# Rank-one theorem and subgraphs of BV functions in Carnot groups ${ }^{\text {h }}$ 

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#### Abstract

We prove a rank-one theorem à $l a \mathrm{G}$. Alberti for the derivatives of vector-valued maps with bounded variation in a class of Carnot groups that includes Heisenberg groups $\mathbb{H}^{n}$ for $n \geq 2$. The main tools are properties relating the horizontal derivatives of a real-valued function with bounded variation and its subgraph.


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## 1. Introduction

One of the main results in the theory of functions with bounded variation (BV) is the rank-one theorem. Recall that a function $u \in L^{1}\left(\Omega, \mathbb{R}^{d}\right)$ has bounded variation in an open set $\Omega \subset \mathbb{R}^{n}\left(u \in B V\left(\Omega, \mathbb{R}^{d}\right)\right)$ if the derivatives $D u$ of $u$ in the sense of distributions are represented by a (matrix-valued) measure with finite total variation. The measure $D u$ can then be decomposed as the sum $D u=D^{a} u+D^{s} u$ of a measure $D^{a} u$, that is absolutely continuous with respect to $\mathscr{L}^{n}$, and a measure $D^{s} u$ that is singular with respect to $\mathscr{L}^{n}$. The Radon-Nikodym derivative $\frac{D^{s} u}{\left|D^{s} u\right|}$ of $D^{s} u$ with respect to its total variation $\left|D^{s} u\right|$ is a $\left|D^{s} u\right|$-measurable map from $\Omega$ to $\mathbb{R}^{d \times n}$. The rank-one theorem states that $\left|D^{s} u\right|$-a.e. this map takes values in the space of rank-one matrices. We refer to [3] for more details on $B V$ functions.

The rank-one theorem was first conjectured by L. Ambrosio and E. De Giorgi in [7] and it has important applications to vectorial variational problems and systems of PDEs. It was proved by G. Alberti in [1] (see also [2,8]): due to its complexity, Alberti's proof is generally regarded as a tour de force in measure theory. Two different proofs of the rank-one theorem were recently found. One is due to G. De Philippis and F. Rindler and follows from a profound PDE result [9], where a rank-one property for maps with bounded deformation ( BD ) was also proved for the first time. At the same time another proof, of a geometric flavor and considerably simpler than those in $[1,9]$, was provided by the second- and third-named authors in [27].

Motivated by these results, in this paper we consider the following natural generalization. Let $X_{1}, \ldots, X_{m}$ be linearly independent vector fields in $\mathbb{R}^{n}, m \leq n$, and let $u$ : $\Omega \rightarrow \mathbb{R}^{d}$ be a function with bounded $H$-variation in an open set $\Omega \subset \mathbb{R}^{n}$, i.e., a vector valued function such that the distributional horizontal derivatives $D_{H} u:=\left(X_{1} u, \ldots, X_{m} u\right)$ are represented by a $d \times m$-matrix valued measure with finite total variation in $\Omega$; consider the singular part $D_{H}^{s} u$ of $D_{H} u$ with respect to $\mathscr{L}^{n}$. Is it true that the Radon-Nikodym derivative $\frac{D_{H}^{s} u}{\left|D_{H}^{s} u\right|}$ is a rank-one matrix $\left|D_{H}^{s} u\right|$-a.e.?

We investigate this question in the setting of Carnot groups $\mathbb{G} \equiv \mathbb{R}^{n}$ (see Section 2) endowed with a left-invariant basis $X_{1}, \ldots, X_{m}$ of the first layer $\mathfrak{g}_{1}$ in the stratification of their Lie algebra. In particular, we find two assumptions on $\mathbb{G}$, that we call properties $\mathscr{C}_{2}$ and $\mathscr{R}$ (see Definitions 2.2 and 5.1, respectively), that ensure the rank-one property for $B V_{H}$ functions in $\mathbb{G}$. We will discuss later the role played by these properties in our argument. Our first main result is the following

Theorem 1.1. Let $\mathbb{G}$ be a Carnot group satisfying properties $\mathscr{C}_{2}$ and $\mathscr{R}$; let $\Omega \subset \mathbb{G}$ be an open set and $u \in B V_{H, l o c}\left(\Omega, \mathbb{R}^{d}\right)$ be a function with locally bounded $H$-variation. Then the singular part $D_{H}^{s} u$ of $D_{H} u$ is a rank-one measure, i.e., the matrix-valued function $\frac{D_{H}^{s} u}{\left|D_{H}^{s} u\right|}(x)$ has rank one for $\left|D_{H}^{s} u\right|$-a.e. $x \in \Omega$.

It is worth pointing out that Theorem 1.1 applies to the $n$-th Heisenberg group $\mathbb{H}^{n}$ provided $n \geq 2$. Recall that Heisenberg groups, defined in Example 2.1 below, are the most notable examples of Carnot groups.

Corollary 1.2. Let $u$ be as in Theorem 1.1 and assume that $\mathbb{G}$ is the Heisenberg group $\mathbb{H}^{n}$, $n \geq 2$; then $D_{H}^{s} u$ is a rank-one measure. More generally, the same holds if $\mathbb{G}$ is a Carnot group of step 2 satisfying property $\mathscr{C}_{2}$.

Corollary 1.2 is an immediate consequence of Theorem 1.1, see Remarks 2.4 and 5.3.
Theorem 1.1 does not directly follow from the outcomes of [9], see Remark 5.5. Its proof follows the geometric strategy devised in [27] and it is based on the relations between a (real-valued) $B V_{H}$ function $u$ in $\mathbb{G}$ and the $H$-perimeter of its subgraph $E_{u}:=\{(x, t): t<u(x)\} \subset \mathbb{G} \times \mathbb{R}$. Recall that a set $E \subset \mathbb{G} \times \mathbb{R}$ has finite $H$-perimeter if its characteristic function $\chi_{E}$ has bounded $H$-variation with respect the vector fields of a basis of the first layer in the Lie algebra stratification of the Carnot group $\mathbb{G} \times \mathbb{R}$. Our second main result is the following

Theorem 1.3. Suppose that $\Omega \subset \mathbb{G}$ is open and bounded and let $u \in L^{1}(\Omega)$. Then $u$ belongs to $B V_{H}(\Omega)$ if and only if its subgraph $E_{u}$ has finite $H$-perimeter in $\Omega \times \mathbb{R}$.

Actually, the proof of Theorem 1.1 requires much finer properties than the one stated in Theorem 1.3. Such properties are stated in Theorems 4.2 and 4.3 in a much more general context than Carnot groups, i.e., for maps with bounded $H$-variation with respect to a generic fixed family of linearly independent vector fields $X_{1}, \ldots, X_{m}$ on $\mathbb{R}^{n}$. Theorem 4.2, from which Theorem 1.3 immediately follows, focuses on the relations between the horizontal (in $\mathbb{R}^{n}$ ) derivatives of $u$ and the horizontal (in $\mathbb{R}^{n} \times \mathbb{R}$ ) derivatives of $\chi_{E_{u}}$. Theorem 4.3 instead deals with the relations between the horizontal normal to $E_{u}$ and the polar vector $\sigma_{u}$ in the decomposition $D_{H} u=\sigma_{u}\left|D_{H} u\right|$, and it also deals with the relations between $D_{H}^{a} u, D_{H}^{s} u$ and the horizontal derivatives of $\chi_{E_{u}}$. When $m=n$ and $X_{i}=\partial_{x_{i}}$ one recovers some results that belong to the folklore of Geometric Measure Theory and are scattered in the literature (see e.g. [28], [11, 4.5.9] and [17, Section 4.1.5]); we tried here to collect them in a more systematic way. We were not able to find references for some of the results we stated.

Property $\mathscr{R}$ ("rectifiability") intervenes in ensuring that the horizontal derivatives of $\chi_{E_{u}}$ are a "rectifiable" measure, see Definition 5.1. This is a non-trivial technical obstruction one has to face when following the strategy of [27]: the rectifiability of sets with finite $H$-perimeter in Carnot groups is indeed a major open problem, which has been solved only in step 2 Carnot groups (see [14,15]) and in the class of Carnot groups of type $\star([26])$. See also [4] for a partial result in general Carnot groups.

Once the rectifiability of $E_{u}$ is ensured, the proof of Theorem 1.1 follows rather easily from the technical Lemma 3.2 below, which is the natural counterpart of the Lemma in [27]. The latter, however, was proved by utilizing the area formula for maps between

[^1]rectifiable subsets of $\mathbb{R}^{n}$, see e.g. [3]. A similar tool is not available in the context of Carnot groups, a fact which forces us to follow a different path. The proof of Lemma 3.2 is indeed achieved by a covering argument that is based on the following result: we state it and postpone to Section 2 the definitions of property $\mathscr{C}_{k}$, the Hausdorff measure $\mathcal{H}^{d}$, the homogeneous dimension $Q$ of $\mathbb{G}$ and of hypersurfaces of class $C_{H}^{1}$ with their horizontal normal.

Theorem 1.4. Let $k \geq 1$ be an integer, $\mathbb{G}$ a Carnot group satisfying property $\mathscr{C}_{k}$ and let $\Sigma_{1}, \ldots, \Sigma_{k}$ be hypersurfaces of class $C_{H}^{1}$ with horizontal normals $\nu_{1}, \ldots, \nu_{k}$. Let also $x \in \Sigma:=\Sigma_{1} \cap \cdots \cap \Sigma_{k}$ be such that $\nu_{1}(x), \ldots, \nu_{k}(x)$ are linearly independent. Then, there exists an open neighborhood $U$ of $x$ such that

$$
0<\mathcal{H}^{Q-k}(\Sigma \cap U)<\infty
$$

In particular, the measure $\mathcal{H}^{Q-k}$ is $\sigma$-finite on the set

$$
\Sigma^{\pitchfork}:=\left\{x \in \Sigma: \nu_{1}(x), \ldots, \nu_{k}(x) \text { are linearly independent }\right\} .
$$

Theorem 1.4, that we prove in Appendix A, is an easy consequence of Theorems A. 3 and A. 5 proved, respectively, in [12] and [22]. Theorem A.5, in particular, states the much deeper property that $\Sigma^{\pitchfork}$ is locally an intrinsic Lipschitz graph. To this aim, one needs the intersection $T_{x} \Sigma_{1} \cap \cdots \cap T_{x} \Sigma_{k}$ of the tangent subgroups to $\Sigma_{i}$ at $x$ to admit a (necessarily commutative) complementary homogeneous subgroup that is horizontal, i.e., contained in $\exp \left(\mathfrak{g}_{1}\right)$. This algebraic property is guaranteed by property $\mathscr{C}_{k}$ ("complementability"), see Remark 2.3. We will provide in Appendix A a proof of Theorem A. 5 which does not rely on the homotopy invariance of the topological degree and is then simpler and shorter than the one in [22].

For the validity of Theorem 1.4, property $\mathscr{C}_{k}$ might seem a restrictive one. We however point out that Theorem 1.4 is no longer valid already when $k=2$ and $\mathbb{G}$ is the first Heisenberg group $\mathbb{H}^{1}$, which does not satisfy $\mathscr{C}_{2}$ : indeed, in this setting the measure $\mathcal{H}^{Q-2}\left(\Sigma^{\pitchfork}\right)$ might be either 0 or $+\infty$ (even locally) as shown by A. Kozhevnikov [19]. See also [20,25].

The fact that Theorem 1.4 does not apply to $\mathbb{H}^{1}$ (actually, to $\mathbb{H}^{1} \times \mathbb{R} \times \mathbb{R}$, see the proof of Lemma 3.2) prevents us from proving the rank-one Theorem 1.1 for $\mathbb{G}=\mathbb{H}^{1}$. This does not follow from [9] either (see Remark 5.6) and, thus, it remains a very interesting open problem.

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## 2. Preliminaries on Carnot groups

### 2.1. Algebraic facts

A Carnot (or stratified) group is a connected, simply connected and nilpotent Lie group whose Lie algebra $\mathfrak{g}$ is stratified, i.e., it has a decomposition $\mathfrak{g}=\mathfrak{g}_{1} \oplus \cdots \oplus \mathfrak{g}_{s}$ such that

$$
\forall j=1, \ldots, s-1 \quad \mathfrak{g}_{j+1}=\left[\mathfrak{g}_{j}, \mathfrak{g}_{1}\right], \quad \mathfrak{g}_{s} \neq\{0\} \quad \text { and } \quad\left[\mathfrak{g}_{s}, \mathfrak{g}\right]=\{0\} .
$$

We refer to the integer $s$ as the step of $\mathbb{G}$ and to $m:=\operatorname{dimg}_{1}$ as its rank; apart from the case in which $\mathbb{G}$ is a Heisenberg group (see Example 2.1), $n$ denotes the topological dimension of $\mathbb{G}$. The group identity is denoted by 0 .

The exponential map $\exp : \mathfrak{g} \rightarrow \mathbb{G}$ is a diffeomorphism and, given a basis $X_{1}, \ldots, X_{n}$ of $\mathfrak{g}$, we often identify $\mathbb{G}$ with $\mathbb{R}^{n}$ by means of exponential coordinates:

$$
\mathbb{R}^{n} \ni x=\left(x_{1}, \ldots, x_{n}\right) \longleftrightarrow \exp \left(x_{1} X_{1}+\cdots+x_{n} X_{n}\right) \in \mathbb{G} .
$$

A one-parameter family $\left\{\delta_{\lambda}\right\}_{\lambda>0}$ of dilations $\delta_{\lambda}: \mathfrak{g} \rightarrow \mathfrak{g}$ is defined by $\delta_{\lambda}(X):=\lambda^{j} X$ for any $X \in \mathfrak{g}_{j}$; notice that $\delta_{\lambda \mu}=\delta_{\lambda} \circ \delta_{\mu}$. By composition with exp one can then define a one-parameter family, for which we use the same symbol $\delta$, of group isomorphisms $\delta_{\lambda}: \mathbb{G} \rightarrow \mathbb{G}$.

