# CORNERS IN NON-EQUIREGULAR SUB-RIEMANNIAN MANIFOLDS 

Enrico Le Donne ${ }^{1}$, Gian Paolo Leonardi ${ }^{2}$, Roberto Monti ${ }^{3}$<br>and Davide Vittone ${ }^{3}$


#### Abstract

We prove that in a class of non-equiregular sub-Riemannian manifolds corners are not length minimizing. This extends the results of [G.P. Leonardi and R. Monti, Geom. Funct. Anal. 18 (2008) 552-582]. As an application of our main result we complete and simplify the analysis in [R. Monti, Ann. Mat. Pura Appl. (2013)], showing that in a 4-dimensional sub-Riemannian structure suggested by Agrachev and Gauthier all length-minimizing curves are smooth.


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## 1. Introduction

One of the major open problems in sub-Riemannian geometry is the regularity of length-minimizing curves. Indeed, no example of a non-smooth minimizer is known, and even the possibility of minimizers with singularities of corner-type has not yet been excluded in full generality (see Problem II in [1] and the discussion in Sect. 4 of [7]).

In [4], the second and third-named authors introduced a shortening technique specifically designed for showing the non-minimality of curves with corner-type singularities (see also the developments in [8]). This technique works for a class of equiregular sub-Riemannian manifolds satisfying the technical condition (1.2) below.

In this paper, we prove the non-minimality of corners in a class of sub-Riemannian manifolds of nonequiregular type. Namely, we show that if the horizontal distribution satisfies condition (1.3) below at the corner point, then the curve is not length minimizing. In this case, the construction of a competitor shorter than the corner is simpler than the one in [4] and relies on the Nagel-Stein-Wainger estimates [9] for the Carnot-Carathèodory distance.

[^0]Let $M$ be an $n$-dimensional differentiable manifold and let $\mathcal{D} \subset T M$ be a smooth subbundle of rank $m$, for some $2 \leq m \leq n$. Then, $\mathcal{D}(x) \subset T_{x} M$ is an $m$-dimensional subspace of the tangent space $T_{x} M$, for all $x \in M$. Let $X_{1}, \ldots, X_{m}$ be a frame of smooth vector fields that form a basis for $\mathcal{D}(x)$, that is $\mathcal{D}(x)=$ $\operatorname{span}\left\{X_{1}(x), \ldots, X_{m}(x)\right\}$ for each $x \in M$. This frame always exists locally. The assumption that the fiber $\mathcal{D}(x)$ has a constant dimension on $M$ plays no role in our argument and can be dropped. Here and in the following, the word "smooth" refers to $C^{\infty}$ regularity.

We denote by $\mathcal{I}=\bigcup_{i \geq 1}\{1, \ldots, m\}^{i}$ the set of admissible multi-indices. For any $\beta=\left(\beta_{1}, \ldots, \beta_{i}\right) \in \mathcal{I}$, for some $i \geq 1$, let us define the iterated commutator

$$
\begin{equation*}
X_{\beta}=\left[X_{\beta_{i}},\left[X_{\beta_{i-1}}, \ldots\left[X_{\beta_{2}}, X_{\beta_{1}}\right] \ldots\right]\right] . \tag{1.1}
\end{equation*}
$$

We say that length $(\beta)=i$ is the length of the multi-index $\beta$. Analogously, we say that length $\left(X_{\beta}\right)=i$ is the length of the iterated commutator $X_{\beta}$. For any point $x \in M$ and $i \geq 1$, let

$$
\mathcal{D}_{i}(x)=\operatorname{span}\left\{X_{\beta}(x) \in T_{x} M: \operatorname{length}\left(X_{\beta}\right)=i\right\} .
$$

Finally, we let $\mathcal{L}_{i}(x)=\mathcal{D}_{1}(x)+\ldots+\mathcal{D}_{i}(x)$ for $i \geq 1$, and we also agree that $\mathcal{L}_{0}(x)=\{0\}$. We assume that $\mathcal{D}$ is bracket generating, i.e., for any $x \in M$ there exists an index $i \in \mathbb{N}$ such that $\mathcal{L}_{i}(x)=T_{x} M$.

