denote by $m$ and $s$ the dimensions of $v$ and $z$, respectively. An orthonormal basis $(X_1, \ldots, X_m, Z_1, \ldots, Z_s)$ of $n$ will be fixed, where $(X_1, \ldots, X_m)$ and $(Z_1, \ldots, Z_s)$ are bases of $v$ and $z$, respectively. We will also choose an orthonormal basis $(X_{ij})_{j=1,\ldots,m}$ of $v$, and an orthonormal basis $(Z_j)_{j=1,\ldots,s}$ of $z$. Since $v_j$'s are all isomorphic to $v$, then we can select $X_{ij}$ such that

$$[X_{il}, X_{ip}] = [X_l, X_p] \quad \text{for every} \quad i = 1, \ldots, n \quad \text{and} \quad 1 \leq l < p \leq m.$$ 

Then we fix graded coordinates in $\mathbb{A}_l^n$ with respect to this basis. We consider

$$(x, y) = \sum_{i=1,\ldots,n} x_{ij} e_{ij} + \sum_{l=1,\ldots,s} y_l E_l$$

where $(e_{ij}, E_l)$ is the canonical basis of $\mathbb{R}^{mn} \times \mathbb{R}^s$. Precisely, a point in $\mathbb{A}_l^n$ of coordinates $(x, y)$ is given by

$$\exp \left( \sum_{i=1,\ldots,n} x_{ij} X_{ij} + \sum_{l=1,\ldots,s} y_l Z_l \right) \in \mathbb{A}_l^n,$$

where $\exp : \mathbb{A}_l^n \rightarrow \mathbb{A}_l^n$ is the standard exponential mapping. Due to (52), for every $i = 1, \ldots, n$, we have

$$[X_{il}, X_{ip}] = [X_l, X_p] = \sum_{k=1}^s b_{lp}^k Z_k.$$

Then, we are in the position to introduce the multi-symplectic form

$$\omega = \sum_{k=1}^s \omega^k E_k, \quad \text{where} \quad \omega^k = \sum_{j=1}^n \sum_{1 \leq l < p \leq m} b_{lp}^k \, dx_{jl} \wedge dx_{jp}. \quad (53)$$

Notice that $\omega$ is an $\mathbb{R}^s$-valued 2-form defined on $\mathbb{R}^{mn}$. A primitive of this form is

$$\theta = \sum_{k=1}^s \theta^k E_k, \quad \text{where} \quad \theta^k = \sum_{j=1}^n \sum_{1 \leq l < p \leq m} \frac{b_{lp}^k}{2} \left( x_{jl} \, dx_{jp} - x_{jp} \, dx_{jl} \right). \quad (54)$$

Then for every Lipschitz curve $c : [a, b] \rightarrow \mathbb{R}^{mn}$, we define the multi-symplectic area swept by $c$ as follows

$$\int_c \theta = \sum_{k=1}^s \left( \int_a^b (c^* \theta^k)(t) \, dt \right) E_k \in \mathbb{R}^k.$$

This definition can be clearly extended to one-dimensional compact Lipschitz manifolds through local parametrizations.
Definition 7.1. Let $A^n_1$ be an Allcock group and let $\mathbb{R}^{mn} \times \mathbb{R}^s$ be the associated graded coordinates, where $\mathbb{R}^{mn}$ is the space of coordinates of the first layer. Then $\omega$ and $\theta$ defined above are the associated multi-simplectic form and its primitive, respectively.

Definition 7.2. If $c_i : [0,1] \rightarrow \mathbb{R}^{mn}$, $i = 0,1$, are Lipschitz loops, we say that $\Gamma : [0,1]^2 \rightarrow \mathbb{R}^{mn}$ is a multi-isotropic homotopy carrying $c_0$ to $c_1$ if $\Gamma$ is Lipschitz, the pull-back $\mathbb{R}^s$-valued 2-form $\Gamma^* \omega$ a.e. vanishes in $[0,1]^2$, $\Gamma(\cdot,0) = c_0$, $\Gamma(\cdot,1) = c_1$ and $\Gamma(0,\cdot) = \Gamma(1,\cdot)$.

All the preceeding notions will play a key role in the proof of the next theorem. Recall that the Euclidean norm will be understood on $\mathbb{R}^{mn}$ and $S^1$ will be thought of as the subset $\{(x,y) \mid x^2 + y^2 = 1\}$ equipped with the Euclidean distance of $\mathbb{R}^2$. The symbol $D$ denotes the closed unit disk of $\mathbb{R}^2$.

Theorem 7.3. Let $A^n_1$ be an Allcock group and let $\omega$ and $\theta$ be the associated forms on $\mathbb{R}^{mn}$ with respect to graded coordinates, where $n \geq 2$. Then there exists a geometric constant $\kappa > 0$ such that for every Lipschitz loop $c : S^1 \rightarrow \mathbb{R}^{mn}$ with $\int_c \theta = 0$ and $c(1,0) = 0$, then there exists a Lipschitz extension $\varphi : D \rightarrow \mathbb{R}^{mn}$ such that $\varphi^*(\omega) = 0$ a.e. in $D$ and there exists

\begin{equation}
(55) \quad \text{Lip}(\varphi) \leq \kappa \text{Lip}(c).
\end{equation}

**Proof.** We consider $\alpha(t) = c(e^{2\pi it})$ and set $\tilde{\alpha} : \mathbb{R} \rightarrow \mathbb{R}^{mn}$ such that $\tilde{\alpha}_{|[0,1]} = \alpha$ and $\tilde{\alpha}(\tau) = \alpha(0) = \alpha(1) = 0$ for every $\tau \in \mathbb{R} \setminus [0,1]$. Clearly we have

\[
\text{Lip}(\alpha) = \text{Lip}(\tilde{\alpha}) \leq 2\pi \text{Lip}(c).
\]

We first change the parametrization by the homotopy

\[
\Gamma_1 : [0,1]^2 \rightarrow \mathbb{R}^{mn}, \quad \Gamma_1(\tau,t) = \tilde{\alpha}((1+t)\tau - t).
\]

Notice that $\Gamma_1$ is Lipschitz and clearly the pull-back form $\Gamma_1^* \omega$ vanishes a.e. in $[0,1]^2$. We have shown that $\alpha$ is isotropically homotopic to $[0,1] \ni \tau \rightarrow \tilde{\alpha}(2\tau - 1)$. Now, we write $\tilde{\alpha}(\tau) = \beta_1(\tau) + \beta_2(\tau)$, where we have set

\[
\beta_1(\tau) = \sum_{j=1}^{m} \tilde{\alpha}_{ij}(\tau) e_{1j} \quad \text{and} \quad \beta_2(\tau) = \sum_{j=1}^{m} \tilde{\alpha}_{ij}(\tau) e_{ij}.
\]