Example 2.1. Apart from Euclidean spaces, which are the only commutative Carnot groups, the most basic examples of Carnot groups are Heisenberg groups. Given an integer $n \geq 1$, the $n$-th Heisenberg group $\mathbb{H}^{n}$ is the $2 n+1$ dimensional Carnot group of step 2 whose Lie algebra is generated by $X_{1}, \ldots, X_{n}, Y_{1}, \ldots, Y_{n}, T$ and the only non-vanishing commutation relations among these generators are given by

$$
\left[X_{j}, Y_{j}\right]=T \quad \text { for any } j=1, \ldots, n
$$

The stratification of the Lie algebra is given by $\mathfrak{g}_{1} \oplus \mathfrak{g}_{2}$, where $\mathfrak{g}_{1}:=\operatorname{span}\left\{X_{j}, Y_{j}: j=\right.$ $1, \ldots n\}$ and $\mathfrak{g}_{2}:=\operatorname{span}\{T\}$. In exponential coordinates

$$
\mathbb{R}^{n} \times \mathbb{R}^{n} \times \mathbb{R} \ni(x, y, t) \longleftrightarrow \exp \left(x_{1} X_{1}+\cdots+y_{n} Y_{n}+t T\right)
$$

one has

$$
X_{j}=\partial_{x_{j}}-\frac{y_{j}}{2} \partial_{t}, \quad Y_{j}=\partial_{y_{j}}+\frac{x_{j}}{2} \partial_{t}, \quad T=\partial_{t}
$$

In this paper, given a Carnot group $\mathbb{G}$ we will frequently deal with products like $\mathbb{G} \times \mathbb{R}^{N}$. Needless to say, this is the Carnot group with algebra $\mathfrak{g} \times \mathbb{R}^{N}$ with product defined by $[(X, t),(Y, s)]=([X, Y], 0)$ for any $X, Y \in \mathfrak{g}, t, s \in \mathbb{R}^{N}$ and whose stratification is given by $\left(\mathfrak{g}_{1} \times \mathbb{R}^{N}\right) \oplus\left(\mathfrak{g}_{2} \times\{0\}\right) \oplus \cdots \oplus\left(\mathfrak{g}_{s} \times\{0\}\right)$.

Definition 2.2. Let $\mathbb{G}$ be a Carnot group with rank $m$ and let $1 \leq k \leq m$ be an integer. We say that $\mathbb{G}$ satisfies the property $\mathscr{C}_{k}$ if the first layer $\mathfrak{g}_{1}$ of its Lie algebra has the following property: for any linear subspace $\mathfrak{w}$ of $\mathfrak{g}_{1}$ of codimension $k$ there exists a commutative complementary subspace in $\mathfrak{g}_{1}$, i.e., a $k$-dimensional subspace $\mathfrak{h}$ of $\mathfrak{g}_{1}$ such that $[\mathfrak{h}, \mathfrak{h}]=0$ and $\mathfrak{g}_{1}=\mathfrak{w} \oplus \mathfrak{h}$.

Remark 2.3. As customary in the literature, we say that $\mathbb{W} \subset \mathbb{G}$ is a vertical plane of codimension $k$ (for some $1 \leq k \leq m)$ if $\mathbb{W}=\exp \left(\mathfrak{w} \oplus \mathfrak{g}_{2} \oplus \cdots \oplus \mathfrak{g}_{s}\right)$ for some linear subspace $\mathfrak{w}$ of $\mathfrak{g}_{1}$ of codimension $k$ (possibly $\mathfrak{w}=\{0\}$ ); such a $\mathbb{W}$ is a homogeneous normal subgroup of $\mathbb{G}$ of topological dimension $n-k$, see also Section 3 below. Then, a Carnot group has the property $\mathscr{C}_{k}$ if and only if, for any vertical plane $\mathbb{W}$ in $\mathbb{G}$, there exists a complementary homogeneous subgroup $\mathbb{H}$ that is horizontal, i.e., such that $\mathbb{H} \subset \exp \left(\mathfrak{g}_{1}\right)$. Notice also that, in this case, $\mathbb{H}$ is necessarily commutative.

Remark 2.4. The Heisenberg group $\mathbb{H}^{n}$ has the property $\mathscr{C}_{k}$ if and only if $1 \leq k \leq n$.
All Carnot groups have the property $\mathscr{C}_{1}$. Free Carnot groups (see e.g. [18]) have the property $\mathscr{C}_{k}$ if and only if $k=1$.

A Carnot group of rank $m$ has the property $\mathscr{C}_{m}$ if and only if $\mathbb{G}$ is Abelian (i.e., $\left.\mathbb{G} \equiv \mathbb{R}^{m}\right)$.

Remark 2.5. It is an easy exercise to show that, if $k \geq 2$ and $\mathbb{G}$ has the property $\mathscr{C}_{k}$, then $\mathbb{G}$ has also the property $\mathscr{C}_{h}$ for any $1 \leq h \leq k$.

Lemma 2.6. Let $N \geq 1$ be an integer and $\mathbb{G}$ be a Carnot group. Then $\mathbb{G}$ has the property $\mathscr{C}_{k}$ if and only if $\mathbb{G} \times \mathbb{R}^{N}$ has the property $\mathscr{C}_{k}$.

Proof. It is clearly enough to prove the statement for $N=1$.
Assume first that $\mathbb{G}$ has the property $\mathscr{C}_{k}$ and let $\mathfrak{w}$ be a $k$-codimensional subspace of the first layer $\mathfrak{g}_{1} \times \mathbb{R}$ of the Lie algebra of $\mathbb{G} \times \mathbb{R}$. We have two cases according to the dimension of $\mathfrak{w}^{\prime}:=\mathfrak{w} \cap\left(\mathfrak{g}_{1} \times\{0\}\right)$ :

- if $\operatorname{dim} \mathfrak{w}^{\prime}=m-k$, using the $\mathscr{C}_{k}$ property of $\mathbb{G}$ one can find a $k$-dimensional commutative subspace $\mathfrak{h}$ of $\mathfrak{g}_{1}$ such that $\mathfrak{g}_{1} \times\{0\}=\mathfrak{w}^{\prime} \oplus(\mathfrak{h} \times\{0\})$. In particular, $\mathfrak{g}_{1} \times \mathbb{R}=\mathfrak{w} \oplus(\mathfrak{h} \times\{0\}) ;$
- if $\operatorname{dim} \mathfrak{w}^{\prime}=m+1-k$, then $\mathfrak{w}=\mathfrak{w}^{\prime} \subset \mathfrak{g}_{1} \times\{0\}$ and, by Remark 2.5, one can find a $(k-1)$-dimensional commutative subspace $\mathfrak{h}$ of $\mathfrak{g}_{1}$ such that $\mathfrak{g}_{1} \times\{0\}=\mathfrak{w} \oplus(\mathfrak{h} \times\{0\})$. In particular, $\mathfrak{g}_{1} \times \mathbb{R}=\mathfrak{w} \oplus(\mathfrak{h} \times \mathbb{R})$.

In both cases we have found a commutative complementary subspace of $\mathfrak{w}$.
Assume now that $\mathbb{G} \times \mathbb{R}$ has the property $\mathscr{C}_{k}$ and let $\mathfrak{w}$ be a $k$-codimensional linear subspace of $\mathfrak{g}_{1}$. Then $\mathfrak{w} \times \mathbb{R}$ is a $k$-codimensional linear subspace of $\mathfrak{g}_{1} \times \mathbb{R}$, hence it admits a $k$-dimensional commutative complementary subspace $\mathfrak{h}$ in $\mathfrak{g}_{1} \times \mathbb{R}$. Denoting by
$\pi: \mathfrak{g}_{1} \times \mathbb{R} \rightarrow \mathfrak{g}_{1}$ the canonical projection, it is readily noticed that $\pi(\mathfrak{h})$ is a $k$-dimensional commutative subspace of $\mathfrak{g}_{1}$ such that $\mathfrak{g}_{1}=\mathfrak{w} \oplus \pi(\mathfrak{h})$. This concludes the proof.

### 2.2. Metric facts

Let $\mathbb{G}$ be a Carnot group with stratified algebra $\mathfrak{g}=\mathfrak{g}_{1} \oplus \cdots \oplus \mathfrak{g}_{s}$. We endow $\mathfrak{g}$ with a positive definite scalar product $\langle\cdot, \cdot\rangle$ such that $\mathfrak{g}_{i} \perp \mathfrak{g}_{j}$ whenever $i \neq j$. We also let $|\cdot|:=\langle\cdot, \cdot\rangle^{1 / 2}$. We fix an orthonormal basis $X_{1}, \ldots, X_{n}$ of $\mathfrak{g}$ adapted to the stratification, i.e., such that $\mathfrak{g}_{j}=\operatorname{span}\left\{X_{m_{j-1}+1}, \ldots, X_{m_{j}}\right\}$ for any $j=1, \ldots, s$, where $m_{j}:=\operatorname{dim}\left(\mathfrak{g}_{1}\right)+\cdots+\operatorname{dim}\left(\mathfrak{g}_{j}\right)$ and $m_{0}:=0$ (in particular, $m_{1}=m$ ).

We will frequently use the homogeneous (pseudo-) norm $\|\cdot\|$ on $\mathbb{G}$ defined in this way: if $x=\exp \left(Y_{1}+\cdots+Y_{s}\right)$ for $Y_{j} \in \mathfrak{g}_{j}$, then

$$
\|x\|:=\sum_{j=1}^{s}\left|Y_{j}\right|^{1 / j}
$$

Clearly one has $\left\|\delta_{\lambda}(x)\right\|=\lambda\|x\|$ for any $x \in \mathbb{G}, \lambda>0$. Homogeneous pseudo-norms arising from different choices of the scalar product $\langle\cdot, \cdot\rangle$ on $\mathbb{G}$ are equivalent.

The group $\mathbb{G}$ is endowed with the Carnot-Carathéodory ( $C C$ ) distance $d$ induced by the family $X_{1}, \ldots, X_{m}$, as we now introduce. Given an interval $I \subset \mathbb{R}$, a Lipschitz curve $\gamma: I \rightarrow \mathbb{G}$ is said to be horizontal if there exist functions $h_{1}, \ldots, h_{m} \in L^{\infty}(I)$ such that for a.e. $t \in I$ we have

$$
\begin{equation*}
\dot{\gamma}(t)=\sum_{i=1}^{m} h_{i}(t) X_{i}(\gamma(t)) . \tag{2.1}
\end{equation*}
$$

Letting $|h|:=\left(h_{1}^{2}+\ldots+h_{m}^{2}\right)^{1 / 2}$, the length of $\gamma$ is defined as

$$
L(\gamma):=\int_{I}|h(t)| d t
$$

It is well-known that for any pair of points $x, y \in \mathbb{G}$ there exists a horizontal curve joining $x$ to $y$. We can therefore define a distance function $d$ letting

$$
d(x, y):=\inf \{L(\gamma): \gamma:[0, T] \rightarrow M \text { horizontal with } \gamma(0)=x \text { and } \gamma(T)=y\} .
$$

It is also well-known that, for any pair $x, y \in \mathbb{G}$, there exists a geodesic joining $x$ and $y$, i.e., a horizontal curve $\gamma$ realizing the infimum in the previous formula. Notice that

$$
d(z x, z y)=d(x, y) \quad \text { and } \quad d\left(\delta_{\lambda}(x), \delta_{\lambda}(y)\right)=\lambda d(x, y) \quad \forall x, y, z \in G, \lambda>0
$$

and that $d(x, y)$ is equivalent to $\left\|x^{-1} y\right\|$.

We denote by $B(x, r)$ open balls of center $x \in \mathbb{G}$ and radius $r>0$ with respect to the CC distance; we also write $B_{r}$ instead of $B(0, r)$, so that $B(x, r)=x B_{r}$. The diameter diam $E$ of $E \subset \mathbb{G}$ and the distance $d\left(E_{1}, E_{2}\right)$ between $E_{1}, E_{2} \subset \mathbb{G}$ is understood with respect to the CC distance.

As customary, for $E \subset \mathbb{G}, d>0$ and $\delta>0$ we set

$$
\begin{aligned}
& \mathcal{H}_{\delta}^{d}(E):=\inf \left\{\sum_{i=1}^{\infty}\left(\operatorname{diam} E_{i}\right)^{d}: E \subset \bigcup_{i=1}^{\infty} E_{i}, \operatorname{diam} E_{i}<\delta\right\} \\
& \mathcal{S}_{\delta}^{d}(E):=\inf \left\{\sum_{i=1}^{\infty}\left(\operatorname{diam} B_{i}\right)^{d}: B_{i} \text { are open balls, } E \subset \bigcup_{i=1}^{\infty} B_{i}, \operatorname{diam} B_{i}<\delta\right\}
\end{aligned}
$$

and we define the $d$-dimensional Hausdorff measure and $d$-dimensional spherical Hausdorff measure of $E$ respectively as

$$
\begin{aligned}
\mathcal{H}^{d}(E) & :=\lim _{\delta \downarrow 0} \mathcal{H}_{\delta}^{d}(E)=\sup _{\delta>0} \mathcal{H}_{\delta}^{d}(E) \\
\mathcal{S}^{d}(E) & :=\lim _{\delta \downarrow 0} \mathcal{S}_{\delta}^{d}(E)=\sup _{\delta>0} \mathcal{S}_{\delta}^{d}(E)
\end{aligned}
$$

The Hausdorff dimension of $E$ is $\inf \left\{d: \mathcal{H}^{d}(E)=0\right\}=\sup \left\{d: \mathcal{H}^{d}(E)=\infty\right\}$. It is well-known that the metric space $(\mathbb{G}, d)$ has Hausdorff dimension $Q:=\sum_{j=1}^{s} j \operatorname{dim} \mathfrak{g}_{j}$ and that, in exponential coordinates and up to multiplicative constants, the measures $\mathcal{H}^{Q}, \mathcal{S}^{Q}$ and $\mathscr{L}^{n}$ coincide, all of them being Haar measures on $\mathbb{G}$.