An absolutely continuous curve $\gamma:[0,1] \rightarrow M$ is said to be horizontal with respect to the distribution $\mathcal{D}$ (or simply $\mathcal{D}$-horizontal) if there exist bounded measurable functions $h_{1}, \ldots, h_{m}:[0,1] \rightarrow \mathbb{R}$ such that

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

Let $g(x ; \cdot)$ be a positive quadratic form (metric) on $\mathcal{D}(x), x \in M$. The length of $\gamma$ in the sub-Riemannian manifold ( $M, \mathcal{D}, g$ ) is defined as

$$
L(\gamma)=\int_{0}^{1} \sqrt{g(\gamma(t) ; \dot{\gamma}(t))} \mathrm{d} t
$$

and the sub-Riemannian distance between two points $x, y \in M$ is defined as

$$
d(x, y):=\inf \{L(\gamma): \gamma \in A C([0,1] ; M) \text { horizontal, } \gamma(0)=x, \gamma(1)=y\} .
$$

When $M$ is connected, the above set is always nonempty because the distribution $\mathcal{D}$ is bracket-generating, and $d$ is a distance on $M$. Finally, we say that a horizontal curve $\gamma$ joining $x$ to $y$ minimizes the sub-Riemannian length (i.e., it is a length minimizer) if $L(\gamma)=d(x, y)$.

Let $\gamma:[0,1] \rightarrow M$ be a $\mathcal{D}$-horizontal curve. When they exist, we denote by $\dot{\gamma}_{L}(t)$ and $\dot{\gamma}_{R}(t)$ the left and right derivative of $\gamma$ at the point $t \in(0,1)$. We say that $\gamma$ has a corner at the point $x=\gamma(t) \in M$, if the left and right derivatives at $t$, do exist and are linearly independent. In [4], it is shown that if the distribution $\mathcal{D}$ is equiregular (i.e., for every $i \geq 1$ the dimension of $\mathcal{D}_{i}(x)$ is constant on $M$ ) and satisfies the condition

$$
\begin{equation*}
\left[\mathcal{D}_{i}, \mathcal{D}_{j}\right] \subset \mathcal{L}_{i+j-1}, \quad \text { for } i, j \geq 2 \text { such that } i+j>4 \tag{1.2}
\end{equation*}
$$

then corners in $(M, \mathcal{D}, g)$ are not length minimizing. In this paper, we prove that if the distribution $\mathcal{D}$ satisfies at some point $x \in M$ the condition

$$
\begin{equation*}
\mathcal{L}_{i}(x) \neq \mathcal{L}_{i-1}(x) \quad \Rightarrow \quad \mathcal{L}_{i+1}(x)=\mathcal{L}_{i}(x), \quad \text { for all } i \geq 2, \tag{1.3}
\end{equation*}
$$

then corners at $x$ are not length minimizing.

Theorem 1.1. Let $\gamma:[0,1] \rightarrow M$ be a horizontal curve with a corner at the point $x=\gamma(t) \in M$, for $t \in(0,1)$. If the distribution $\mathcal{D}$ satisfies (1.3) at $x$, then $\gamma$ is not length minimizing in $(M, \mathcal{D}, g)$.

The proof of Theorem 1.1, the main result of this paper, is presented in Section 2. After a blow-up argument, we can assume that $M=\mathbb{R}^{n}$, that $\mathcal{D}$ is a 2 -dimensional distribution of planes in $\mathbb{R}^{n}$, and that $\gamma:[-1,1] \rightarrow \mathbb{R}^{n}$ is a corner at the point $0 \in \mathbb{R}^{n}$ of the type

$$
\gamma(t)=\left\{\begin{aligned}
-t x & \text { if } t \in[-1,0] \\
t y & \text { if } t \in(0,1]
\end{aligned}\right.
$$

where $x, y \in \mathbb{R}^{n}$ are linearly independent. We prove the non-minimality of $\gamma$ by an inductive argument on the dimension $n \geq 2$. In the inductive step, we use assumption (1.3) and known estimates on the sub-Riemannian distance to find a competitor shorter than the corner.

We found the basic idea of the Proof of Theorem 1.1 starting from a question raised by Agrachev and Gauthier during the meeting Geometric control and sub-Riemannian geometry held in Cortona in May 2012. They suggested the following situation, in order to find a nonsmooth length-minimizing curve. On the manifold $M=\mathbb{R}^{4}$, let $\Delta$ be the distribution of 2-planes spanned point-wise by the vector fields

$$
\begin{equation*}
X_{1}=\frac{\partial}{\partial x_{1}}+2 x_{2} \frac{\partial}{\partial x_{3}}+x_{3}^{2} \frac{\partial}{\partial x_{4}}, \quad X_{2}=\frac{\partial}{\partial x_{2}}-2 x_{1} \frac{\partial}{\partial x_{3}} \tag{1.4}
\end{equation*}
$$