Taking into account formula (53) for the multi-simplectic 2-form $\omega$, we notice that

\[
\omega(\beta_1(\tau) \wedge \beta_2(\tau)) = \sum_{k=1}^{s} \omega^k (\beta_1(\tau) \wedge \beta_2(\tau)) E_k = 0,
\]

since $dx_{1i} \wedge dx_{jp} (e_{1r} \wedge e_{iu}) = 0$ whenever $i \neq 1$. Therefore, defining the homotopy

\[
\Gamma_2 : [0,1]^2 \rightarrow \mathbb{R}^{mn}, \quad \Gamma_2(\tau,t) = \beta_1(2\tau - 1 + t) + \beta_2(2\tau - 1),
\]

it follows that for a.e. $(\tau,t) \in [0,1]^2$, we have

\[
\Gamma_2^* \omega = \omega(\partial_\tau \Gamma_2 \wedge \partial_t \Gamma_2) = 2\omega(\dot{\beta}_2(2\tau - 1) \wedge \dot{\beta}_1(2\tau - 1 + t)) = 0.
\]
Then $\tilde{\alpha}(2\tau - 1)$ is multi-isotropically homotopic to the product curve

$$(\beta_1\beta_2)(\tau) = \beta_1(2\tau) + \beta_2(2\tau - 1).$$

Now, we move $\beta_1$ to a multi-symplectically orthogonal space defining

$$\beta_1(\tau) = \sum_{j=1}^{m} \tilde{\alpha}_1(\tau)e_{2j}.$$ 

Thus, we get $\omega(\beta_1 \wedge \beta_1) = 0$. Notice that to perform this move from $\beta_1$ to $\beta_1$, we have used the assumption $n \geq 2$. We define the isotropic homotopy $\Gamma_3 : [0, 1]^2 \rightarrow \mathbb{R}^{mn}$,

$$\Gamma_3(\tau, t) = \cos(t\pi/2)\beta_1(2\tau) + \sin(t\pi/2)\beta_1(2\tau) + \beta_2(2\tau - 1).$$

Now, we observe that

$$\omega_k(\beta_1 \wedge \dot{\beta}_1) = \sum_{j=1}^{n} \sum_{1 \leq l < p \leq m} b_{lp}^k (dx_{jl} \wedge dx_{jp})(\beta_1 \wedge \dot{\beta}_1)$$

and similarly we have

$$\omega_k(\beta_1 \wedge \dot{\beta}_1) = \sum_{j=1}^{n} \sum_{1 \leq l < p \leq m} b_{lp}^k (dx_{2l} \wedge dx_{2p})(\beta_1 \wedge \dot{\beta}_1),$$

hence $\omega(\beta_1 \wedge \dot{\beta}_1) = \omega(\bar{\beta}_1 \wedge \dot{\beta}_1)$ a.e. in $[0, 1]$. Thus, a simple computation yields

$$\Gamma_3^*\omega = \omega(\partial_3 \Gamma_3 \wedge \partial_\tau \Gamma_3)$$

$$= -\pi \sin(t\pi/2) \cos(t\pi/2) \left(\omega(\beta_1(2\tau) \wedge \dot{\beta}_1(2\tau)) - \omega(\beta_1(2\tau) \wedge \dot{\beta}_1(2\tau))\right)$$

$$+ \cos(t\pi/2) \omega(\beta_1(2\tau) \wedge \dot{\beta}_1(2\tau))$$

$$= \cos(t\pi/2) \omega(\beta_1(2\tau) \wedge \dot{\beta}_1(2\tau))$$

for every $t \in [0, 1]$ and a.e. $\tau \in [0, 1]$. On the other hand, $\bar{\beta}_1(2\tau)$ vanishes when $2^{-1} \leq \tau \leq 1$ and $\beta_2(2\tau - 1) = 0$ for a.e. $\tau \in [0, 2^{-1}]$, this implies that $\Gamma_3^*\omega = 0$ a.e. in $[0, 1]^2$. We have proved that $\beta_1\beta_2$ is multi-isotropically homotopic to

$$(\beta_1\beta_2)(\tau) = \beta_1(2\tau) + \beta_2(2\tau - 1).$$

To construct the next homotopy, we use the fact that the algebra $\mathfrak{n}$ is surjective on isotropic loops. Let us consider the $L^\infty$ curve

$$\sigma = \sum_{k=1}^{s} \omega^k \left(\beta_1\beta_2 \wedge \frac{d}{ds}(\beta_1\beta_2)\right) Z_k.$$
We wish to show that $\sigma \in \text{Av}_0$, namely $\int_0^1 \sigma = 0$. To do so, we use the vanishing of the multi-simplectic area $\int_c \theta = 0$, that yields
\[
\int_{S^1} c^* \theta = \int_0^1 \alpha^* \theta = 0.
\]
Then we get
\[
\int_0^1 (\beta_1 \beta_2)^* \theta = \int_0^{1/2} (\beta_1)^* \theta + \int_0^{1/2} \beta_1^* \theta + \int_{1/2}^1 \beta_2^* \theta = \int_0^1 (\beta_1 \beta_2)^* \theta = \int_0^1 \alpha^* \theta = 0.
\]
and the equalities
\[
(\beta_1 \beta_2)^* \theta = \frac{1}{2} \omega \left( (\beta_1 \beta_2) \wedge \frac{d}{ds} (\beta_1 \beta_2) \right) = \frac{1}{2} \sigma
\]
prove our claim, namely, $\sigma \in \text{Av}_0$. Since $n$ is surjective on isotropic loops, we set
\[
(57) \quad \lambda = \text{length}(c) = \int_0^1 |\dot{\alpha}(t)| \, dt \leq 2\pi \text{Lip}(c) = 2\pi L.
\]
and apply Definition 6.3. Then we can find a loop $a \in \text{Lip}([0, 1], v)$ that is isotropically homotopic to a point and in particular satisfies
\[
(58) \quad [a, \dot{a}] = \sigma \quad \text{and} \quad \max_{[0, 1]} |a| \leq C \left( \lambda + \frac{\|\sigma\|_{L^\infty}}{\lambda} \right).
\]
We write $a = \sum_{i=1}^m a_i X_i$, then
\[
\sum_{1 \leq l < p \leq m} \det \begin{pmatrix} a_l & a_p \\ \dot{a}_l & \dot{a}_p \end{pmatrix} b_{lp}^k = \omega^k \left( (\beta_1 \beta_2) \wedge \frac{d}{ds} (\beta_1 \beta_2) \right)
\]
for every $k = 1, \ldots, s$. We introduce the Lipschitz loop $q = \sum_{j=1}^m a_j e_{1j}$ in $\mathbb{R}^m$ and observe that
\[
\omega(q(\tau) \wedge (\beta_1 \beta_2)(\tau)) = 0.
\]
Furthermore, by definition of $\omega^k$, one immediately gets
\[
(59) \quad \omega^k (q \wedge \dot{q}) = \sum_{1 \leq l < p \leq m} \det \begin{pmatrix} a_l & a_p \\ \dot{a}_l & \dot{a}_p \end{pmatrix} b_{lp}^k.
\]
It follows that
\[
\omega (q \wedge \dot{q}) = \omega \left( (\beta_1 \beta_2) \wedge \frac{d}{d\tau} (\beta_1 \beta_2) \right) \quad \text{a.e. in } [0, 1].
\]
Due to the previous condition, arguing as for $\Gamma_3$, one easily finds that the Lipschitz mapping $\Gamma_4 : [0, 1]^2 \longrightarrow \mathbb{R}^m$ defined as
\[
\Gamma_4(\tau, t) = \cos(t\pi/2) (\beta_1 \beta_2)(\tau) + \sin(t\pi/2) q(\tau)
\]
is a multi-isotropic homotopy between $\beta_1 \beta_2$ and $q$. According to Definition 6.3, there exists an isotropic homotopy $\Gamma : [0, 1]^2 \longrightarrow v$ such that $\Gamma(\cdot, 0) = a$, $\Gamma(\cdot, 1) \equiv \xi$, $\Gamma(0, \cdot) = \Gamma(1, \cdot)$ and $[\partial \tau \Gamma, \partial_t \Gamma] = 0$ a.e. in $[0, 1]^2$. Writing $\Gamma = \sum_{i=1}^n \sum_{j=1}^m \Gamma_{ij} X_{ij}$ and
defining $\Gamma_5 : [0, 1]^2 \to \mathbb{R}^{mn}$ as $\Gamma_5 = \sum_{i=1}^n \sum_{j=1}^m \Gamma_{ij} e_{ij}$, we have that $\Gamma_5^* \omega = 0$ and estimate (48) yields a constant $\bar{C}_1 > 0$ such that

$$\text{Lip}(\Gamma_5) \leq \bar{C}_1 \left( \lambda + \frac{\|\sigma\|_{L^\infty}}{\lambda} \right).$$