## 3. Intrinsic regular hypersurfaces in Carnot groups

We say that a continuous real function $f$ on an open set $\Omega \subset \mathbb{G}$ is of class $C_{H}^{1}$ if its horizontal derivatives $X_{1} f, \ldots, X_{m} f$ are continuous in $\Omega$. In this case we write $f \in C_{H}^{1}(\Omega)$ and we set $\nabla_{H} f:=\left(X_{1} f, \ldots, X_{m} f\right)$.

A set $S \subset \mathbb{G}$ is a $C_{H}^{1}$ hypersurface if for any $x \in S$ there exist an open neighborhood $U$ of $x$ and $f \in C_{H}^{1}(U)$ such that

$$
S \cap U=\{y \in U: f(y)=0\} \quad \text { and } \quad \nabla_{H} f \neq 0 \text { on } U .
$$

In this case, we define the horizontal normal to $x$ as $\nu_{S}(x):=\frac{\nabla_{H} f(x)}{\left|\nabla_{H} f(x)\right|} \in \mathbb{R}^{m}$. The normal $\nu_{S}(x)=\left(\left(\nu_{S}(x)\right)_{1}, \ldots,\left(\nu_{S}(x)\right)_{m}\right)$ is defined up to sign and it can be canonically identified with a horizontal vector at $x$ by

$$
\nu_{S}(x)=\left(\nu_{S}(x)\right)_{1} X_{1}(x)+\cdots+\left(\nu_{S}(x)\right)_{m} X_{m}(x) .
$$

A $C_{H}^{1}$ hypersurface has locally finite $\mathcal{H}^{Q-1}$-measure, see e.g. [30] and the references therein. ${ }^{1}$

The hyperplane $\nu_{S}(x)^{\perp}$ in $\mathfrak{g}$ is a Lie subalgebra. The associated subgroup $T_{x} S:=$ $\exp \left(\nu_{S}(x)^{\perp}\right)$ is called tangent subgroup to $S$ at $x$ : we point out the well-known property that

$$
\begin{equation*}
\forall \varepsilon>0 \exists \bar{r}=\bar{r}(x, \varepsilon)>0 \text { such that } \forall r \in(0, \bar{r}) \quad\left(x^{-1} S\right) \cap B_{r} \subset\left(T_{x} S\right)_{\varepsilon r} \cap B_{r}, \tag{3.2}
\end{equation*}
$$

where for $E \subset \mathbb{G}$ and $\delta>0$ we denote by $E_{\delta}$ the $\delta$-neighborhood of $E$. A proof of (3.2), using the fact that in exponential coordinates $T_{x} S=\left\{(\xi, \eta) \in \mathbb{R}^{n}=\mathbb{R}^{m} \times \mathbb{R}^{n-m}: \xi \perp\right.$ $\left.\nu_{S}(x)\right\}$, is implicitly contained in the proof of Lemma A.4. Notice also that

$$
T_{x} S=\exp \left(\left\{X \in \mathfrak{g}_{1}: X f(x)=0\right\} \oplus \mathfrak{g}_{2} \cdots \oplus \mathfrak{g}_{s}\right)
$$

in particular, while $\nu_{S}(x)$ depends on the scalar product $\langle\cdot, \cdot\rangle$ on $\mathfrak{g}$, the subgroup $T_{x} S$ is intrinsic.

The tangent group $T_{x} S$ is a vertical plane of codimension 1 (or vertical hyperplane), where we say that $\mathbb{W} \subset \mathbb{G}$ is a vertical plane of codimension $k, 1 \leq k \leq m$, if $\mathbb{W}=$ $\exp \left(\mathfrak{w} \oplus \mathfrak{g}_{2} \oplus \cdots \oplus \mathfrak{g}_{s}\right)$ for some linear subspace $\mathfrak{w}$ of $\mathfrak{g}_{1}$ of codimension $k$ (possibly $\mathfrak{w}=\{0\}$ ). Such a $\mathbb{W}$ is a homogeneous normal subgroup of $\mathbb{G}$ of topological dimension $n-k$ and Hausdorff dimension $Q-k$. The intersection of vertical planes is always a vertical plane (of possibly higher codimension).

The following simple lemma will be used in the proof of Lemma 3.2.

Lemma 3.1. Let $\mathbb{W} \subset \mathbb{G}$ be a vertical plane of codimension $k$ and let $x \in \mathbb{W}, r>0$ and $\varepsilon \in(0,1)$ be fixed. Then, the set $\mathbb{W} \cap B(x, r)$ can be covered by a family of balls $\left\{B\left(y_{\ell}, \varepsilon r\right)\right\}_{\ell \in L}$ of radius $\varepsilon r$ with cardinality $\# L \leq(4 / \varepsilon)^{Q-k}$.

Proof. By dilation and translation invariance, it is not restrictive to assume that $x=0$ and $r=1$. Let $\left\{y_{\ell}\right\}_{\ell \in L}$ be a maximal family of points of $\mathbb{W} \cap B(0,1)$ such that the balls $B\left(y_{\ell}, \varepsilon / 2\right)$ are pairwise disjoint; working by contradiction, it can be easily seen that the family $\left\{B\left(y_{\ell}, \varepsilon\right)\right\}_{\ell \in L}$ covers $\mathbb{W} \cap B(0,1)$. The measure $\mathcal{H}^{Q-k}$ is locally finite on $\mathbb{W}$ (see e.g. $[21,23,24])$, is left-invariant and it is $(Q-k)$-homogeneous with respect to dilations. In particular, setting $M:=\mathcal{H}^{Q-k}(\mathbb{W} \cap B(0,1))$, we have

$$
\left(\frac{\varepsilon}{2}\right)^{Q-k} M \# L=\sum_{\ell \in L} \mathcal{H}^{Q-k}\left(\mathbb{W} \cap B\left(y_{\ell}, \varepsilon / 2\right)\right) \leq \mathcal{H}^{Q-k}(\mathbb{W} \cap B(0,2))=2^{Q-k} M
$$

which proves the claim.

[^2]A key tool in the proof of the rank-one Theorem 1.1 is the following Lemma 3.2 which, in turn, uses Theorem 1.4, whose proof is instead postponed to Appendix A. We denote by $\pi: \mathbb{G} \times \mathbb{R} \rightarrow \mathbb{G}$ the canonical projection $\pi(x, t)=x$.

Lemma 3.2. Let $\mathbb{G}$ be a Carnot group satisfying property $\mathscr{C}_{2}$. Let $\Sigma_{1}, \Sigma_{2}$ be $C_{H}^{1}$ hypersurfaces in $\mathbb{G} \times \mathbb{R}$ with unit normals $\nu_{\Sigma_{1}}, \nu_{\Sigma_{2}}$. Then, the set

$$
R:=\left\{\begin{array}{ll} 
& \pi(q)=\pi(p), \\
p \in \Sigma_{1}: \exists q \in \Sigma_{2} \text { such that }\left(\nu_{\Sigma_{1}}(p)\right)_{m+1}=\left(\nu_{\Sigma_{2}}(q)\right)_{m+1}=0 \\
& \nu_{\Sigma_{1}}(p) \neq \pm \nu_{\Sigma_{2}}(q)
\end{array}\right\}
$$

is $\mathcal{H}^{Q}$-negligible.
Proof. Let us consider the distances $d_{\mathbb{G} \times \mathbb{R}}$ and $d_{\mathbb{G} \times \mathbb{R} \times \mathbb{R}}$ on (respectively) $\mathbb{G} \times \mathbb{R}$ and $\mathbb{G} \times \mathbb{R} \times \mathbb{R}$ defined by

$$
\begin{array}{ll}
d_{\mathbb{G} \times \mathbb{R}}\left((x, t),\left(x^{\prime}, t^{\prime}\right)\right):=d\left(x, x^{\prime}\right)+\left|t-t^{\prime}\right| & \forall x, x^{\prime} \in \mathbb{G}, t, t^{\prime} \in \mathbb{R} \\
d_{\mathbb{G} \times \mathbb{R} \times \mathbb{R}}\left((x, t, s),\left(x^{\prime}, t^{\prime}, s^{\prime}\right)\right):=d\left(x, x^{\prime}\right)+\left|t-t^{\prime}\right|+\left|s-s^{\prime}\right| & \forall x, x^{\prime} \in \mathbb{G}, t, t^{\prime}, s, s^{\prime} \in \mathbb{R},
\end{array}
$$

where $d$ is the Carnot-Carathéodory distance on $\mathbb{G}$. Such distances are left-invariant and homogeneous, hence they are equivalent to the Carnot-Carathéodory distances on $\mathbb{G} \times \mathbb{R}$ and $\mathbb{G} \times \mathbb{R} \times \mathbb{R}$; in particular, it is enough to prove the statement when the Hausdorff measure $\mathcal{H}^{Q}$ is the one induced by $d_{\mathbb{G} \times \mathbb{R}}$ on $\mathbb{G} \times \mathbb{R}$. We use the same notation $B(a, r)$ for balls of radius $r>0$ in either $\mathbb{G}, \mathbb{G} \times \mathbb{R}$ or $\mathbb{G} \times \mathbb{R} \times \mathbb{R}$, according to which group the center $a$ belongs to.

The sets

$$
\begin{aligned}
& \widetilde{\Sigma}_{1}:=\left\{(x, t, s) \in \mathbb{G} \times \mathbb{R} \times \mathbb{R}:(x, t) \in \Sigma_{1}, s \in \mathbb{R}\right\} \\
& \widetilde{\Sigma}_{2}:=\left\{(x, t, s) \in \mathbb{G} \times \mathbb{R} \times \mathbb{R}:(x, s) \in \Sigma_{2}, t \in \mathbb{R}\right\}
\end{aligned}
$$

are clearly $C_{H}^{1}$ hypersurfaces in $\mathbb{G} \times \mathbb{R} \times \mathbb{R}$ and, moreover,

$$
\begin{aligned}
& \nu_{\widetilde{\Sigma}_{1}}(x, t, s)=\left(\left(\nu_{\Sigma_{1}}(x, t)\right)_{1}, \ldots,\left(\nu_{\Sigma_{1}}(x, t)\right)_{m},\left(\nu_{\Sigma_{1}}(x, t)\right)_{m+1}, 0\right) \\
& \nu_{\widetilde{\Sigma}_{2}}(x, t, s)=\left(\left(\nu_{\Sigma_{2}}(x, s)\right)_{1}, \ldots,\left(\nu_{\Sigma_{2}}(x, s)\right)_{m}, 0,\left(\nu_{\Sigma_{2}}(x, s)\right)_{m+1}\right) .
\end{aligned}
$$

Let us define

$$
\begin{aligned}
\widetilde{R}:= & \left\{P \in \widetilde{\Sigma}_{1} \cap \widetilde{\Sigma}_{2}:\left(\nu_{\widetilde{\Sigma}_{1}}(P)\right)_{m+1}=\left(\nu_{\widetilde{\Sigma}_{2}}(P)\right)_{m+2}=0 \text { and } \nu_{\widetilde{\Sigma}_{1}}(P) \neq \pm \nu_{\widetilde{\Sigma}_{2}}(P)\right\} \\
= & \left\{(x, t, s) \in \widetilde{\Sigma}_{1} \cap \widetilde{\Sigma}_{2}:\left(\nu_{\Sigma_{1}}(x, t)\right)_{m+1}=\left(\nu_{\Sigma_{2}}(x, s)\right)_{m+1}=0\right. \\
& \text { and } \left.\nu_{\Sigma_{1}}(x, t) \neq \pm \nu_{\Sigma_{2}}(x, s)\right\} .
\end{aligned}
$$

By construction we have $\widetilde{\pi}(\widetilde{R})=R$, where $\widetilde{\pi}: \mathbb{G} \times \mathbb{R} \times \mathbb{R} \rightarrow \mathbb{G} \times \mathbb{R}$ is the group homomorphism defined by $\widetilde{\pi}(x, t, s):=(x, t)$; moreover the measure $\mathcal{H}^{Q}\llcorner\widetilde{R}$ is $\sigma$-finite by Theorem 1.4 (notice that we are also using Lemma 2.6). We are going to show that $\mathcal{H}^{Q}(\widetilde{\pi}(T))=0$ for any fixed $T \subset \widetilde{R}$ such that $\mathcal{S}^{Q}(T)<\infty$; this is clearly enough to conclude.