The distribution $\Delta$ satisfies (1.3). We fix on $\Delta$ the quadratic form $g$ making $X_{1}$ and $X_{2}$ orthonormal.
Let $\alpha>0$ be a parameter and consider the initial and final points $y=(-1, \alpha, 0,0) \in \mathbb{R}^{4}$ and $x=(1, \alpha, 0,0) \in$ $\mathbb{R}^{4}$, respectively. Agrachev and Gauthier asked whether the corner $\gamma:[-1,1] \rightarrow \mathbb{R}^{4}$ joining $y$ to $x$

$$
\begin{equation*}
\gamma_{1}(t)=t, \quad \gamma_{2}(t)=\alpha|t|, \quad \gamma_{3}(t)=0, \quad \gamma_{4}(t)=0, \quad t \in[-1,1] \tag{1.5}
\end{equation*}
$$

is, for small $\alpha>0$, a length minimizer in $\left(\mathbb{R}^{4}, \Delta, g\right)$. The presence of the variable $x_{3}$ in the coefficients of the vector field $X_{1}$ in (1.4) is the technical obstruction for the application of the results of [4].

In [6], the curve $\gamma$ in (1.5) was shown not to be length minimizing for $\alpha \neq 1$, by the explicit construction of a shorter competitor. This answered the above question in the negative. The case $\alpha=1$, however, was left open.

In Section 3, as an application of Theorem 1.1, we prove the following result.
Theorem 1.2. Let $g$ be any smooth metric on $\Delta$. In the sub-Riemannian manifold $\left(\mathbb{R}^{4}, \Delta, g\right)$ all length minimizing curves are smooth and, in particular, no corner is length minimizing.

The Proof of Theorem 1.2 relies on Theorem 1.1. The inductive base is provided by the regularity of geodesics in the first Heisenberg group. This proof covers in particular the case $\alpha=1$ in (1.5) and is simpler than the one in [6].

## 2. Proof of Theorem 1.1

The first step of the proof is a blow-up argument that closely follows [4].
Let $\gamma:[-1,1] \rightarrow M$ be a horizontal curve with a corner at the point $x=\gamma(0) \in M$. We can choose smooth and linearly independent vector fields $X_{1}, X_{2} \in \mathcal{D}$ such that $X_{1}(x)=\dot{\gamma}_{R}(0)$ and $X_{2}(x)=-\dot{\gamma}_{L}(0)$ and we complete $X_{1}, X_{2}$ to a (local) frame $X_{1}, \ldots, X_{m}$ for $\mathcal{D}$. Then we complete $X_{1}, \ldots, X_{m}$ to a frame $X_{1}, \ldots, X_{n}$ for $T M$ in the following way. We choose iterated commutators $X_{m+1}, \ldots, X_{n} \in\left\{X_{\beta}: \beta \in \mathcal{I}\right.$, length $\left.(\beta) \geq 2\right\}$ such that $X_{1}(x), \ldots, X_{n}(x)$ are linearly independent. This choice is possible because $\mathcal{D}$ is bracket generating at $x$. We can also assume that $j \leq k$ implies length $\left(X_{j}\right) \leq \operatorname{length}\left(X_{k}\right)$.

In a neighbourhood of $x \in M$, we fix exponential coordinates of the first type induced by the frame $X_{1}, \ldots, X_{n}$ starting from $x$. Then we can identify $M$ with $\mathbb{R}^{n}, X_{1}, \ldots, X_{n}$ with vector fields on $\mathbb{R}^{n}$, and $x$ with $0 \in \mathbb{R}^{n}$. The fact that we have exponential coordinates of the first type means that for $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ (in fact, for $x$ belonging to a neighbourhood of $0 \in \mathbb{R}^{n}$ ) we have

$$
\begin{equation*}
x=\exp \left(\sum_{i=1}^{n} x_{i} X_{i}\right)(0) \tag{2.1}
\end{equation*}
$$

Here, the exponential mapping is defined by $\exp (X)(0)=\gamma(1)$ where $\gamma$ is the solution of $\dot{\gamma}=X(\gamma)$ and $\gamma(0)=0$.
We assign to the coordinate $x_{i}$ the weight $w_{i}=\operatorname{length}\left(X_{i}\right), i=1, \ldots, n$. Then we have $w_{1}=\ldots=w_{m}=1$. The natural dilations on $\mathbb{R}^{n}$ adapted to the frame $X_{1}, \ldots, X_{n}$ are

$$
\begin{equation*}
\delta_{\lambda}(x)=\left(\lambda^{w_{1}} x_{1}, \lambda^{w_{2}} x_{2}, \ldots, \lambda^{w_{n}} x_{n}\right), \quad x \in \mathbb{R}^{n}, \lambda>0 \tag{2.2}
\end{equation*}
$$

Let $X=X_{\beta}$ be any iterated commutator of the vector fields $X_{1}, \ldots, X_{m}$. Then we have

$$
\begin{equation*}
X=\sum_{i=1}^{n} a_{i}(x) \frac{\partial}{\partial x_{i}} \tag{2.3}
\end{equation*}
$$

where $a_{i} \in C^{\infty}\left(\mathbb{R}^{n}\right), i=1, \ldots, n$, are smooth functions that have the structure described in the following proposition.