Defining $\xi = \sum_{i,j} \xi_{ij} X_{ij}$ and $q_0 = \sum_{i,j} \xi_{ij} e_{ij} \in \mathbb{R}^{mn}$, one immediately check that $\Gamma_5$ makes $q$ multi-isotropically homotopic the a point $q_0 \in \mathbb{R}^{mn}$. Thus, pasting all the previous multi-isotropic homotopies, we get the following mapping $H : [0, 1]^2 \to \mathbb{R}^{mn}$, defined as

$$H(\tau, t) = \begin{cases} 
\Gamma_1(\tau, 5t) & 0 \leq t < 1/5 \\
\Gamma_2(\tau, 5t - 1) & 1/5 \leq t < 2/5 \\
\Gamma_3(\tau, 5t - 2) & 2/5 \leq t < 3/5 \\
\Gamma_4(\tau, 5t - 3) & 3/5 \leq t < 4/5 \\
\Gamma_5(\tau, 5t - 4) & 4/5 \leq t < 1 
\end{cases}.$$

It is clearly both continuous and a.e. differentiable in $[0, 1]^2$. Moreover, we also have $H^* (\omega) = 0$ a.e.in $[0, 1]^2$. By convexity of $[0, 1]^2$ and triangle inequality, one easily notices the following estimate

$$\text{Lip}(H) \leq 5 \sum_{j=1}^5 \text{Lip}(\Gamma_j).$$

Direct computations show that

$$\begin{cases} 
\text{Lip}(\Gamma_1) \leq 6\pi L \\
\text{Lip}(\Gamma_2) \leq 6\pi L \\
\text{Lip}(\Gamma_3) \leq 6\pi L + \frac{\pi}{2} \max_{[0, 1]} |\alpha| 
\end{cases},$$

where we recall that we have defined $L = \text{Lip}(c)$. Taking into account the definition of Allcock group, we have to estimate $\|\sigma\|_{L^\infty}$. Thus, taking into account (56), a direct computation yields

$$|\omega^k (\gamma(\tau) \wedge \dot{\gamma}(\tau))| \leq 2nm \left( \max_{1 \leq l < p \leq m} |b_{lp}| \right) |\gamma(\tau)| |\dot{\gamma}(\tau)|$$

for every $k = 1, \ldots, s$, hence estimates

$$\max_{[0, 1]} |\overline{\beta_1 \beta_2}| \leq \max_{[0, 1]} |\alpha| \quad \text{and} \quad \left\| \frac{d}{d\tau}(\overline{\beta_1 \beta_2}) \right\|_{L^\infty([0, 1], \mathbb{R}^{mn})} \leq 2\|\dot{\alpha}\|_{L^\infty([0, 1], \mathbb{R}^{mn})} \leq 4\pi L,$$

along with $\max_{[0, 1]} |\alpha| \leq \lambda$, lead us to the estimate

$$\|\sigma\|_{L^\infty} \leq 8\pi \sqrt{s} nm \left( \max_{1 \leq l < p \leq m} |b_{lp}| \right) \lambda L = C_0 \lambda L.$$

As a consequence, by virtue of (57) and (60), we get a constant $C_1 > 0$ such that

$$\text{Lip}(\Gamma_5) \leq C_1 L.$$
In addition, by (57) and (58), observing that \( \max_{[0,1]} |a| = \max_{[0,1]} |q| \), we also obtain
\[
\max_{[0,1]} |q| \leq C (2\pi + C_0) L,
\]
where \( C_0 > 0 \) is a geometric constant depending on the group. It is also obvious that \( \text{Lip}(q) \leq \text{Lip}(\Gamma_5) \leq C_1 L \). Thus, joining all the previous estimates, we get a new constant \( C_2 > 0 \) such that
\[
\text{Lip}(\Gamma_4) \leq C_2 L.
\]
By last inequality along with (63), all previous estimates for \( \text{Lip}(\Gamma_j) \) and applying (61), we have found a geometric constant \( \kappa_1 > 0 \), only depending on the group, such that
\[
(64) \quad \text{Lip}(H) \leq \kappa_1 L.
\]
Then \( H \) is a multi-isotropic homotopy such that \( H(\cdot,0) = \alpha : [0,1] \to \mathbb{R}^{mn} \) and \( H(\cdot,1) \equiv q_0 \). The condition \( H(0,t) = H(1,t) \) for every \( t \in [0,1] \) implies that
\[
\varphi(\rho e^{2\pi i \theta}) = \begin{cases} H(\theta, 2(1 - \rho)) & \text{if } 1/2 \leq \rho \leq 1 \\ q_0 & \text{if } 0 \leq \rho \leq 1/2 \end{cases}
\]
is well defined on the closed unit disk \( D \subset \mathbb{R}^2 \). Furthermore, a direct computation shows that \( \text{Lip}(\varphi|_{D \setminus B_1/2}) \leq c_0 \text{Lip}(H) \), for a suitable geometric constant \( c_0 > 0 \), where \( B_{1/2} = \{ z \in \mathbb{R}^2 | |z| < 1/2 \} \). Since \( \varphi|_{B_1/2} \equiv q_0 \), then
\[
\text{Lip}(\varphi) \leq c_0 \text{Lip}(H) \leq c_0 \kappa_1 L
\]
Clearly \( \varphi_{S^1} \equiv c \) and \( \varphi^* \omega = 0 \) a.e. in \( D \). Our claim is achieved. \( \square \)

**Remark 7.4.** The previous theorem extends Theorem 2.3 of [1], where the standard symplectic space \( \mathbb{R}^{2n} \) is replaced with the \( \mathbb{R}^{mn} \) equipped with the multi-symplectic form \( \omega = \sum_{k=1}^s \omega^k E_k \).

Next, we will show how Theorem 7.3 leads us to a Lipschitz extension theorem. We will use some abstract tools in metric spaces, following the work by Lang and Schlichenmaier, [21].

**Definition 7.5.** We say that a metric space \( Y \) is *Lipschitz m-connected* for some \( m \in \mathbb{N} \) if there exists a constant \( c_m > 0 \) such that any Lipschitz map \( \Gamma : S^m \to Y \) has a Lipschitz extension \( \Phi : D^{m+1} \to Y \) with estimate \( \text{Lip}(\Phi) \leq c_m \text{Lip}(\Gamma) \).