For any $P \in T$ and $i=1,2$, the tangent space $T_{P} \widetilde{\Sigma}_{i}$ equals $\mathbb{W}_{i} \times \mathbb{R} \times \mathbb{R}$ for a suitable vertical hyperplane $\mathbb{W}_{i}$ of $\mathbb{G}$. In particular, setting $\mathbb{W}=\mathbb{W}(P):=\mathbb{W}_{1} \cap \mathbb{W}_{2}$, we have by (3.2) that for any $P \in T$ and any $\varepsilon \in(0,1)$ there exists $\bar{r}=\bar{r}(\varepsilon, P)>0$ such that

$$
\begin{align*}
\left(P^{-1} T\right) \cap B(0, r) & \subset(\mathbb{W} \times \mathbb{R} \times \mathbb{R})_{\varepsilon r} \cap B(0, r)  \tag{3.3}\\
& =\left(\mathbb{W}_{\varepsilon r} \times \mathbb{R} \times \mathbb{R}\right) \cap B(0, r) \quad \text { for any } r \in(0, \bar{r}) .
\end{align*}
$$

Notice also that $\mathbb{W}$ is a vertical plane of codimension 2 in $\mathbb{G}$. Let $\varepsilon>0$ be fixed and set

$$
T_{j}:=\left\{P \in T: \bar{r}(\varepsilon, P) \geq \frac{1}{j}\right\}, \quad j=1,2, \ldots
$$

Since $T_{j} \uparrow T$, the proof will be accomplished by showing that for any fixed $j$

$$
\begin{equation*}
\mathcal{H}^{Q}\left(\widetilde{\pi}\left(T_{j}\right)\right)<C \varepsilon, \tag{3.4}
\end{equation*}
$$

where $C>0$ is a constant that will be determined in the sequel.
Let us prove (3.4). Fix $\delta \in\left(0, \frac{1}{j}\right)$; since $\mathcal{H}^{Q}\left(T_{j}\right) \leq \mathcal{H}^{Q}(T)<+\infty$, one can find a (countable or finite) family $\left\{B\left(\widetilde{P}_{i}, r_{i} / 2\right)\right\}_{i}$ of balls in $\mathbb{G} \times \mathbb{R} \times \mathbb{R}$ such that $0<r_{i}<\delta$,

$$
T_{j} \subset \bigcup_{i} B\left(\widetilde{P}_{i}, r_{i} / 2\right) \quad \text { and } \quad \sum_{i}\left(r_{i} / 2\right)^{Q} \leq \sum_{i}\left(\operatorname{diam} B\left(\widetilde{P}_{i}, r_{i} / 2\right)\right)^{Q} \leq C_{1}
$$

where $C_{1}:=\mathcal{H}^{Q}(T)+1$. We can also assume that $T_{j} \cap B\left(\widetilde{P}_{i}, r_{i} / 2\right)$ is non-empty for any $i$. Choosing $P_{i} \in T_{j} \cap B\left(\widetilde{P}_{i}, r_{i} / 2\right)$, for any $i$ the balls $B\left(P_{i}, r_{i}\right)$ have then the following properties:

$$
\begin{equation*}
P_{i} \in T_{j}, \quad 0<r_{i}<\delta, \quad T_{j} \subset \bigcup_{i} B\left(P_{i}, r_{i}\right) \quad \text { and } \quad \sum_{i} r_{i}^{Q} \leq 2^{Q} C_{1} \tag{3.5}
\end{equation*}
$$

Setting $\mathbb{W}_{i}:=\mathbb{W}\left(P_{i}\right)$, by (3.3) we have

$$
\begin{align*}
\left(P_{i}^{-1} T_{j}\right) \cap B\left(0, r_{i}\right) & \subset\left(\left(\mathbb{W}_{i}\right)_{\varepsilon r_{i}} \times \mathbb{R} \times \mathbb{R}\right) \cap B\left(0, r_{i}\right)  \tag{3.6}\\
& =\left(\left(\mathbb{W}_{i}\right)_{\varepsilon r_{i}} \cap B\left(0, r_{i}\right)\right) \times\left(-r_{i}, r_{i}\right) \times\left(-r_{i}, r_{i}\right)
\end{align*}
$$

By Lemma 3.1, for any $i$ we can find a family of balls $\left\{B\left(y_{i, \ell}, \varepsilon r_{i}\right)\right\}_{\ell \in L_{i}}$ such that

$$
\forall \ell \in L_{i} y_{i, \ell} \in \mathbb{W}_{i}, \quad \# L_{i} \leq(8 / \varepsilon)^{Q-2} \quad \text { and } \quad \mathbb{W}_{i} \cap B\left(0,2 r_{i}\right) \subset \bigcup_{\ell \in L_{i}} B\left(y_{i, \ell}, \varepsilon r_{i}\right)
$$

In particular

$$
\begin{equation*}
\left(\mathbb{W}_{i}\right)_{\varepsilon r_{i}} \cap B\left(0, r_{i}\right) \subset\left(\mathbb{W}_{i} \cap B\left(0, r_{i}+\varepsilon r_{i}\right)\right)_{\varepsilon r_{i}} \subset \bigcup_{\ell \in L_{i}} B\left(y_{i, \ell}, 2 \varepsilon r_{i}\right) \tag{3.7}
\end{equation*}
$$

Let us also fix points $\left\{\tau_{k}\right\}_{k \in K_{i}} \subset\left(-r_{i}, r_{i}\right)$ such that $\# K_{i} \leq 2 \varepsilon^{-1}$ and

$$
\begin{equation*}
\left(-r_{i}, r_{i}\right) \subset \bigcup_{k \in K_{i}}\left(\tau_{k}-2 \varepsilon r_{i}, \tau_{k}+2 \varepsilon r_{i}\right) \tag{3.8}
\end{equation*}
$$

By (3.6), (3.7) and (3.8) we get
$\left(P_{i}^{-1} T_{j}\right) \cap B\left(0, r_{i}\right) \subset \bigcup_{\substack{\ell \in L_{i} \\ k, h \in K_{i}}} B\left(y_{i, \ell}, 2 \varepsilon r_{i}\right) \times\left(\tau_{k}-2 \varepsilon r_{i}, \tau_{k}+2 \varepsilon r_{i}\right) \times\left(\tau_{h}-2 \varepsilon r_{i}, \tau_{h}+2 \varepsilon r_{i}\right)$.
For any $\ell \in L_{i}$ and $k, h, h^{\prime} \in K_{i}$ one has

$$
\begin{aligned}
& \widetilde{\pi}\left(B\left(y_{i, \ell}, 2 \varepsilon r_{i}\right) \times\left(\tau_{k}-2 \varepsilon r_{i}, \tau_{k}+2 \varepsilon r_{i}\right) \times\left(\tau_{h}-2 \varepsilon r_{i}, \tau_{h}+2 \varepsilon r_{i}\right)\right) \\
= & \widetilde{\pi}\left(B\left(y_{i, \ell}, 2 \varepsilon r_{i}\right) \times\left(\tau_{k}-2 \varepsilon r_{i}, \tau_{k}+2 \varepsilon r_{i}\right) \times\left(\tau_{h^{\prime}}-2 \varepsilon r_{i}, \tau_{h^{\prime}}+2 \varepsilon r_{i}\right)\right) \\
= & B\left(y_{i, \ell}, 2 \varepsilon r_{i}\right) \times\left(\tau_{k}-2 \varepsilon r_{i}, \tau_{k}+2 \varepsilon r_{i}\right) \\
\subset & B\left(\left(y_{i, \ell}, \tau_{k}\right), 4 \varepsilon r_{i}\right)
\end{aligned}
$$

which, using (3.5), implies that

$$
\begin{aligned}
\widetilde{\pi}\left(T_{j}\right) & \subset \bigcup_{i} \widetilde{\pi}\left(T_{j} \cap B\left(P_{i}, r_{i}\right)\right) \\
& \subset \bigcup_{i} \bigcup_{\substack{\ell \in L_{i} \\
k, h \in K_{i}}} \widetilde{\pi}\left(P_{i}\left(B\left(y_{i, \ell}, 2 \varepsilon r_{i}\right) \times\left(\tau_{k}-2 \varepsilon r_{i}, \tau_{k}+2 \varepsilon r_{i}\right) \times\left(\tau_{h}-2 \varepsilon r_{i}, \tau_{h}+2 \varepsilon r_{i}\right)\right)\right) \\
& \subset \bigcup_{i} \bigcup_{\substack{\ell \in L_{i} \\
k \in K_{i}}} \widetilde{\pi}\left(P_{i}\right) B\left(\left(y_{i, \ell}, \tau_{k}\right), 4 \varepsilon r_{i}\right) \\
& =\bigcup_{i} \bigcup_{\substack{\ell \in L_{i} \\
k \in K_{i}}} B\left(p_{i \ell k}, 4 \varepsilon r_{i}\right)
\end{aligned}
$$

where $p_{i \ell k}:=\widetilde{\pi}\left(P_{i}\right)\left(y_{i, \ell}, \tau_{k}\right) \in \mathbb{G} \times \mathbb{R}$. Using again (3.5) we obtain that

$$
\mathcal{H}_{4 \varepsilon \delta}^{Q}\left(T_{j}\right) \leq \sum_{i} \# L_{i} \# K_{i}\left(8 \varepsilon r_{i}\right)^{Q} \leq \varepsilon \sum_{i} 2^{6 Q-5} r_{i}^{Q} \leq 2^{7 Q-5} C_{1} \varepsilon
$$

which, by the arbitrariness of $\delta \in\left(0, \frac{1}{j}\right)$, gives the claim (3.4).

## 4. Functions with bounded $\boldsymbol{H}$-variation and subgraphs

Let $X=\left(X_{1}, \ldots, X_{m}\right)$ be an $m$-tuple of linearly independent vector fields in $\mathbb{R}^{n}$; for $i=1, \ldots, m$ and $j=1, \ldots, n$ we consider smooth functions $a_{i j}$ such that

$$
X_{i}(x)=\sum_{j=1}^{n} a_{i j}(x) \partial_{x_{j}}
$$

The model case is of course that of a Carnot group $\mathbb{G} \equiv \mathbb{R}^{n}$ endowed with a left-invariant basis $X_{1}, \ldots, X_{m}$ of the first layer $\mathfrak{g}_{1}$ in the Lie algebra stratification; in the present section, however, we work in higher generality.

One of the main purposes of this paper is the study of functions with bounded $H$-variation $([6,13])$, that we are going to introduce only very briefly. In this section, $\Omega$ is an open subset of $\mathbb{R}^{n}$ and, given $\varphi \in C^{1}\left(\Omega, \mathbb{R}^{m}\right)$, we let $\operatorname{div}_{X} \varphi:=\sum_{i=1}^{m} X_{i}^{*} \varphi_{i}$ where $X_{i}^{*}$ denotes the formal adjoint operator of the vector field $X_{i}$. Given a $\mathbb{R}^{m}$-valued function $f$ on $\Omega$ and a $\mathbb{R}^{m}$-valued measure $\mu$ on $\Omega$ we use the compact notation $\int_{\Omega} f \cdot d \mu$ for the sum $\int_{\Omega} f_{1} d \mu_{1}+\cdots+\int_{\Omega} f_{m} d \mu_{m}$.

Definition 4.1. We say that $u \in L_{l o c}^{1}(\Omega)$ is a function of locally bounded $H$-variation in $\Omega$, and we write $u \in B V_{H, l o c}(\Omega)$, if there exists a vector valued Radon measure $D_{H} u=$ $\left(D_{X_{1}} u, \ldots, D_{X_{m}} u\right)$ with locally finite total variation such that for every $\varphi \in C_{c}^{1}\left(\Omega ; \mathbb{R}^{m}\right)$ we have

$$
\begin{equation*}
\int_{\Omega} \varphi \cdot d D_{H} u=-\int_{\Omega} u \operatorname{div}_{X} \varphi d \mathscr{L}^{n} \tag{4.9}
\end{equation*}
$$

Moreover, if $u \in L^{1}(\Omega)$, we say that $u$ has bounded $H$-variation in $\Omega\left(u \in B V_{H}(\Omega)\right)$ if $D_{H} u$ has finite total variation $\left|D_{H} u\right|$ on $\Omega$.

We say that $E \subset \Omega$ has finite $H$-perimeter in $\Omega$ if its characteristic function $\chi_{E}$ belongs to $B V_{H}(\Omega)$.

We recall that the total variation $|\mu|$ of a $\mathbb{R}^{d}$-valued measure $\mu=\left(\mu_{1}, \ldots, \mu_{d}\right)$ is defined for Borel sets $B$ as

$$
\begin{aligned}
|\mu|(B) & :=\sup \left\{\sum_{\ell=1}^{\infty}\left|\mu\left(B_{\ell}\right)\right|:\left(B_{\ell}\right)_{\ell} \text { disjoint Borel subsets of } B\right\} \\
& =\sup \left\{\int_{B} \varphi \cdot d \mu: \varphi: B \rightarrow \mathbb{R}^{d} \text { Borel function, }|\varphi| \leq 1\right\}
\end{aligned}
$$

If $A \Subset \Omega$ is open and $u \in B V_{H, l o c}(\Omega)$, one can easily prove that

$$
\left|D_{H} u\right|(A)=\sup \left\{\int_{A} u \operatorname{div}_{X} \varphi d \mathscr{L}^{n}: \varphi \in C_{c}^{1}\left(A ; \mathbb{R}^{m}\right),|\varphi| \leq 1\right\}
$$

actually, $u \in B V_{H}(A)$ if and only if the supremum on the right-hand side is finite. The total variation is lower-semicontinuous with respect to the $L_{l o c}^{1}$ convergence; moreover (see $[16,13]$ ), for any $u \in B V_{H}(\Omega)$ there exists a sequence $\left(u_{h}\right)_{h}$ in $C^{\infty}(\Omega) \cap B V_{H}(\Omega)$ such that

$$
\begin{align*}
& u_{h} \rightarrow u \text { in } L^{1}(\Omega) \\
& \left|D_{H} u_{h}\right|(\Omega) \rightarrow\left|D_{H} u\right|(\Omega) \\
& \left|D_{X_{i}} u_{h}\right|(\Omega) \rightarrow\left|D_{X_{i}} u\right|(\Omega) \quad \forall i=1, \ldots, m  \tag{4.10}\\
& \left|\left(D_{H} u_{h}, \mathscr{L}^{n}\right)\right|(\Omega) \rightarrow\left|\left(D_{H} u, \mathscr{L}^{n}\right)\right|(\Omega) .
\end{align*}
$$

The aim of this section is the study of the relations occurring between a function $u \in B V_{H}(\Omega)$ and its subgraph

$$
E_{u}:=\{(x, t) \in \Omega \times \mathbb{R}: t<u(x)\} \subset \Omega \times \mathbb{R} .
$$

We introduce the family $\widetilde{X}=\left(\widetilde{X}_{1}, \ldots, \widetilde{X}_{m+1}\right)$ of linearly independent vector fields in $\mathbb{R}^{n+1}$ defined for $(x, t) \in \mathbb{R}^{n} \times \mathbb{R}$ by

$$
\begin{aligned}
& \widetilde{X}_{i}(x, t):=\left(X_{i}(x), 0\right) \in \mathbb{R}^{n+1} \equiv \mathbb{R}^{n} \times \mathbb{R} \quad \text { if } i=1, \ldots, m \\
& \widetilde{X}_{m+1}(x, t):=\partial_{t}
\end{aligned}
$$

If $U \subset \mathbb{R}^{n+1}$ is open and $u \in B V_{H, l o c}(U)$ with respect to the family $\widetilde{X}$ we write $D_{\widetilde{H}} u:=$ $\left(D_{\widetilde{X}_{1}} u, \ldots, D_{\widetilde{X}_{m+1}} u\right)$.

The following result is the natural generalization of some classical facts about Euclidean functions of bounded variation, see e.g. [17, Section 4.1.5]. We denote by $\pi: \mathbb{R}^{n+1} \rightarrow \mathbb{R}^{n}$ the canonical projection $\pi(x, t)=x ; \pi_{\#}$ denotes the associated pushforward of measures.

Theorem 4.2. Suppose $\Omega$ is bounded in $\mathbb{R}^{n}$ and let $u \in L^{1}(\Omega)$. Then $u$ belongs to $B V_{H}(\Omega)$ if and only if its subgraph $E_{u}$ has finite $H$-perimeter (with respect to the family $\widetilde{X}$ ) in $\Omega \times \mathbb{R}$.

Moreover, writing $D_{\widetilde{H}}^{\prime} \chi_{E_{u}}:=\left(D_{\widetilde{X}_{1}} \chi_{E_{u}}, \ldots, D_{\widetilde{X}_{m}} \chi_{E_{u}}\right)$, then the following statements hold:
(i) $\pi_{\#} D_{\widetilde{X}_{i}} \chi_{E_{u}}=D_{X_{i}} u$ for any $i=1, \ldots, m$;
(ii) $\pi_{\#} \partial_{t} \chi_{E_{u}}=-\mathscr{L}^{n}$;
(iii) $\pi_{\#}\left|D_{\widetilde{X}_{i}} \chi_{E_{u}}\right|=\left|D_{X_{i}} u\right|$ for any $i=1, \ldots, m$;
(iv) $\pi_{\#}\left|\partial_{t} \chi_{E_{u}}\right|=\mathscr{L}^{n}$;
(v) $\pi_{\#}\left|D_{\widetilde{H}}^{\prime} \chi_{E_{u}}\right|=\left|D_{H} u\right|$;
(vi) $\pi_{\#}\left|D_{\widetilde{H}} \chi_{E_{u}}\right|=\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|$.

Proof. Suppose first that $\chi_{E_{u}} \in B V_{H}(\Omega \times \mathbb{R})$ with respect to the family $\widetilde{X}$. We need to fix a sequence $\left(g_{h}\right)_{h}$ in $C_{c}^{\infty}(\mathbb{R})$ such that $g_{h}$ is even, $g_{h} \equiv 1$ on $[0, h], g_{h} \equiv 0$ on $[h+1,+\infty)$ and $\int_{\mathbb{R}} g_{h}(t) d t=2 h+1$. Let $\varphi \in C_{c}^{1}\left(\Omega, \mathbb{R}^{m}\right)$ with $|\varphi| \leq 1$ be fixed. By the Dominated Convergence Theorem we have

$$
\begin{aligned}
\int_{\Omega \times \mathbb{R}} \varphi(x) \cdot d\left(D_{\widetilde{H}}^{\prime} \chi_{E_{u}}\right)(x, t) & =\lim _{h \rightarrow+\infty} \int_{\Omega \times \mathbb{R}} g_{h}(t) \varphi(x) \cdot d\left(D_{\widetilde{H}}^{\prime} \chi_{E_{u}}\right)(x, t) \\
& =-\lim _{h \rightarrow+\infty} \int_{\Omega \times \mathbb{R}} \chi_{E_{u}}(x, t) g_{h}(t) \operatorname{div}_{X} \varphi(x) d \mathscr{L}^{n+1}(x, t) \\
& =-\lim _{h \rightarrow+\infty} \int\left(\int_{\Omega}^{u(x)} g_{h}(t) d t\right) \operatorname{div}_{X} \varphi(x) d \mathscr{L}^{n}(x) .
\end{aligned}
$$

For every $z \in \mathbb{R}$ and every $h \in \mathbb{N}$ we have

$$
\int_{-\infty}^{z} g_{h}(t) d t \leq|z|+h+\frac{1}{2} \quad \text { and } \quad \lim _{h \rightarrow+\infty}\left(\int_{-\infty}^{z} g_{h}(t) d t-h-\frac{1}{2}\right)=z
$$

using the fact that $\int_{\Omega} \operatorname{div}_{X} \varphi(x) d \mathscr{L}^{n}(x)=0$, by the Dominated Convergence Theorem we obtain

$$
\begin{align*}
\int_{\Omega \times \mathbb{R}} \varphi(x) \cdot d\left(D_{\widetilde{H}}^{\prime} \chi_{E_{u}}\right)(x, t) & =-\lim _{h \rightarrow+\infty} \int_{\Omega}\left(\int_{-\infty}^{u(x)} g_{h}(t) d t-h-\frac{1}{2}\right) \operatorname{div}_{X} \varphi(x) d \mathscr{L}^{n}(x) \\
& =-\int_{\Omega} u(x) \operatorname{div}_{X} \varphi(x) d \mathscr{L}^{n}(x)  \tag{4.11}\\
& =\int_{\Omega} \varphi(x) \cdot d\left(D_{H} u\right)(x) .
\end{align*}
$$

In particular, $u \in B V_{H}(\Omega)$ and, for any open set $A \subset \Omega$,

$$
\begin{align*}
& \left|D_{H} u\right|(A) \leq\left|D_{\widetilde{H}}^{\prime} \chi_{E_{u}}\right|(A \times \mathbb{R})  \tag{4.12}\\
& \left|D_{X_{i}} u\right|(A) \leq\left|D_{\widetilde{X}_{i}} \chi_{E_{u}}\right|(A \times \mathbb{R}) \quad \text { for any } i=1, \ldots, m .
\end{align*}
$$

Before passing to the reverse implication we observe two facts. First, for any $\varphi \in C_{c}^{1}(\Omega)$ one has

$$
\begin{align*}
\int_{\Omega \times \mathbb{R}} \varphi(x) d\left(\partial_{t} \chi_{E_{u}}\right)(x, t) & =\lim _{h \rightarrow+\infty} \int_{\Omega \times \mathbb{R}} \varphi(x) g_{h}(t) d\left(\partial_{t} \chi_{E_{u}}\right)(x, t) \\
& =-\lim _{h \rightarrow+\infty} \int_{\Omega \times \mathbb{R}} \varphi(x) g_{h}^{\prime}(t) \chi_{E_{u}}(x, t) d \mathscr{L}^{n+1}(x, t) \\
& =-\lim _{h \rightarrow+\infty} \int_{\Omega} \varphi(x)\left(\int_{-\infty}^{u(x)} g_{h}^{\prime}(t) d t\right) d \mathscr{L}^{n}(x)  \tag{4.13}\\
& =-\lim _{h \rightarrow+\infty} \int_{\Omega} \varphi(x) g_{h}(u(x)) d \mathscr{L}^{n}(x) \\
& =-\int_{\Omega} \varphi d \mathscr{L}^{n}
\end{align*}
$$

whence, for any open set $A \subset \Omega$,

$$
\begin{equation*}
\mathscr{L}^{n}(A) \leq\left|\partial_{t} \chi_{E_{u}}\right|(A \times \mathbb{R}) \tag{4.14}
\end{equation*}
$$

Second, if $\varphi \in C_{c}^{1}\left(\Omega, \mathbb{R}^{m+1}\right)$ one has by (4.11) and (4.13)

$$
\int_{\Omega \times \mathbb{R}} \varphi(x) \cdot d\left(D_{\widetilde{H}} \chi_{E_{u}}\right)(x, t)=\int_{\Omega} \varphi(x) \cdot d\left(D_{H} u,-\mathscr{L}^{n}\right)(x)
$$

which gives for any open set $A \subset \Omega$

$$
\begin{equation*}
\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|(A) \leq\left|D_{\widetilde{H}} \chi_{E_{u}}\right|(A \times \mathbb{R}) \tag{4.15}
\end{equation*}
$$

Suppose now that $u \in B V_{H}(\Omega)$. Let $A \subset \Omega$ be open and let $\varphi \in C_{c}^{1}(A \times \mathbb{R})$ and $i=1, \ldots, m$ be fixed. Let $\left(u_{h}\right)_{h}$ be a sequence in $C^{\infty}(A) \cap B V_{H}(A)$ satisfying (4.10) (with $A$ in place of $\Omega$ ); then

$$
\begin{align*}
& \int_{A \times \mathbb{R}} \varphi d\left(D_{\widetilde{X}_{i}} \chi_{E_{u_{h}}}\right) \\
& =-\int_{A \times \mathbb{R}} \chi_{E_{u_{h}}}(x, t) \widetilde{X}_{i}^{*} \varphi(x, t) d \mathscr{L}^{n+1}(x, t) \\
& =-\int_{A}\left(\int_{-\infty}^{u_{h}(x)} \sum_{j=1}^{n} \partial_{x_{j}}\left(a_{i j}(x) \varphi(x, t)\right) d t\right) d \mathscr{L}^{n}(x) \tag{4.16}
\end{align*}
$$

$$
\begin{aligned}
& =-\int_{A}\left(\sum_{j=1}^{n} \partial_{x_{j}} \int_{-\infty}^{u_{h}(x)} a_{i j}(x) \varphi(x, t) d t-\sum_{j=1}^{n} a_{i j}(x) \varphi\left(x, u_{h}(x)\right) \partial_{x_{j}} u_{h}(x)\right) d \mathscr{L}^{n}(x) \\
& =\int_{A} \varphi\left(x, u_{h}(x)\right) X_{i} u_{h}(x) d \mathscr{L}^{n}(x)
\end{aligned}
$$

where we used the fact that $x \mapsto a_{i j}(x) \int_{-\infty}^{u_{h}(x)} \varphi(x, t) d t$ is in $C_{c}^{1}(A)$. In a similar way

$$
\begin{align*}
\int_{A \times \mathbb{R}} \varphi d\left(\partial_{t} \chi_{E_{u_{h}}}\right) & =-\int_{A}\left(\int_{-\infty}^{u_{h}(x)} \partial_{t} \varphi(x, t) d t\right) d \mathscr{L}^{n}(x)  \tag{4.17}\\
& =-\int_{A} \varphi\left(x, u_{h}(x)\right) d \mathscr{L}^{n}(x)
\end{align*}
$$

Formulas (4.16) and (4.17) imply that for any $\varphi \in C_{c}^{1}\left(A \times \mathbb{R}, \mathbb{R}^{m+1}\right)$

$$
\int_{A \times \mathbb{R}} \varphi \cdot d\left(D_{\widetilde{H}} \chi_{E_{u_{h}}}\right)=\int_{A} \varphi\left(x, u_{h}(x)\right) \cdot d\left(D_{H} u_{h},-\mathscr{L}^{n}\right)(x) .
$$

Since $\chi_{E_{u_{h}}} \rightarrow \chi_{E_{u}}$ in $L^{1}(A \times \mathbb{R})$ we obtain

$$
\begin{align*}
\left|D_{\widetilde{H}} \chi_{E_{u}}\right|(A \times \mathbb{R}) & \leq \liminf _{h \rightarrow+\infty}\left|D_{\widetilde{H}} \chi_{E_{u_{h}}}\right|(A \times \mathbb{R}) \leq \lim _{h \rightarrow+\infty}\left|\left(D_{H} u_{h},-\mathscr{L}^{n}\right)\right|(A)  \tag{4.18}\\
& =\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|(A)<+\infty
\end{align*}
$$

which proves that $\chi_{E_{u}} \in B V_{\widetilde{H}}(\Omega \times \mathbb{R})$, as desired. Notice that, using the lower semicontinuity in a similar way, one also gets

$$
\begin{align*}
& \left|D_{\widetilde{H}}^{\prime} \chi_{E_{u}}\right|(A \times \mathbb{R}) \leq\left|D_{H} u\right|(A) \\
& \left|D_{\widetilde{X}_{i}} \chi_{E_{u}}\right|(A \times \mathbb{R}) \leq\left|D_{X_{i}} u\right|(A) \text { for any } i=1, \ldots, m  \tag{4.19}\\
& \left|\partial_{t} \chi_{E_{u}}\right|(A \times \mathbb{R}) \leq \mathscr{L}^{n}(A)<+\infty
\end{align*}
$$

Eventually, statements (i) and (ii) follow from (4.11) and (4.13), while statements (iii)-(vi) are consequences of formulas (4.12), (4.14), (4.15), (4.18) and (4.19).

Let us introduce some further notation. For $u \in B V_{H, l o c}(\Omega)$ we decompose its distributional horizontal derivatives as $D_{H} u=D_{H}^{a} u+D_{H}^{s} u$, where $D_{H}^{a} u$ is absolutely continuous with respect to $\mathscr{L}^{n}$ and $D_{H}^{s} u$ is singular with respect to $\mathscr{L}^{n}$. We also write $D_{H}^{a} u=X u \mathscr{L}^{n}$ for some function $X u \in L_{l o c}^{1}\left(\Omega, \mathbb{R}^{m}\right)$.

We also consider the polar decomposition $D_{H} u=\sigma_{u}\left|D_{H} u\right|$, where $\sigma_{u}: \Omega \rightarrow \mathbb{S}^{m-1}$ is a $\left|D_{H} u\right|$-measurable function. In case $u=\chi_{E}$ is the characteristic function of a set
$E \subset \Omega \times \mathbb{R}$ of locally finite $\widetilde{H}$-perimeter in $\Omega \times \mathbb{R}$ we write $D_{\widetilde{H}} \chi_{E}=\nu_{E}\left|D_{\widetilde{H}} \chi_{E}\right|$ for some Borel function $\nu_{E}=\left(\left(\nu_{E}\right)_{1}, \ldots,\left(\nu_{E}\right)_{m+1}\right)$ called horizontal inner normal to $E$.