**Definition 7.6.** Let \( (X,Y) \) be a couple of metric spaces. We say that \( (X,Y) \) has the *Lipschitz extension property* if there exists \( C > 0 \) such that for every subset \( Z \subset X \) and every Lipschitz map \( f : Z \to Y \), there exists a Lipschitz extension \( \overline{f} : X \to Y \) such that \( \text{Lip}(\overline{f}) \leq C \text{Lip}(f) \).

**Theorem 7.7** (Lang-Schlichenmaier, [21]). Let \( X \) and \( Y \) be two metric spaces and suppose that the Nagata dimension of \( X \) is less than or equal to \( n \) and that \( Y \) is complete. If \( Y \) is Lipschitz m-connected for \( m = 0, 1, \ldots, n-1 \), then the pair \( (X,Y) \) has the Lipschitz extension property.
Taking into account that Nagata dimension of \( \mathbb{R}^2 \) is clearly two, by Theorem 7.7 it follows that both Lipschitz 0-connectedness and Lipschitz 1-connectedness of \( A^n \) imply Corollary 1.4. The former property is a consequence of the fact that \( A^n \), as any stratified group, is connected by geodesics. The latter is proved in the following

**Theorem 7.8.** \( A^n \) is 1-connected for every \( n \geq 2 \).

**Proof.** Let \( \Gamma : S^1 \rightarrow A^n \) be a Lipschitz loop. Up to a left translation, that preserves the Lipschitz constant of \( \Gamma \), we can assume that \( \Gamma(1,0) = e \), where \( e \) is the unit element of \( A^n \). We introduce the 1-periodic mappings \( a : [0,1] \rightarrow V_1 \) and \( b : [0,1] \rightarrow V_2 \) such that

\[
\Gamma(e^{2\pi it}) = \exp \left( a(t) + b(t) \right).
\]

We consider our fixed basis \( (X_{ij}) \) of \( V_1 \), that satisfies (52), along with the orthonormal basis \( (Z_k) \) of \( V_2 \). We define

\[
a(t) = \sum_{1 \leq i \leq n \atop 1 \leq j \leq m} \alpha_{ij}(t) X_{ij} \quad \text{and} \quad b(t) = \sum_{k=1}^{s} \beta_k(t) Z_k
\]

where \( \alpha = (\alpha_{ij}) : [0,1] \rightarrow \mathbb{R}^{mn} \) and \( \beta = (\beta_k) : [0,1] \rightarrow \mathbb{R}^s \) are Lipschitz loops that satisfy \( \alpha(0) = \alpha(1) = 0 \) and \( \beta(0) = \beta(1) = 0 \). Since \( \Gamma \) is Lipschitz, Theorem 1.1 implies the a.e. validity of contact equations

\[
(65) \quad \sum_{k=1}^{s} \beta_k Z_k = \frac{1}{2} [a, \dot{a}] = \frac{1}{2} \sum_{j=1}^{n} \sum_{1 \leq l < p \leq m} \sum_{k=1}^{s} b_{lp}^k \left( \alpha_{jl} \dot{\alpha}_{jp} - \alpha_{jp} \dot{\alpha}_{jl} \right) Z_k = \sum_{k=1}^{s} \alpha^\ast(\theta^k) Z_k
\]

As a consequence, the 1-periodicity of \( \beta \) yields

\[
(66) \quad \int_{\alpha} \theta = \sum_{k=1}^{s} \left( \int_{0}^{1} (\alpha^\ast(\theta^k))(t) \, dt \right) E_k = 0.
\]

We define the curve \( c = (c_{ij}) : S^1 \rightarrow \mathbb{R}^{mn} \) defined by

\[
c_{ij}(e^{2\pi it}) = \alpha_{ij}(t),
\]

therefore \( c \) has vanishing multi-symplectic area, due to (66), and \( c(1,0) = 0 \). This allows us to apply Theorem 7.3 getting a Lipschitz extension \( \varphi : D \rightarrow \mathbb{R}^{mn} \) of \( c : S^1 \rightarrow \mathbb{R}^{mn} \) such that

\[
(67) \quad \varphi^\ast \omega = 0 \quad \text{a.e. in } D.
\]

To construct the extension of \( \Gamma \), we introduce the function

\[
A : D \rightarrow V_1, \quad A(x_1, x_2) = \sum_{1 \leq i \leq n \atop 1 \leq j \leq m} \varphi_{ij}(x_1, x_2) X_{ij}.
\]

We observe that \( a(t) = A(e^{2\pi it}) \) and define \( z_k(e^{2\pi it}) = \beta_k(t) \), where \( z_k : S^1 \rightarrow \mathbb{R}^s \) is Lipschitz continuous for every \( k = 1, \ldots, s \). To achieve our claim, we have to find
a Lipschitz extension $T : D \to V_2$ such that $T|_{S^1} = \sum_{k=1}^{s} z_k Z_k$ and the following contact equations a.e. hold

\begin{equation}
\partial x_1 T = \frac{1}{2} [A, \partial x_1 A] \quad \text{and} \quad \partial x_2 T = \frac{1}{2} [A, \partial x_2 A].
\end{equation}

In fact, if such a function $B$ exists, then applying Theorem 1.2 in the case $N$ is closed Euclidean disk $D$ of $\mathbb{R}^2$, hence the mapping

$$\Phi = \exp (A + T) : D \to \mathbb{R}^n$$

is Lipschitz continuous and

$$\text{Lip}(\Phi) \leq C \text{Lip}(A) = C \text{Lip}(\varphi),$$

for a suitable geometric constant. Furthermore, $\Phi|_{S^1} = \Gamma$, hence Theorem 7.3 yields another geometric constant $\kappa$ such that

\begin{equation}
\text{Lip}(\Phi) \leq \kappa C \text{Lip}(c).
\end{equation}

By Lemma 3.3 there exists $C_1 > 0$ such that

\begin{equation}
\text{Lip}(\Phi) \leq \kappa C C_1 \text{Lip}(\Gamma).
\end{equation}

Thus, to conclude the proof, we are left to show the existence of $T : D \to V_2$ satisfying the contact equations (68) and the boundary condition $T|_{S^1} = \sum_{k=1}^{s} z_k Z_k$.

First, for every $k = 1, \ldots, s$, we consider the 1-form $f_k = f_k^1 dx_1 + f_k^2 dx_2$ of components

$$f_k^i = \frac{1}{2} \sum_{j=1}^{n} \sum_{1 \leq l < p \leq m} b_{lp}^k (\varphi_{jl} \partial x_i \varphi_{jp} - \varphi_{jp} \partial x_i \varphi_{jl}) \in L^\infty(B_1),$$

where $B_1$ is the unit open ball of $\mathbb{R}^2$. Smoothing the function $\varphi$ with $\varphi^\varepsilon = \varphi * \zeta_\varepsilon$, we obtain the approximating forms $g_{k\varepsilon} = g_{k\varepsilon}^1 dx_1 + g_{k\varepsilon}^2 dx_2$ of components

$$g_{k\varepsilon}^i = \frac{1}{2} \sum_{j=1}^{n} \sum_{1 \leq l < p \leq m} b_{lp}^k (\varphi_{jl}^\varepsilon \partial x_i \varphi_{jp} - \varphi_{jp}^\varepsilon \partial x_i \varphi_{jl}^\varepsilon),$$

where $\zeta_\varepsilon$ is a standard mollifier and $g_{k\varepsilon}$ is a smooth 1-form on $B_{1-\varepsilon}$. It follows that

$$dg_{k\varepsilon} = (\partial x_1 g_{k\varepsilon}^2 - \partial x_2 g_{k\varepsilon}^1) dx_1 \wedge dx_2$$

$$= \sum_{j=1}^{n} \sum_{1 \leq l < p \leq m} b_{lp}^k (\partial x_1 \varphi_{jl}^\varepsilon \partial x_2 \varphi_{jp}^\varepsilon - \partial x_2 \varphi_{jl}^\varepsilon \partial x_1 \varphi_{jp}^\varepsilon) dx_1 \wedge dx_2$$

$$= (\varphi^\varepsilon)^* \omega^k.$$
As a consequence, setting $f_{\kappa\varepsilon}^i = f_k^i * \zeta_\varepsilon$ and $f_{\kappa\varepsilon} = f_{\kappa\varepsilon}^1 dx_1 + f_{\kappa\varepsilon}^2 dx_2$ one gets
$$df_{\kappa\varepsilon} = 0 \quad \text{on } B_{1-\varepsilon}.$$  
This gives the existence of a unique function $\tilde{\psi}_{\kappa\varepsilon} \in C^\infty(B_{1-\varepsilon})$ such that $\tilde{\psi}_{\kappa\varepsilon}(0,0) = 0$ and $d\tilde{\psi}_{\kappa\varepsilon} = f_{\kappa\varepsilon}$ in $B_{1-\varepsilon}$. The family of functions $\tilde{\psi}_{\kappa\varepsilon}$ is uniformly Lipschitz continuous and uniformly bounded, hence Ascoli-Arzelà’s theorem gives the existence of a Lipschitz function $\tilde{\psi}_k : D \rightarrow \mathbb{R}$ that is the uniform limit of $\tilde{\psi}_{\kappa\varepsilon}$ on compact sets of $B_1$ and clearly extends to the closure $D$. Furthermore, $d\tilde{\psi}_k = f_k$ as distributions, hence taking into account that $\tilde{\psi}_k$ is Lipschitz, it follows that $d\tilde{\psi}_k = f_k$ a.e. in $D$. Thus, defining $\psi_k : D \rightarrow \mathbb{R}$ as $\psi_k = \tilde{\psi}_k + z_k(1,0) - \tilde{\psi}_k(1,0)$ for every $k = 1, \ldots, s$, we get
$$\psi_k(1,0) = z_k(1,0) \quad \text{and} \quad d\psi_k = \varphi^* \theta_k \quad \text{a.e. in } D.$$  
As a result, the function $T : D \rightarrow V_2$ given by $T = \sum_{k=1}^s z_k Z_k$ satisfies the condition $T(1,0) = \sum_{k=1}^s z_k(1,0) Z_k$ along with the contact equations (68). The last step is to show the validity of the boundary condition $T|_{S^1} = \sum_{k=1}^s z_k Z_k$.

To do this, we recall that the mapping $\Phi : D \rightarrow A^n_{\alpha}$
$$\Phi = \exp (A + T)$$
is Lipschitz continuous. In particular, $\Phi|_{S^1} : S^1 \rightarrow A^n_{\alpha}$ and $\Gamma : S^1 \rightarrow A^n_{\alpha}$ are Lipschitz curves with $\Phi(1,0) = \Gamma(1,0) = e$. Then these curves must coincide by uniqueness of the horizontal lifting. In fact, we have
$$\begin{cases} 
\Phi(e^{2\pi i t}) = \exp \left( a(t) + \tilde{b}(t) \right) \\
\Gamma(e^{2\pi i t}) = \exp \left( a(t) + b(t) \right)
\end{cases},$$
where $b(0) = \tilde{b}(0) = 0$. By Theorem 1.1 since both $\Phi|_{S^1}$ and $\Gamma$ are Lipschitz, then both $b$ and $\tilde{b}$ satisfy the contact equations
$$\dot{b} = \frac{1}{2} [a, \dot{a}] = \tilde{b} \quad \text{a.e. in } [0,1]$$
and clearly $b = \tilde{b}$. This shows that $\Phi|_{S^1} = \Gamma$ and leads us to the conclusion. \(\blacksquare\)

Remark 7.9. From definition of 1-connectedness, one easily observes that Theorem 7.8 exactly coincides with Theorem 1.3.

8. Quadratic isoperimetric inequalities

In this section, we prove the validity of quadratic isoperimetric inequalities in Allcock groups. We will follow conventions and notation of Section 7.

Let $g$ be a left invariant Riemannian metric defined on $A^n_{\alpha}$ such that the fixed basis $(X_{ij}, Z_{k})_{i,j,k}$ of $A^n_{\alpha}$ is orthonormal. Denote by $\rho$ the associated Carnot-Carathéodory distance defined on $A^n_{\alpha}$ as
$$\rho_0(x,y) = \inf_{\Gamma : x \rightarrow y, \Gamma \subset H A^n_{\alpha}} \int_0^1 \sqrt{g(\Gamma(t))(\Gamma'(t),\Gamma'(t))} \, dt$$
where $\Gamma : [0, 1] \to \mathbb{A}^n$. Notice that if $\Gamma(t) = \exp (\gamma_1(t) + \gamma_2(t))$, $\gamma_i(t) \in V_i$ and also $\Gamma'(t) = \sum_i a_{ij}(t) X_{ij} \in L^1((0, 1), V_i)$, then

$$
\int_0^1 \sqrt{g(\Gamma(t))(\Gamma'(t), \Gamma'(t))} \, dt = \int_0^1 \|a(t)\| \, dt = \int_0^1 \|\gamma'_1(t)\| \, dt.
$$

The last equality follows from [13], taking into account that $\Gamma$ is horizontal. The symbol $| \cdot |$ above also denotes the Hilbert norm in $\mathbb{A}^n$ that makes $(X_{ij}, Z_k)$ orthonormal. Abusing notation, we will use the same symbol to define a norm on $\mathbb{A}^n$ as follows

$$
|x - y| := |\exp^{-1}(x) - \exp^{-1}(y)|,
$$

for every $x, y \in \mathbb{A}^n$. We denote by $\varrho$ the Riemannian distance associated to $g$ and notice that $\varrho \leq \rho_0$.

**Proposition 8.1.** Let $\mathbb{V}$ be a horizontal subgroup of $\mathbb{A}^n$. Then for every $x, y \in \mathbb{V}$, we have $\rho_0(x, y) = \varrho(x, y) = |x - y|$.

**Proof.** Let $x = \exp \xi$ and $y = \exp \eta$ and consider the curve $\tilde{\Gamma}(t) = \exp (\xi + t(\eta - \xi))$, where $\xi, \eta \in V$, where $V$ is the Lie algebra of $\mathbb{V} \subset \exp(V_i)$. Then $[\xi, \eta] = 0$ and $\tilde{\Gamma}$ is horizontal. Furthermore,

$$
\int_0^1 \sqrt{g(\tilde{\Gamma}(t))(\tilde{\Gamma}'(t), \tilde{\Gamma}'(t))} \, dt = \int_0^1 |\gamma'_1(t)| \, dt = |\xi - \eta| = |x - y|.
$$

If $\Gamma = \exp \circ (\gamma_1 + \gamma_2)$ is any absolutely continuous curve connecting $x$ and $y$, then $\gamma_1(0) = \xi$ and $\gamma_1(1) = \eta$ and clearly

$$
\int_0^1 \sqrt{g(\Gamma(t))(\Gamma'(t), \Gamma'(t))} \, dt = \int_0^1 |\gamma'_1(t)| \, dt \geq |\xi - \eta| = |x - y|.
$$

This concludes the proof. $\square$

**Proposition 8.2.** Let $\Gamma : S^1 \to \mathbb{A}^n$ be a Lipschitz mapping defined as

$$
\Gamma = \exp \circ \left( \sum_{ij} c_{ij} X_{ij} + \sum_k z_k Z_k \right).
$$

Then we have $\text{length}_{\varrho} \Gamma = \text{length}_{\rho_0} \Gamma = \text{length}_{|\cdot|}(c)$.