The following result is basically a consequence of Theorem 4.2.
Theorem 4.3. Let $u \in B V_{H}(\Omega)$ and define

$$
\begin{aligned}
& S:=\left\{(x, t) \in \Omega \times \mathbb{R}:\left(\nu_{E_{u}}\right)_{m+1}(x, t)=0\right\} \\
& T:=\left\{(x, t) \in \Omega \times \mathbb{R}:\left(\nu_{E_{u}}\right)_{m+1}(x, t) \neq 0\right\} .
\end{aligned}
$$

Then, the following identities hold

$$
\begin{align*}
& \nu_{E_{u}}(x, t)=\left(\sigma_{u}(x), 0\right) \quad \text { for }\left|D_{\widetilde{H}} \chi_{E_{u}}\right| \text { a.e. }(x, t) \in S ;  \tag{4.20}\\
& \nu_{E_{u}}(x, t)=\frac{(X u(x),-1)}{\sqrt{1+|X u(x)|^{2}}} \text { for }\left|D_{\widetilde{H}} \chi_{E_{u}}\right| \text {-a.e. }(x, t) \in T ;  \tag{4.21}\\
& \pi_{\#}\left(D_{\widetilde{H}} \chi_{E_{u}}\llcorner S)=\left(D_{H}^{s} u, 0\right)\right.  \tag{4.22}\\
& \pi_{\#}\left(D_{\widetilde{H}} \chi_{E_{u}}\llcorner T)=\left(D_{H}^{a} u,-\mathscr{L}^{n}\right) .\right. \tag{4.23}
\end{align*}
$$

Proof. Thanks to Theorem 4.2 (vi) we can disintegrate the measure $\left|D_{\widetilde{H}} \chi_{E_{u}}\right|$ with respect to $\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|$ (see e.g. [3, Theorem 2.28]): for every $x \in \Omega$ there exists a probability measure $\mu_{x}$ on $\mathbb{R}$ such that for every Borel function $g \in L^{1}\left(\Omega \times \mathbb{R},\left|D_{\widetilde{H}} \chi_{E_{u}}\right|\right)$

$$
\int_{\Omega \times \mathbb{R}} g(x, t) d\left|D_{\widetilde{H}} \chi_{E_{u}}\right|(x, t)=\int_{\Omega}\left(\int_{\mathbb{R}} g(x, t) d \mu_{x}(t)\right) d\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|(x) .
$$

It follows that for any Borel function $\varphi: \Omega \rightarrow \mathbb{R}$

$$
\begin{align*}
\int_{\Omega} \varphi(x) d\left(D_{H} u,-\mathscr{L}^{n}\right)(x) & =\int_{\Omega} \varphi(x) d \pi_{\#}\left(\nu_{E_{u}}\left|D_{\widetilde{H}} \chi_{E_{u}}\right|\right)(x) \\
& =\int_{\Omega \times \mathbb{R}} \varphi(x) \nu_{E_{u}}(x, t) d\left|D_{\widetilde{H}} \chi_{E_{u}}\right|(x, t)  \tag{4.24}\\
& =\int_{\Omega} \varphi(x)\left(\int_{\mathbb{R}} \nu_{E_{u}}(x, t) d \mu_{x}(u)\right) d\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|(x) .
\end{align*}
$$

Since $D_{H}^{a} u$ and $D_{H}^{s} u$ are mutually singular we have

$$
\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|=\left|\left(D_{H}^{a} u,-\mathscr{L}^{n}\right)\right|+\left|\left(D_{H}^{s} u, 0\right)\right|=\sqrt{1+|X u|^{2}} \mathscr{L}^{n}+\left|D_{H}^{s} u\right|
$$

and (4.24) gives

$$
\begin{align*}
& \int_{\Omega} \varphi d\left((X u,-1) \mathscr{L}^{n}+\left(\sigma_{u}, 0\right)\left|D_{H}^{s} u\right|\right)  \tag{4.25}\\
= & \int_{\Omega} \varphi(x)\left(\int_{\mathbb{R}} \nu_{E_{u}}(x, t) d \mu_{x}(t)\right) d\left(\sqrt{1+|X u|^{2}} \mathscr{L}^{n}+\left|D_{H}^{s} u\right|\right)(x) . \tag{4.26}
\end{align*}
$$

Denote by $I$ a subset of $\Omega$ such that $\mathscr{L}^{n}(I)=0$ and $\left|D_{H}^{s} u\right|(\Omega \backslash I)=0$. Considering Borel test functions $\varphi$ such that $\varphi=0$ in $\Omega \backslash I$, we deduce that for $\left|D_{H}^{s} u\right|$-a.e. $x \in I$ one has

$$
\left(\sigma_{u}(x), 0\right)=\int_{\mathbb{R}} \nu_{E_{u}}(x, t) d \mu_{x}(t)
$$

Taking on both sides the scalar product with $\left(\sigma_{u}(x), 0\right)$ we get

$$
\left\langle\left(\sigma_{u}(x), 0\right), \int_{\mathbb{R}} \nu_{E_{u}}(x, t) d \mu_{x}(t)\right\rangle=1
$$

and, since $\mu_{x}(\mathbb{R})=1$ and (for $\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|$-a.e. $\left.x \in \Omega\right)\left|\nu_{E_{u}}(x, t)\right|=1$ for $\mu_{x}$-a.e. $t$, we deduce that

$$
\nu_{E_{u}}(x, t)=\left(\sigma_{u}(x), 0\right) \quad \text { for }\left|D_{H}^{s} u\right| \text {-a.e. } x \in I \text { and } \mu_{x} \text {-a.e. } t \in \mathbb{R},
$$

i.e.,

$$
\begin{equation*}
\nu_{E_{u}}(x, t)=\left(\sigma_{u}(x), 0\right) \quad \text { for }\left|D_{\widetilde{H}} \chi_{E_{u}}\right| \text {-a.e. }(x, t) \in I \times \mathbb{R} . \tag{4.27}
\end{equation*}
$$

Taking into account again (4.25) and letting $\varphi$ be such that $\varphi=0$ on $I$ we instead obtain

$$
\begin{aligned}
& \int_{\Omega} \varphi \frac{(X u,-1)}{\sqrt{1+|X u|^{2}}} \sqrt{1+|X u|^{2}} d \mathscr{L}^{n} \\
= & \int_{\Omega} \varphi(x)\left(\int_{\mathbb{R}} \nu_{E_{u}}(x, t) d \mu_{x}(t)\right) \sqrt{1+|X u(x)|^{2}} d \mathscr{L}^{n}(x)
\end{aligned}
$$

Consequently, for $\mathscr{L}^{n}$-a.e. $x \in \Omega \backslash I$ we have

$$
\int_{\mathbb{R}} \nu_{E_{u}}(x, t) d \mu_{x}(t)=\frac{(X u(x),-1)}{\sqrt{1+|X u(x)|^{2}}}
$$

Reasoning as before we deduce that

$$
\nu_{E_{u}}(x, t)=\frac{(X u(x),-1)}{\sqrt{1+|X u(x)|^{2}}} \quad \text { for } \mathscr{L}^{n} \text {-a.e. } x \in \Omega \backslash I \text { and } \mu_{x} \text {-a.e. } t \in \mathbb{R},
$$

or equivalently

$$
\begin{equation*}
\nu_{E_{u}}(x, t)=\frac{(X u(x),-1)}{\sqrt{1+|X u(x)|^{2}}} \quad \text { for }\left|D_{\widetilde{H}} \chi_{E_{u}}\right| \text {-a.e. }(x, t) \in(\Omega \backslash I) \times \mathbb{R} \tag{4.28}
\end{equation*}
$$

Formula (4.27) implies that $\left|D_{\widetilde{H}} \chi_{E_{u}}\right|$-a.e. $(x, t) \in I \times \mathbb{R}$ belongs to $S$ and that $\left|D_{\widetilde{H}} \chi_{E_{u}}\right|$-a.e. $(x, t) \in T$ belongs to $(\Omega \backslash I) \times \mathbb{R}$. Similarly, (4.28) says that $\left|D_{\widetilde{H}} \chi_{E_{u}}\right|$-a.e. $(x, t) \in(\Omega \backslash I) \times \mathbb{R}$ belongs to $T$ and that $\left|D_{\widetilde{H}} \chi_{E_{u}}\right|$-a.e. $(x, t) \in S$ belongs to $I \times \mathbb{R}$. Since $S$ and $T$ are disjoint, this is enough to conclude (4.20) and (4.21). Statement (4.22) now easily follows because

$$
\pi_{\#}\left(D_{\widetilde{H}} \chi_{E_{u}}\llcorner S)=\pi_{\#}\left(\nu_{E_{u}}\left|D_{\widetilde{H}} \chi_{E_{u}}\right|\llcorner(I \times \mathbb{R}))=\left(\sigma_{u}, 0\right)\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|\left\llcorner I=\left(D_{H}^{s} u, 0\right)\right.\right.\right.
$$

Similarly, one has

$$
\begin{aligned}
\pi_{\#}\left(D_{\widetilde{H}} \chi_{E_{u}}\llcorner T)\right. & =\pi_{\#}\left(\nu_{E_{u}}\left|D_{\widetilde{H}} \chi_{E_{u}}\right|\llcorner((\Omega \backslash I) \times \mathbb{R}))\right. \\
& =\frac{(X u,-1)}{\sqrt{1+|X u|^{2}}}\left|\left(D_{H} u,-\mathscr{L}^{n}\right)\right|\left\llcorner(\Omega \backslash I)=(X u,-1) \mathscr{L}^{n},\right.
\end{aligned}
$$

which gives (4.23).

## 5. The rank-one theorem for $\boldsymbol{B} V_{H}$ functions in Carnot groups

We now use the results of the previous section in the setting of a Carnot group $\mathbb{G}$. We utilize the notation of Section 2 ; in particular, we identify $\mathbb{G} \equiv \mathbb{R}^{n}$ by exponential coordinates and a left-invariant basis $X_{1}, \ldots, X_{m}$ of $\mathfrak{g}_{1}$ is fixed. The vector fields $\widetilde{X}_{1}, \ldots, \widetilde{X}_{m+1}$ on $\mathbb{G} \times \mathbb{R}$ are defined as in the previous section; notice that they form a basis of the first layer of the Lie algebra of $\mathbb{G} \times \mathbb{R}$. The homogeneous dimension of $\mathbb{G} \times \mathbb{R}$ is $Q+1$.

A set $R \subset \mathbb{G}$ is $H$-rectifiable if $\mathcal{H}^{Q-1}(R)<\infty$ and there exists a (finite or countable) family $\left(\Sigma_{i}\right)_{i}$ of $C_{H}^{1}$ hypersurfaces in $\mathbb{G}$ such that

$$
\mathcal{H}^{Q-1}\left(R \backslash \bigcup_{i} \Sigma_{i}\right)=0
$$

We define the horizontal normal $\nu_{R}$ to $R$ as

$$
\nu_{R}(x):=\nu_{\Sigma_{i}}(x) \quad \text { if } x \in R \cap \Sigma_{i} \backslash \cup_{j<i} \Sigma_{j}
$$

The normal $\nu_{R}$ is well-defined (up to sign) $\mathcal{H}^{Q-1}$-a.e. on $R .{ }^{2}$

[^3]Definition 5.1. We say that a Carnot group $\mathbb{G}$ satisfies property $\mathscr{R}$ if the following holds. For any bounded open set $\Omega \subset \mathbb{G}$ and any $u \in B V_{H}(\Omega)$, the distributional $\widetilde{X}$-derivatives $D_{\widetilde{H}} \chi_{E_{u}}$ of the characteristic function of the subgraph $E_{u}$ of $u$ can be represented as

$$
\begin{equation*}
D_{\widetilde{H}} \chi_{E_{u}}=\nu_{\partial_{H}^{*} E_{u}} \theta \mathcal{S}^{Q}\left\llcorner\partial_{H}^{*} E_{u}\right. \tag{5.29}
\end{equation*}
$$

for some $H$-rectifiable set $\partial_{H}^{*} E_{u}$ in $\Omega \times \mathbb{R}$ and some positive density $\theta \in L^{1}\left(\partial_{H}^{*} E_{u}, \mathcal{S}^{Q}\right)$. We call $\partial_{H}^{*} E_{u}$ the $H$-reduced boundary of $E_{u}$.

Notice that, in Definition 5.1, the measure $D_{\widetilde{H}} \chi_{E_{u}}$ has finite total variation by Theorem 4.2.

Remark 5.2. In view of Theorem 1.3, for the validity of property $\mathscr{R}$ in $\mathbb{G}$ it is enough that a rectifiability theorem holds for sets with finite $H$-perimeter in $\mathbb{G} \times \mathbb{R}$; namely, it suffices that any set $E$ with finite $H$-perimeter in $\mathbb{G} \times \mathbb{R}$ satisfies $D_{\widetilde{H}} \chi_{E}=\nu_{\partial_{H}^{*} E} \theta \mathcal{S}^{Q}\left\llcorner\partial_{H}^{*} E\right.$ for some $H$-rectifiable set $\partial_{H}^{*} E$ and some positive density $\theta \in L^{1}\left(\partial_{H}^{*} E, \mathcal{S}^{Q}\right)$. We conjecture that this, in turn, is equivalent to the validity of a rectifiability theorem for sets with finite $H$-perimeter in $\mathbb{G}$; in particular, we conjecture that property $\mathscr{R}$ is equivalent to the rectifiability theorem in $\mathbb{G}$.

Remark 5.3. If $\mathbb{G}$ is a Carnot group of step 2 , then $\mathbb{G}$ satisfies property $\mathscr{R}$ : this follows from the fact that $\mathbb{G} \times \mathbb{R}$ is also a step 2 Carnot group and that the rectifiability theorem holds in any step 2 Carnot group, see [15].