**Proof.** Using the sub-Riemannian area formula, one gets

$$
\int_a^b \varrho(\exp(\dot{\gamma_1}(t))) \, dt = \int_a^b \rho_0(\exp(\dot{\gamma_1}(t))) \, dt = \int_G N(\Gamma, y) \, d\mathcal{H}^1(y),
$$

then Theorem 2.10.13 of [13] gives

$$
\int_a^b \varrho(\exp(\dot{\gamma_1}(t))) \, dt = \int_a^b \rho_0(\exp(\dot{\gamma_1}(t))) \, dt = V_a^b \Gamma = \text{length}_{\rho_0}(\Gamma) = \text{length}_{\varrho}(\Gamma).
$$

Taking into account Proposition 8.1, it follows that $\rho_0(\exp(\dot{\gamma_1}(t))) = |\gamma_1(t)| = |\dot{c}(t)|$.

$\square$
Definition 8.3. Let $L : \mathbb{R}^2 \longrightarrow \mathbb{A}^n$ be an h-homomorphism. Then the jacobian of $L$ is given by

$$J(L) = \frac{\mathcal{H}^2_{\rho_0}(L(A))}{\mathcal{H}^2_{|L|}(A)}$$

where $A$ is any set of positive measure in the plane.

The previous definition of jacobian has been introduced in [22] for h-homomorphisms of stratified groups. Notice that it does not depend on the choice of the set $A$.

Remark 8.4. Let $L : \mathbb{R}^2 \longrightarrow \mathbb{A}^n$ be an h-homomorphism. Then $L(\mathbb{R}^2) = \mathbb{V}$ is a horizontal subgroup of $\mathbb{A}^n$ and Proposition 8.1 implies that

$$\mathcal{H}^2_{\rho_0}(L(A)) = \mathcal{H}^2_{|L|}(L(A)),$$

then the classical area formula yields

$$J(L) = |L_1 \wedge L_2|$$

where $|L_1 \wedge L_2|$ is the classical Euclidean jacobian of $L$ and $L_i = L(e_i)$.

Proposition 8.5. Let $f : D \longrightarrow \mathbb{A}^n$ be a Lipschitz mapping defined as

$$f = \exp \left( \sum_{ij} \varphi_{ij} X_{ij} + \sum_k \psi_k Z_k \right).$$

Let $N(f, y) = \mathcal{H}^0(f^{-1}(y))$ be the multiplicity function. Then we have

$$\int_{\mathbb{A}^n} N(f, y) \mathcal{H}^2_{\rho_0}(y) = \int_D J(Df(x)) \, dx = \int_D |\partial_{x_1} \varphi \wedge \partial_{x_2} \varphi| \, dx.$$  

Proof. Since the Pansu differential $Df(x) : \mathbb{R}^2 \longrightarrow \mathbb{A}^n$ is an h-homomorphism, it has the following matrix representation

$$Df(x) = \begin{pmatrix}
\nabla \varphi_{11} \\
\nabla \varphi_{12} \\
\nabla \varphi_{13} \\
\vdots \\
\nabla \varphi_{nm} \\
0 \\
\vdots \\
0
\end{pmatrix}.$$ 

In view of Remark 8.4, the sub-Riemannian area formula of [22] concludes the proof. □

Theorem 8.6. Let $\mathbb{A}^n$ be an Allcock group and let $\omega$ and $\theta$ be the associated forms on $\mathbb{R}^{mn}$ with respect to fixed graded coordinates, where $n \geq 2$. Then there exists a geometric constant $K > 0$ such that for every Lipschitz loop $c : S^1 \longrightarrow \mathbb{R}^{mn}$, with $\int_c \theta = 0$ and
c(1,0) = 0 one can find a Lipschitz extension \( \varphi : D \rightarrow \mathbb{R}^{mn} \) such that \( \varphi^* \omega = 0 \) a.e. in \( D \) and

\[
\mathcal{H}_1^2(\varphi(D)) \leq \int_D |\partial_{x_1} \varphi \land \partial_{x_2} \varphi| \, dx \leq K \text{length}_1(c)^2,
\]

which length_1(c) is the length of \( c \) with respect to the Euclidean norm \( |\cdot| \) in \( \mathbb{R}^{mn} \).

**Proof.** We will continue the argument used in the proof of Theorem 7.3, exploiting the same notation. We first recall that

\[
\text{length}_1(c) = \int_0^1 |\alpha'(t)| \, dt
\]

is the length of \( c : S^1 \rightarrow \mathbb{R}^{mn} \), where \( \alpha(t) = c(e^{i2\pi t}) \). From the proof of Theorem 7.3, we recall the definition of \( \varphi : D \rightarrow \mathbb{R}^{mn} \) as

\[
\varphi(\rho e^{i2\pi \theta}) = \begin{cases} H(\theta, 2(1 - \rho)) & \text{if } 1/2 \leq \rho \leq 1 \\ q_0 & \text{if } 0 \leq \rho \leq 1/2 \end{cases},
\]

where \( H : [0, 1]^2 \rightarrow \mathbb{R}^{mn} \) is given by

\[
H(\tau, t) = \Gamma_{k+1}(\tau, 5t - k) \quad \text{if } \frac{k}{5} \leq t < \frac{k+1}{5} \quad \text{and } k = 0, 1, 2, 3, 4.
\]

If we set \( \tilde{\varphi}(\rho, \theta) = \varphi(\rho e^{i2\pi \theta}) \), then we have

\[
\int_D |\varphi_{x_1} \land \varphi_{x_2}| \, dx = \int_{1/2}^1 \int_0^1 J \tilde{\varphi}(\rho, \theta) \, d\theta \, d\rho,
\]

where \( |v \land w| \) is the Hilbert norm on 2-vectors of \( \Lambda^2(\mathbb{R}^{mn}) \) with respect to the canonical basis. This is a consequence of the change of variable \( \phi(\rho, \theta) = \rho e^{i2\pi \theta} \) and the fact that \( |\phi_{\rho} \land \phi_{\theta}| = \rho \). It follows that

\[
\int_D J \varphi(x) \, dx = 2 \int_{1/2}^1 \int_0^1 |H_{\theta}(\theta, 2(1 - \rho)) \land H_{\rho}(\theta, 2(1 - \rho))| \, d\theta \, d\rho.
\]

Taking into account that

\[
H(\theta, 2(1 - \rho)) = \Gamma_{k+1}(\theta, 10(1 - \rho) - k) \quad \text{if } 1 - \frac{k+1}{10} < \rho \leq 1 - \frac{k}{10} \quad \text{and } k = 0, 1, 2, 3, 4,
\]

a simple change of variable yields

\[
\int_D |\varphi_{x_1} \land \varphi_{x_2}| \, dx = \sum_{k=0}^4 \int_0^1 \int_0^1 |(\Gamma_{k+1})_\tau(\theta, t) \land (\Gamma_{k+1})_t(\theta, t)| \, dt \, d\theta.
\]

Now, we use the explicit formulas of \( \Gamma_j \) given in the proof of Theorem 7.3. From definition of \( \Gamma_1 \) it is obvious that \( |(\Gamma_1)_\tau \land (\Gamma_1)_t| = 0 \) a.e. in the unit square \([0,1]^2\),
denoted by $Q$. Simple computations yield the following estimates
\[
\int_Q |(\Gamma_2)_\tau \wedge (\Gamma_2)_t| d\tau dt = 2 \int_Q |\dot{\beta}_1(2\tau - 1 + t) \wedge \dot{\beta}_2(2\tau - 1)| d\tau dt \\
\leq 2 \int_0^1 \left( \int_0^1 |\dot{\beta}_1(2\tau - 1 + t)| dt \right) |\dot{\beta}_2(2\tau - 1)| d\tau \\
\leq 2 \int_{\mathbb{R}} \left( \int_{\mathbb{R}} |\dot{\alpha}'(t)| dt \right) |\dot{\alpha}'(2\tau - 1)| d\tau = 2 \left( \int_0^1 |\alpha'(t)| dt \right)^2.
\]
The area contributed by $\Gamma_3$ is given by
\[
\int_Q |(\Gamma_3)_\tau \wedge (\Gamma_3)_t| d\tau dt = \pi \int_Q \left( \cos \left( \frac{\pi t}{2} \right) \dot{\beta}_1(2\tau) + \sin \left( \frac{\pi t}{2} \right) \dot{\beta}_1(2\tau) + \dot{\beta}_2(2\tau - 1) \right) \\
\wedge \left( - \sin \left( \frac{\pi t}{2} \right) \beta_1(2\tau) + \cos \left( \frac{\pi t}{2} \right) \beta_1(2\tau) \right) d\tau dt \\
\leq \pi \int_Q \left( \cos \left( \frac{\pi t}{2} \right) \dot{\beta}_1(2\tau) + \sin \left( \frac{\pi t}{2} \right) \dot{\beta}_1(2\tau) + \dot{\beta}_2(2\tau - 1) \right) \\
\wedge \left( - \sin \left( \frac{\pi t}{2} \right) \beta_1(2\tau) + \cos \left( \frac{\pi t}{2} \right) \beta_1(2\tau) \right) d\tau dt \\
\leq 2\pi \int_0^1 \left( |\beta_1(2\tau)| + |\dot{\beta}_2(2\tau - 1)| \right) |\beta_1(2\tau)| d\tau \\
\leq \pi \int_0^1 |\dot{\alpha}(t)| dt \int_0^2 \left( |\dot{\alpha}'(t)| + |\dot{\alpha}'(t - 1)| \right) dt \\
= 3\pi \left( \int_0^1 |\dot{\alpha}(t)| dt \right)^2.
\]
Concerning the isotropic homotopy $\Gamma_4$, we get
\[
\int_Q |(\Gamma_4)_\tau \wedge (\Gamma_4)_t| d\tau dt = \frac{\pi}{2} \int_Q \left( \cos \left( \frac{\pi t}{2} \right) \beta_1(2\tau)(\tau) + \sin \left( \frac{\pi t}{2} \right) q'(\tau) \right) \\
\wedge \left( - \sin \left( \frac{\pi t}{2} \right) \beta_1(\tau) + \cos \left( \frac{\pi t}{2} \right) q(\tau) \right) d\tau dt \\
\leq \frac{\pi}{2} \int_0^1 \left( |(\beta_1(\tau)| + |q'(\tau)| \right) \left( |\beta_1(\tau)| + |q(\tau)| \right) d\tau
\]
By definition of $\beta_1\beta_2$, one easily checks that
\[(74) \max_{[0,1]} |(\beta_1,\beta_2)| \leq 2 \int_0^1 |\dot{\alpha}(t)| dt \text{ and } |(\beta_1,\beta_2)'(t)| \leq 2 \left( |\dot{\alpha}'(2t)| + |\dot{\alpha}'(2t - 1)| \right).
\]
The curve $q(t)$ is given by the proof of Theorem [7.3] where it has been obtained applying Definition [6.3] with $\lambda = \text{length}_{\alpha_1}(c)$. It follows that there exists $C > 0$ such that $|q(0)| \leq C \lambda$. Thus, we get
\[
|q(t)| \leq C\lambda + \frac{C}{\lambda} \int_0^t |\sigma(t)| dt.
\]
Recall that
\[ \sigma(t) = \sum_{k=1}^{s} \omega^k \left( (\beta_1 \beta_2)(t) \wedge (\beta_1 \beta_2)'(t) \right) Z_k , \]
then bilinearity of \( \omega^k \)'s yields \( C_1 > 0 \) such that
\[ |\sigma(t)| \leq C_1 |(\beta_1 \beta_2)(t)| |(\beta_1 \beta_2)'(t)| , \]
hence in view of (74), it follows that
\[ (75) \quad |\sigma(t)| \leq 4C_1 \lambda \left( |\tilde{\alpha}'(2t)| + |\tilde{\alpha}'(2t - 1)| \right) . \]
As a result, we get
\[ \max_{t \in [0,1]} |q(t)| \leq C \lambda (1 + 4C_1) , \]
that implies
\[ \int_Q |(\Gamma_4)_{\tau} \wedge (\Gamma_4)_i| d\tau dt \leq \lambda \frac{\pi}{2} (2 + C + 4C_1 C) \int_0^1 \left( |(\beta_1 \beta_2)'(\tau)| + |q'(\tau)| \right) d\tau . \]
Taking into account that
\[ \int_0^1 |\dot{q}(t)| dt \leq \frac{C}{\lambda} \int_0^1 |\sigma(t)| dt \leq 2C C_1 \int_0^1 |(\beta_1 \beta_2)'(t)| \leq 4C C_1 \lambda , \]
we are then lead to the following
\[ \int_Q |(\Gamma_4)_{\tau} \wedge (\Gamma_4)_i| d\tau dt \leq \pi (2 + C + 4C_1 C) (1 + 2CC_1) \lambda^2 . \]
Finally, by definition of Allcock group and applying Definition 6.3, we have
\[ \int_Q |(\Gamma_5)_{\tau} \wedge (\Gamma_5)_i| d\tau dt \leq C \left( \int_0^1 |\dot{q}(t)| dt \right)^2 \leq 16 C^3 C_1^2 \lambda^2 . \]
Joining the previous estimates with (73), estimate (72) follows. \( \square \)

**Remark 8.7.** The previous theorem could be seen as a completion of Theorem 7.3, where we have fixed our attention on the area enclosed by the extension \( \varphi \).