Remark 5.4. If (5.29) holds, then

$$
\left|D_{\widetilde{H}} \chi_{E_{u}}\right|=\theta \mathcal{S}^{Q}\left\llcorner\partial_{H}^{*} E_{u} \quad \text { and } \quad \nu_{E_{u}}=\nu_{\partial_{H}^{*} E_{u}} \mathcal{S}^{Q} \text {-a.e. on } \partial_{H}^{*} E_{u} .\right.
$$

Proof of Theorem 1.1. Without loss of generality one can assume that $u=\left(u_{1}, \ldots, u_{d}\right) \in$ $B V_{H}\left(\Omega, \mathbb{R}^{d}\right)$. It is not restrictive to assume that $\Omega$ is bounded. For any $i=1, \ldots, d$ we write $D_{H}^{s} u_{i}=\sigma_{i}\left|D_{H}^{s} u_{i}\right|$ for a $\left|D_{H}^{s} u_{i}\right|$-measurable map $\sigma_{i}: \Omega \rightarrow \mathbb{S}^{m-1}$; notice that, using the notation of Section 4, the equality $\sigma_{i}=\sigma_{u_{i}}$ holds $\left|D^{s} u_{i}\right|$-almost everywhere. We also let $E_{i}:=\left\{(x, t) \in \Omega \times \mathbb{R}: t<u_{i}(x)\right\}$ be the subgraph of $u_{i}$, that has finite $H$-perimeter in $\Omega \times \mathbb{R}$ by Theorem 4.2. Denoting by $\partial_{H}^{*} E_{i}$ the $H$-reduced boundary of $E_{i}$ and writing $\nu_{i}=\nu_{E_{i}}$ for the measure theoretic inner normal to $E_{i}$, we have by Theorem 4.3 and Remark 5.4 that

$$
\left|D_{H}^{s} u_{i}\right|=\pi_{\#}\left(\theta_{i} \mathcal{S}^{Q}\left\llcorner S_{i}\right) \quad \text { for some positive } \theta_{i} \in L^{1}\left(\partial_{H}^{*} E_{i}, \mathcal{S}^{Q}\right)\right.
$$

[^4]where $S_{i}:=\left\{p \in \partial_{H}^{*} E_{i}:\left(\nu_{i}(p)\right)_{m+1}=0\right\}$ and $\pi_{\#}$ denotes push-forward of measures through the projection $\pi$ defined by $\mathbb{G} \times \mathbb{R} \ni(x, t) \mapsto x \in \mathbb{G}$. By rectifiability, we can assume that $\partial_{H}^{*} E_{i}$ is contained in the union $\cup_{\ell \in \mathbb{N}} \Sigma_{\ell}^{i}$ of $C_{H}^{1}$ hypersurfaces $\Sigma_{\ell}^{i}$ in $\mathbb{G} \times \mathbb{R}$.

Using Theorem 4.3, Remark 5.4 and Lemma 3.2 the following properties hold for $\mathcal{S}^{Q}$-a.e. $p \in S_{1} \cup \cdots \cup S_{d}$ :

$$
\begin{align*}
& \text { if } p \in S_{i} \text {, then } \nu_{i}(p)=\left(\sigma_{i}(\pi(p)), 0\right)  \tag{5.30}\\
& \text { if } p \in \Sigma_{\ell}^{i} \text {, then } \nu_{i}(p)= \pm \nu_{\Sigma_{\ell}^{i}}(p)  \tag{5.31}\\
& \text { if } p \in \Sigma_{\ell}^{i} \text { and } \exists q \in S_{j} \cap \Sigma_{k}^{j} \cap \pi^{-1}(\pi(p)) \text {, then } \nu_{\Sigma_{\ell}^{i}}(p)= \pm \nu_{\Sigma_{k}^{j}}(q) \text {. } \tag{5.32}
\end{align*}
$$

Up to modifying each $S_{i}$ on a $\mathcal{S}^{Q}$-negligible set and each $\sigma_{i}$ on a $\left|D_{H}^{s} u_{i}\right|$-negligible set, we can assume that (5.30), (5.31) and (5.32) hold for any $p \in S_{1} \cup \cdots \cup S_{d}$ and that, for any $i=1, \ldots, d, \sigma_{i}=0$ on $\Omega \backslash \pi\left(S_{i}\right)$.

Since $D_{H}^{s} u=\left(\sigma_{1}\left|D_{H}^{s} u_{1}\right|, \ldots, \sigma_{d}\left|D_{H}^{s} u_{d}\right|\right)$ and $\left|D_{H}^{s} u\right|$ is concentrated on $\pi\left(S_{1}\right) \cup \cdots \cup$ $\pi\left(S_{d}\right)$, it is enough to prove that the matrix-valued function $\left(\sigma_{1}, \ldots, \sigma_{d}\right)$ has rank 1 on $\pi\left(S_{1}\right) \cup \cdots \cup \pi\left(S_{d}\right)$. This follows if we prove that the implication

$$
i, j \in\{1, \ldots, d\}, i \neq j, x \in \pi\left(S_{i}\right) \Longrightarrow \sigma_{j}(x) \in\left\{0, \sigma_{i}(x),-\sigma_{i}(x)\right\}
$$

holds. If $i, j, x$ are as above and $x \notin \pi\left(S_{j}\right)$, then $\sigma_{j}(x)=0$. Otherwise, $x \in \pi\left(S_{i}\right) \cap \pi\left(S_{j}\right)$, i.e., there exist $p \in S_{i}$ and $\ell \in \mathbb{N}$ such that $\pi(p)=x$ and $\sigma_{i}(x)= \pm \nu_{\Sigma_{\ell}^{i}}(p)$ and there exist $q \in S_{j}$ and $k \in \mathbb{N}$ such that $\pi(q)=x$ and $\sigma_{j}(x)= \pm \nu_{\Sigma_{k}^{j}}(p)$. By (5.32) we obtain $\sigma_{j}(x)= \pm \sigma_{i}(x)$, as wished.

Remark 5.5. As an easy consequence of Remark 2.4 and Remark 5.3, Theorem 1.1 holds for the Heisenberg group $\mathbb{H}^{n}$ provided $n \geq 2$. This result does not directly follow from [9], as we now briefly explain using the notation of Example 2.1 and restricting for simplicity to $n=2$, the general case $n \geq 2$ being a straightforward generalization.

Let $u \in B V_{H}\left(\Omega, \mathbb{R}^{m}\right)$ for some open set $\Omega \subset \mathbb{H}^{2}$. It can be easily seen that the matrixvalued measure $\left(\mu_{1}, \mu_{2}, \mu_{3}, \mu_{4}\right):=D_{H} u=\left(X_{1} u, X_{2} u, Y_{1} u, Y_{2} u\right)$ satisfies the equations

$$
\mathscr{A} \mu:=\left(\begin{array}{l}
X_{1} \mu_{2}-X_{2} \mu_{1} \\
Y_{1} \mu_{4}-Y_{2} \mu_{3} \\
X_{1} \mu_{4}-Y_{2} \mu_{1} \\
Y_{1} \mu_{2}-X_{2} \mu_{3} \\
X_{1} \mu_{3}-Y_{1} \mu_{1}+Y_{2} \mu_{2}-X_{2} \mu_{4}
\end{array}\right)=0
$$

in the sense of distributions. Write the first-order differential operator $\mathscr{A}$ (the horizontal curl in $\mathbb{H}^{2}$, see [5, Example 3.12]) in the form

$$
\mathscr{A}=A_{1} \partial_{x_{1}}+A_{2} \partial_{x_{2}}+A_{3} \partial_{y_{1}}+A_{4} \partial_{y_{2}}+A_{5} \partial_{t}
$$

for suitable $A_{j}=A_{j}(x, y, t)$ and consider the wave cone $\Lambda_{\mathscr{A}}(x, y, t)$ (see [9]) associated with $\mathscr{A}$

$$
\Lambda_{\mathscr{A}}(x, y, t):=\bigcup_{\xi \in \mathbb{R}^{5} \backslash\{0\}} \operatorname{ker} \mathbb{A}_{x, y, t}(\xi), \quad \text { where } \mathbb{A}_{x, y, t}(\xi):=2 \pi i \sum_{j=1}^{5} A_{j}(x, y, t) \xi_{j}
$$

One can readily check that

$$
\mathbb{A}_{x, y, t}(\xi)=0 \quad \text { for } \xi:=\left(\frac{y}{2},-\frac{x}{2}, 1\right) \in \mathbb{R}^{5} \backslash\{0\}
$$

i.e., the wave cone $\Lambda_{\mathscr{A}}(x, y, t)$ is the full space for any $(x, y, t) \in \mathbb{H}^{2}$. In particular, [9, Theorem 1.1] gives no information on the polar decomposition of $D_{H}^{s} u$.

Remark 5.6. The rank-one property for $B V$ functions in the first Heisenberg group remains a very interesting open question, since it does not follow either from Theorem 1.1 (because property $\mathscr{C}_{2}$ fails for $\mathbb{H}^{1}$ ) or from [9, Theorem 1.1], as we now explain.

Let $u \in B V_{H}\left(\Omega, \mathbb{R}^{m}\right)$ for some open set $\Omega \subset \mathbb{H}^{1}$; we use again the notation of Example 2.1 and we set $p=(x, y, t) \in \mathbb{H}^{1} \equiv \mathbb{R}^{3}$. One can check that $\left(\mu_{1}, \mu_{2}\right):=D_{H} u=$ ( $X u, Y u$ ) satisfies

$$
\mathscr{A} \mu:=\binom{Y X \mu_{1}-2 X Y \mu_{1}+X X \mu_{2}}{Y Y \mu_{1}-2 Y X \mu_{2}+X Y \mu_{2},}=0
$$

in the sense of distributions. Now $\mathscr{A}$ (the horizontal curl in $\mathbb{H}^{1}$, see [5, Example 3.11]) is a second-order differential operator that one can write as

$$
\mathscr{A}=\sum_{|\alpha|=2} A_{\alpha}(p) \partial^{\alpha}
$$

where $\alpha \in \mathbb{N}^{3}$ is a multi-index and $\partial^{\alpha}=\partial_{x}^{\alpha_{1}} \partial_{y}^{\alpha_{2}} \partial_{t}^{\alpha_{3}}$. As before, one can define the wave cone

$$
\Lambda_{\mathscr{A}}(p)=\bigcup_{\xi \in \mathbb{R}^{3} \backslash\{0\}} \operatorname{ker} \mathbb{A}_{p}(\xi), \quad \text { where } \mathbb{A}_{p}(\xi)=(2 \pi i)^{2} \sum_{|\alpha|=2} A_{\alpha}(p) \xi^{\alpha}
$$

Again, one has

$$
\mathbb{A}_{p}(\xi)=0 \quad \text { for } \xi:=\left(\frac{y}{2},-\frac{x}{2}, 1\right) \in \mathbb{R}^{3} \backslash\{0\}
$$

and the wave cone $\Lambda_{\mathscr{A}}(x, y, t)$ is the full space.

## Appendix A. Intersection of regular hypersurfaces vs. intrinsic Lipschitz graphs

## A.1. Intrinsic Lipschitz graphs

We follow [12]. Let $\mathbb{W}, \mathbb{H}$ be homogeneous (i.e., invariant under dilations) complementary subgroups of $\mathbb{G}$, i.e., such that $\mathbb{W} \cap \mathbb{H}=\{0\}$ and $\mathbb{G}=\mathbb{W} \mathbb{H}$. In particular, for any $x \in \mathbb{G}$ there exist unique $x_{\mathbb{W}} \in \mathbb{W}$ and $x_{\mathbb{H}} \in \mathbb{H}$ such that $x=x_{\mathbb{W}} x_{\mathbb{H}}$. Recall (see e.g. [12, Remark 2.3]) that any homogeneous subgroup $\mathbb{W}$ is stratified, that is, its Lie algebra $\mathfrak{w}$ is a subalgebra of $\mathfrak{g}$ and $\mathfrak{w}=\mathfrak{w}_{1} \oplus \cdots \oplus \mathfrak{w}_{s}$ where $\mathfrak{w}_{i}=\mathfrak{w} \cap \mathfrak{g}_{i}$. Moreover, the metric (Hausdorff) dimension of $\mathbb{W}$ is $Q_{\mathbb{W}}:=\sum_{i=1}^{s} i \operatorname{dim} \mathfrak{w}_{i}$.

The intrinsic graph of a function $\phi: \mathbb{W} \rightarrow \mathbb{H}$ is defined by

$$
\operatorname{gr} \phi:=\{w \phi(w): w \in \mathbb{W}\}
$$

We introduce the homogeneous cones $C_{\mathbb{W}, \mathbb{H}}(x, \alpha)$ of center $x \in \mathbb{G}$ and aperture $\alpha>0$ as

$$
C_{\mathbb{W}, \mathbb{H}}(x, \alpha):=x C_{\mathbb{W}, \mathbb{H}}(0, \alpha) \quad \text { where } \quad C_{\mathbb{W}, \mathbb{H}}(0, \alpha):=\left\{y \in \mathbb{G}:\left\|x_{\mathbb{W}}\right\| \leq \alpha\left\|x_{\mathbb{H}}\right\|\right\} .
$$

Definition A.1. A function $\phi: \mathbb{W} \rightarrow \mathbb{H}$ is intrinsic Lipschitz if there exists $\alpha>0$ such that

$$
\forall x \in \operatorname{gr} \phi \quad \operatorname{gr} \phi \cap C_{\mathbb{W}, \mathbb{H}}(x, \alpha)=\{x\} .
$$

We say that $S \subset \mathbb{G}$ is an intrinsic Lipschitz graph if there exists an intrinsic Lipschitz $\operatorname{map} \phi: \mathbb{W} \rightarrow \mathbb{H}$ such that $S=\operatorname{gr} \phi$.