**Proof of Theorem 1.5.** Up to left translation, we can assume that \( \Gamma(1,0) \) coincides with the unit element. Thus, we define
\[ \Gamma = \exp \circ \left( \sum_{i=1}^{n} \sum_{j=1}^{n} c_{ij} X_{ij} + \sum_{k=1}^{s} z_k Z_k \right) , \]
where \( c : S^1 \rightarrow \mathbb{R}^{mn} \). Arguing as in the beginning of the proof of Theorem 7.3, it follows that \( \int \theta = 0 \) and \( c(1,0) = 0 \). Then we apply Theorem 8.6 getting a Lipschitz extension \( \varphi : D \rightarrow \mathbb{R}^{mn} \) such that
\[ (76) \quad \int_D |\partial_{x_1} \varphi \wedge \partial_{x_2} \varphi| dx \leq K \text{length}_{\tilde{i},i}(c)^2 . \]
Finally, Proposition 8.2 and Proposition 8.5 lead us to the conclusion. \( \square \)
Example 8.8. Let us assume the setting in the proof of Theorem 1.5, where the Allcock group is the 7-dimensional Heisenberg group $H^2$. According to this proposition, we set $\phi = (\varphi, \psi)$. We equip $H^2$ with suitable graded coordinates $(x_1, \ldots, x_5)$ such that

$$\nabla \phi(x) = \begin{pmatrix} \nabla \varphi_1 \\ \nabla \varphi_2 \\ \nabla \varphi_3 \\ \nabla \varphi_4 \\ \varphi_1 \nabla \varphi_3 - \varphi_3 \nabla \varphi_1 + \varphi_2 \nabla \varphi_3 - \varphi_3 \nabla \varphi_2 \end{pmatrix},$$

(77)

where the contact equations imply that $\nabla \psi$ is equal to the last row of (77). Then a simple computation yields

$$|\partial_{x_1} \phi \wedge \partial_{x_2} \phi| \leq \sqrt{1 + 3 |\varphi|^2} |\partial_{x_1} \varphi \wedge \partial_{x_2} \varphi|.$$ 

Now, if $\varphi$ is the extension provided by Theorem 8.6 and $f = \exp\left(\sum_{j=1}^{4} \varphi_j X_j + \psi X_5\right)$ is the corresponding Lipschitz mapping, then

$$\mathcal{H}^2_{\varphi}(f(\mathcal{D})) \leq C \left(1 + 2 \max_{\mathcal{D}} |\varphi| \right) \text{length}_{\mathcal{H}}(c)^2$$

for a suitable geometric constant $C > 0$.

References

[1] D. Allcock, An isoperimetric inequality for the Heisenberg groups, GAFA, 8, 219-233, (1998)
[2] J. Alonso, Inégalités isopérimétriques et quasi-isométries, C.R.Acad.Sci.Paris, 311, 761-764, (1990)
[3] F.J. Almgren, Jr., The homotopy groups of the integral cycle groups, Topology, 1, (1962), 257-299.
[4] Z.M. Balogh, R. Hofer-Iシーズンger, J.T. Tyson, Lifts of Lipschitz maps and horizontal fractals in the Heisenberg group, Ergodic Theory Dynam. Systems, 26, 621-651, (2006)
[5] M.R. Bridson, The geometry of the word problem, Invitations to geometry and topology, Oxford University Press, p.29-91, (2002)
[6] J. Burillo, J. Taback, Equivalence of geometric and combinatorial Dehn functions. New York J. Math. 8, 169-179, (2002)
[7] L. Capogna, M. Cowling, Conformality and $Q$-Harmonicity in Carnot groups, Duke Math. J. 135, n.3, 455-479, (2006)
[8] D.B.A. Epstein, J.W. Cannon, D.F. Holt, S.V.F. Levy, M.S. Paterson, W.P. Thurston, Word Processing in Groups, Jones and Bartlett, Boston-London, (1992)
[9] L. Chavel, Riemannian Geometry. A Modern Introduction, Second edition, Cambridge University Press, (2006).
[10] J. Cheeger, B. Kleiner, Differentiating maps into $L^1$ and the geometry of BV functions, to appear on Annals of Mathematics
[11] N.S. Dairbekov, Mappings with bounded distortion of two-step Carnot groups, Proc. Anal. Geom., 122-155, Sobolev Institute Press, Novosibirsk, (2000)
[12] L.C. Evans, R.F. Gariepy, Measure theory and fine properties of functions, Studies in Advanced Mathematics, CRC Press, Boca Raton, (1992)
[13] H. Federer, Geometric Measure Theory, Springer, (1969).
[14] G.B. Folland, E.M. Stein, Hardy Spaces on Homogeneous groups, Princeton University Press, (1982)
[15] N. Garofalo, D. M. Nhieu, Lipschitz continuity, global smooth approximation and extension theorems for Sobolev functions in Carnot-Carathéodory spaces, Jour. Anal. Math., 74, 67-97 (1998)

[16] M. Gromov, Hyperbolic groups, in “Essays in Group Theory” (S. M. Gersten, ed.), MSRI Publications 8, 75-263, Springer-Verlag, (1987)

[17] M. Gromov, Asymptotic invariants for infinite groups, in “Geometric Group Theory”, vol. 2 (G. A. Niblo, M. A. Roller, eds.), London Mathematical Society Lecture Notes, 182, Cambridge University Press, (1993)

[18] M. Gromov, Carnot-Carathéodory spaces seen from within, Subriemannian Geometry, Progress in Mathematics, 144. ed. by A. Bellaiche and J. Risler, Birkhauser Verlag, Basel, 1996.

[19] P. Hajlasz, P. Koskela, Sobolev met Poincare, Memoirs of the American Mathematical Society, 688, (2000).

[20] A. Korányi, H. M. Reimann, Foundation for the Theory of Quasiconformal Mappings on the Heisenberg Group, Adv. Math., 111, 1-87, (1995).

[21] U. Lang, T. Schlichenmaier, Nagata dimension, quasisymmetric embeddings, and Lipschitz extensions, Int. Math. Res. Not. 58, 3625-3655, (2005).

[22] V. Magnani, Differentiability and Area formula on stratified Lie groups, Houston Jour. Math., 27, n.2, 297-323, (2001)

[23] V. Magnani, Lipschitz continuity, Aleksandrov theorem and characterizations for H-convex functions, Math. Ann., 334, n.1, 199-233, (2006)

[24] V. Magnani, Elements of Geometric Measure Theory on Sub-Riemannian groups, PhD theses series of Scuola Normale Superiore, (2002)

[25] V. Magnani, Towards Differential Calculus in stratified groups, arXiv:math/0701322v2, (2007)

[26] A. Y. Olshanskii, M. V. Sapir, Quadratic isometric functions of the Heisenberg group, A combinatorial proof. J. Math. Sci., 93, n.6, 921-927, (1999)

[27] M. Rumin, Un complexe de formes différentielles sur les variétés de contact, C. R. Acad. Sci. Paris Sér. I Math. 310, n.6, 401-404, (1990)

[28] P. Pansu, Métriques de Carnot-Carathéodory quasiisométries des espaces symétriques de rang un, Ann. Math., 129, 1-60, (1989)

[29] P. Petersen, Riemannian Geometry, Second edition, Springer, New York, (2006)

[30] H. M. Reimann, H. M., F. Ricci The complexified Heisenberg group, Proceedings on Analysis and Geometry (Russian) (Novosibirsk Akademgorodok, 1999), 465–480, Izdat. Ross. Akad. Nauk Sib. Otd. Inst. Mat., Novosibirsk, (2000).

[31] R. Young, Scaled relators and Dehn functions for nilpotent groups, arXiv:math/0601297v3, (2006).

[32] R. Young, Filling inequalities for nilpotent groups, arXiv:math/0608174v4 (2008)

[33] V. S. Varadarajan, Lie groups, Lie algebras and their representation, Springer-Verlag, New York, (1984).

Valentino Magnani, Dipartimento di Matematica, Largo Pontecorvo 5, I-56127, Pisa
E-mail address: magnani@dm.unipi.it