Remark A.2. We will later use the following equivalent definition of intrinsic Lipschitz continuity: $\phi: \mathbb{W} \rightarrow \mathbb{H}$ is intrinsic Lipschitz if and only if there exists $\beta>0$ such that

$$
\forall x \in \operatorname{gr} \phi \quad \operatorname{gr} \phi \cap D(x, \mathbb{H}, \beta)=\{x\}
$$

where the homogeneous cone $D(x, \mathbb{H}, \beta)$ is defined by

$$
D(x, \mathbb{H}, \beta):=x D(\mathbb{H}, \beta) \quad \text { and } \quad D(\mathbb{H}, \beta):=\bigcup_{h \in \mathbb{H}} \overline{B(h, \beta d(h, 0))} .
$$

Indeed, it is enough to observe that, for any $\alpha>0$ and $\beta>0$, there exist $\beta_{\alpha}>0$ and $\alpha_{\beta}>0$ such that

$$
C_{\mathbb{W}, \mathbb{H}}(x, \alpha) \supset D\left(\mathbb{H}, \beta_{\alpha}\right) \quad \text { and } \quad D(\mathbb{H}, \beta) \supset C_{\mathbb{W}, \mathbb{H}}\left(x, \alpha_{\beta}\right) .
$$

This, in turn, is a consequence of a homogeneity argument based on the following fact: if $S:=\{x \in \mathbb{G}:\|x\|=1\}$ and

$$
A_{\alpha}:=S \cap \operatorname{int}\left(C_{\mathbb{W}, \mathbb{H}}(x, \alpha)\right), \quad B_{\beta}:=S \cap \operatorname{int}(D(\mathbb{H}, \beta)),
$$

then $\left\{A_{\alpha}\right\}_{\alpha>0}$ and $\left\{B_{\beta}\right\}_{\beta>0}$ are monotone families of (relatively) open subsets of $S$ such that the intersection

$$
\bigcap_{\alpha>0} A_{\alpha}=\bigcap_{\beta>0} B_{\beta}=\mathbb{H} \cap S
$$

is a compact set.
The following result will be used in the proof of Theorem 1.4.
Theorem A. 3 ([12, Theorem 3.9]). Let $\mathbb{W}, \mathbb{H}$ be homogeneous complementary subgroups of $\mathbb{G}$, let $\phi: \mathbb{W} \rightarrow \mathbb{H}$ be intrinsic Lipschitz and let $\alpha>0$ be as in Definition A.1. Then there exists a positive $C=C(\mathbb{W}, \mathbb{H}, \alpha)$ such that

$$
\frac{1}{C} r^{Q_{\mathrm{W}}} \leq \mathcal{H}^{Q_{\mathrm{W}}}(\operatorname{gr} \phi \cap B(x, r)) \leq C r^{Q_{\mathrm{W}}} \quad \forall x \in \operatorname{gr} \phi, r>0
$$

## A.2. Transversal intersections of $C_{H}^{1}$ hypersurfaces are intrinsic Lipschitz graphs

The aim of this section is proving Theorem A.5, due to V. Magnani [22], for which we need the preparatory Lemma A.4. Actually, its use could be avoided by utilizing a local version of Theorem A. 3 which, even though not explicitly stated there, would easily follow adapting the techniques of [12]. We note however that Lemma A.4, and (A.33) in particular, provides also a proof of (3.2).

Lemma A.4. Let $\Omega \subset \mathbb{G}$ be open, $f \in C_{H}^{1}(\Omega), \bar{x} \in \Omega$ and let $A:=\nabla_{H} f(\bar{x})$. Then, for any $\varepsilon>0$ there exist an open set $U \subset \Omega$ with $\bar{x} \in U$ and a function $g \in C_{H}^{1}(\mathbb{G})$ such that
(i) $g=f$ on $U$;
(ii) $\left|\nabla_{H} g-A\right|<\varepsilon$ on $\mathbb{G}$.

Proof. Without loss of generality we can assume that $\bar{x}=0$. We preliminarily fix a smooth function $\chi: \mathbb{G} \rightarrow[0,1]$ such that $\chi \equiv 1$ on $B_{1}$ and $\chi \equiv 0$ on $\mathbb{G} \backslash B_{2}$. For any $r>0$, the functions $\chi_{r}:=\chi \circ \delta_{1 / r}$ satisfy

$$
0 \leq \chi_{r} \leq 1, \quad \chi \equiv 1 \text { on } B_{r}, \quad \chi \equiv 0 \text { on } \mathbb{G} \backslash B_{2 r}, \quad\left|\nabla_{H} \chi_{r}\right| \leq \frac{C}{r}
$$

for some positive $C$ independent of $r$.
Let $\varepsilon>0$ be fixed. We fix $r>0$ such that $\left|\nabla_{H} f-A\right|<\varepsilon$ on $B_{2 r}$. With this choice, setting $\lambda(x):=A_{1} x_{1}+\cdots+A_{m} x_{m}$ (where $x$ is represented in exponential coordinates) we prove that

$$
\begin{equation*}
|f(x)-\lambda(x)|<2 \varepsilon r \quad \text { for any } x \in B_{2 r} \tag{А.33}
\end{equation*}
$$

Indeed, for any $x \in B_{2 r}$ there exists a horizontal curve $\gamma:[0,1] \rightarrow \mathbb{G}$ such that $\gamma(0)=0$, $\gamma(1)=x$ and $L(\gamma)<2 r$. By definition, there exists $h \in L^{\infty}\left([0,1], \mathbb{R}^{m}\right)$ such that

$$
\dot{\gamma}(t)=\sum_{i=1}^{m} h_{i}(t) X_{i}(\gamma(t)) \quad \text { for a.e. } t \in[0,1] .
$$

Moreover, for any $i=1, \ldots, m$ we have $\int_{0}^{1} h_{i}=x_{i}$, because in exponential coordinates one has $X_{i}(x)=\partial_{x_{i}}+\sum_{\ell>m+1} a_{i \ell} \partial_{x_{\ell}}$ (see e.g. [29]). It follows that

$$
\begin{aligned}
|f(x)-\lambda(x)| & =\left|\int_{0}^{1} \sum_{i=1}^{m} h_{i}(t) X_{i} f(\gamma(t)) d t-\int_{0}^{1} \sum_{i=1}^{m} A_{i} h_{i}(t) d t\right| \\
& \leq \int_{0}^{1}|h(t)|\left\|\nabla_{H} f(\gamma(t))-A\right\| d t \\
& <2 \varepsilon r .
\end{aligned}
$$

We now define $g:=\chi_{r} f+\left(1-\chi_{r}\right) \lambda$; statement (i) is readily checked, while for (ii)

$$
\begin{aligned}
\left|\nabla_{H} g-A\right| & =\left|\chi_{r} \nabla_{H} f+\left(1-\chi_{r}\right) A+(f-\lambda) \nabla_{H} \chi_{r}-A\right| \\
& \leq \chi_{r}\left|\nabla_{H} f-A\right|+|f-\lambda|\left|\nabla_{H} \chi_{r}\right| \\
& \leq \varepsilon+2 C \varepsilon .
\end{aligned}
$$

The proof is then accomplished.
We can now prove the main result of this section. Since property $\mathscr{C}_{1}$ holds in any Carnot group, when $k=1$ Theorem A. 5 states in particular that hypersurfaces of class $C_{H}^{1}$ in a Carnot group $\mathbb{G}$ are locally intrinsic Lipschitz graphs of codimension 1.

Theorem A. 5 ([22, Theorem 1.4]). Let $\mathbb{G}$ be a Carnot group of rank $m$ and let $\Sigma_{1}, \ldots, \Sigma_{k}$, $k \leq m$, be hypersurfaces of class $C_{H}^{1}$ with horizontal normals $\nu_{1}, \ldots, \nu_{k}$; let $x \in \Sigma:=\Sigma_{1} \cap$ $\cdots \cap \Sigma_{k}$ be such that $\nu_{1}(x), \ldots, \nu_{k}(x)$ are linearly independent. Consider the vertical plane $\mathbb{W}:=T_{x} \Sigma_{1} \cap \cdots \cap T_{x} \Sigma_{k}$ of codimension $k$ and assume that there exists a complementary homogeneous horizontal subgroup $\mathbb{H}$ such that $\mathbb{G}=\mathbb{W} \mathbb{H}$. Then, there exists an open neighborhood $U$ of $x$ and an intrinsic Lipschitz $\phi: \mathbb{W} \rightarrow \mathbb{H}$ such that

$$
\Sigma \cap U=\operatorname{gr} \phi \cap U
$$

Proof. We work in exponential coordinates associated with an adapted basis $X_{1}, \ldots, X_{n}$ of $\mathfrak{g}$ such that

```
\(\mathbb{H}=\exp \left(\operatorname{span}\left\{X_{1}, \ldots, X_{k}\right\}\right), \quad \mathbb{W}=\exp \left(\left(\operatorname{span}\left\{X_{k+1}, \ldots, X_{s}\right\}\right) \oplus \mathfrak{g}_{2} \oplus \cdots \oplus \mathfrak{g}_{s}\right)\).
```

By definition we can find an open neighborhood $U$ of $x$ and $f=\left(f_{1}, \ldots, f_{k}\right) \in C_{H}^{1}\left(U, \mathbb{R}^{k}\right)$ such that $\Sigma \cap U=\{x \in U: f(x)=0\} \cap U$ and the $m \times k$ matrix-valued function $\nabla_{H} f$ has rank $k$ in $U$. Actually, by our choice of the basis the $k \times k$ minor $M:=$ $\left(X_{1} f(x), \ldots, X_{k} f(x)\right)$ has rank $k$.

Let $\varepsilon$ be a positive number, to be fixed later and only depending on $M$. By Lemma A.4, possibly restricting $U$ we can assume that $f$ is defined on the whole $\mathbb{G}$, that $f \in C_{H}^{1}\left(\mathbb{G}, \mathbb{R}^{k}\right)$ and $\left|\nabla_{H} f-\nabla_{H} f(x)\right|<\varepsilon$; in particular,

$$
\left|\left(X_{1} f, \ldots, X_{k} f\right)-M\right|<\varepsilon \quad \text { on } \mathbb{G} .
$$

It is enough to prove that the level set $R:=\{x \in \mathbb{G}: f(x)=0\}$ is an intrinsic Lipschitz graph. We divide the proof of this claim into two steps.

Step 1: $R$ is the intrinsic graph of some $\phi: \mathbb{W} \rightarrow \mathbb{H}$. It is enough to show that, for any $w \in \mathbb{W}$, there exists a unique $h \in \mathbb{H}$ such that $f(w h)=0$; in particular, this allows to define the map $\phi$ by $\phi(w):=h$.

The map $\left(h_{1}, \ldots, h_{k}\right) \longleftrightarrow \exp \left(h_{1} X_{1}+\cdots+h_{k} X_{k}\right)$ is a group isomorphism between $\mathbb{H}$ and $\mathbb{R}^{k}$. Upon identifying $\mathbb{H}$ and $\mathbb{R}^{k}$ in this way, for any $w \in \mathbb{W}$ we can consider $f_{w}: \mathbb{R}^{k} \rightarrow \mathbb{R}^{k}$ defined by $f_{w}(h):=f(w h)$. This map is of class $C^{1}$ and

$$
\nabla f_{w}(h)=\left(X_{1} f(w h), \ldots, X_{k} f(w h)\right)
$$

We have $\left|\nabla f_{w}-M\right|<\varepsilon$ which, if $\varepsilon$ is small enough, implies that $f_{w}$ is a $C^{1}$ diffeomorphism of $\mathbb{R}^{k}$ : see e.g. the argument in $[11,3.1 .1] .{ }^{3}$ This concludes the proof of Step 1 ; we notice also that, possibly reducing $\varepsilon$, there exists $c>0$ such that (see again in [11, 3.1.1])

$$
\begin{equation*}
\left|f\left(w h_{1}\right)-f\left(w h_{2}\right)\right|=\left|f_{w}\left(h_{1}\right)-f_{w}\left(h_{2}\right)\right| \geq c\left|h_{1}-h_{2}\right| \quad \forall h_{1}, h_{2} \in \mathbb{R}^{k} \tag{А.34}
\end{equation*}
$$

Step 2: $\phi$ is intrinsic Lipschitz. By Remark A. 2 it is enough to prove that

$$
\operatorname{gr} \phi \cap D(x, \mathbb{H}, \beta)=\{x\} \quad \text { for any } x \in \mathbb{G}
$$

for a suitable $\beta>0$ that we will choose in a moment.
Let then $x \in \operatorname{gr} \phi$ be fixed; consider $x^{\prime} \in D(x, \mathbb{H}, \beta)$, so that $x^{\prime}=x y$ for some $y \in D(\mathbb{H}, \beta)$. By definition, there exists $h \in \mathbb{H}$ such that

$$
d\left(0, h^{-1} y\right)=d(h, y) \leq \beta d(h, 0)
$$

Denoting by $L$ the Lipschitz constant of $f$ we deduce using (A.34) that

[^5]\[

$$
\begin{aligned}
\left|f\left(x^{\prime}\right)\right| & =\left|f\left(x h h^{-1} y\right)-f(x)\right| \\
& \geq|f(x h)-f(x)|-\left|f\left(x h h^{-1} y\right)-f(x h)\right| \geq c\|h\|-L d(h, y) \geq(\tilde{c}-\beta L) d(0, h)
\end{aligned}
$$
\]

for some $\tilde{c}>0$. In particular, if $\beta$ is small enough, one can have $f\left(x^{\prime}\right)=0$ only if $h=0$, which immediately gives $x^{\prime}=x$. This concludes the proof.

## We can eventually prove Theorem 1.4.

Proof of Theorem 1.4. By property $\mathscr{C}_{k}$ and Remark 2.3, the vertical plane $\mathbb{W}:=T_{x} \Sigma_{1} \cap$ $\cdots \cap T_{x} \Sigma_{k}$ admits a complementary horizontal homogeneous subgroup $\mathbb{H}$. One can then easily conclude using Theorems A. 3 and A. 5.

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[^2]:    ${ }^{1}$ Actually, this also follows from Theorem 1.4 with $k=1$.

[^3]:    ${ }^{2}$ The key property to prove this assertion is that the set of points where two $C_{H}^{1}$ hypersurfaces intersect transversally is $\mathcal{H}^{Q-1}$-negligible: this fact holds true in any equiregular Carnot-Carathéodory space, see

[^4]:    e.g. [10]. Actually, in view of Theorem 1.1 we could restrict to the setting of Carnot groups satisfying property $\mathscr{C}_{2}$, where the claim follows from Theorem 1.4.

[^5]:    ${ }^{3}$ The careful reader will notice that the argument in [11,3.1.1] works also when the parameter $\delta$ introduced therein is $+\infty$.