Proposition 2.1. There exist polynomials $p_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}$ and functions $r_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}, i=1, \ldots, n$, such that:
(i) $\quad a_{i}(x)=p_{i}(x)+r_{i}(x), x \in \mathbb{R}^{n}$;
(ii) $p_{i}\left(\delta_{\lambda}(x)\right)=\lambda^{w_{i}-\text { length }(X)} p_{i}(x)$;
(iii) $\lim _{\lambda \rightarrow \infty} \lambda^{w_{i}-\text { length }(X)} r_{i}\left(\delta_{1 / \lambda}(x)\right)=0, x \in \mathbb{R}^{n}$.

Proposition 2.1 can be proved as in ([5], p. 306). We omit the details, here. For $\lambda>0$, we let

$$
\begin{equation*}
X^{\lambda}(x)=\sum_{i=1}^{n} \lambda^{w_{i}-\text { length }(X)} a_{i}\left(\delta_{1 / \lambda}(x)\right) \frac{\partial}{\partial x_{i}}, \quad x \in \mathbb{R}^{n} \tag{2.4}
\end{equation*}
$$

The mapping $X \mapsto X^{\lambda}$ is bracket-preserving. Namely, for any multi-index $\beta \in \mathcal{I}$ and for $i=1, \ldots, m$ we have

$$
\begin{equation*}
\left[X_{i}, X_{\beta}\right]^{\lambda}=\left[X_{i}^{\lambda}, X_{\beta}^{\lambda}\right], \quad \lambda>0 \tag{2.5}
\end{equation*}
$$

We let $\mathcal{D}^{\lambda}=\operatorname{span}\left\{X_{1}^{\lambda}, \ldots, X_{m}^{\lambda}\right\}, \mathcal{D}_{i}^{\lambda}=\operatorname{span}\left\{X_{\beta}^{\lambda}:\right.$ length $\left.(\beta)=i\right\}$, and $\mathcal{L}_{i}^{\lambda}=\mathcal{D}_{1}^{\lambda}+\ldots+\mathcal{D}_{i}^{\lambda}$. By (2.5), from (1.3) we deduce that at the point $x=0$ we have

$$
\begin{equation*}
\mathcal{L}_{i}^{\lambda} \neq \mathcal{L}_{i-1}^{\lambda} \quad \Rightarrow \quad \mathcal{L}_{i+1}^{\lambda}=\mathcal{L}_{i}^{\lambda}, \quad \text { for } i \geq 2 \tag{2.6}
\end{equation*}
$$

By Proposition 2.1, for any iterated commutator $X=X_{\beta}$ as in (2.3), we can define the vector field $X^{\infty}$ in $\mathbb{R}^{n}$

$$
X^{\infty}(x)=\lim _{\lambda \rightarrow \infty} X^{\lambda}(x)=\sum_{i=1}^{n} p_{i}(x) \frac{\partial}{\partial x_{i}}, \quad x \in \mathbb{R}^{n}
$$

where $p_{i}, i=1, \ldots, n$, are polynomials such that $p_{i} \circ \delta_{\lambda}=\lambda^{w_{i}-\operatorname{length}(X)} p_{i}$. In particular, if $w_{i}<\operatorname{length}(X)$ then $p_{i}=0$. Passing to the limit as $\lambda \rightarrow \infty$ in (2.5), we see that also the mapping $X \mapsto X^{\infty}$ is bracket-preserving, i.e., $\left[X_{i}, X_{\beta}\right]^{\infty}=\left[X_{i}^{\infty}, X_{\beta}^{\infty}\right]$. Then at the point $x=0$, condition (2.6) holds also for $\lambda=\infty$.

Let $g(x ; \cdot)$ be a metric on $\mathcal{D}(x)$. On the distribution $\mathcal{D}^{\lambda}, \lambda>0$, we introduce the metric $g^{\lambda}(x ; \cdot)$ defined by

$$
g^{\lambda}\left(x ; X^{\lambda}\right)=g\left(\delta_{1 / \lambda}(x) ; X\right), \quad x \in \mathbb{R}^{n},
$$

and on $\mathcal{D}^{\infty}=\operatorname{span}\left\{X_{1}^{\infty}, \ldots, X_{m}^{\infty}\right\}$ we introduce the metric $g^{\infty}(x ; \cdot)$ defined by

$$
g^{\infty}\left(x ; X^{\infty}\right)=\lim _{\lambda \rightarrow \infty} g^{\lambda}\left(x ; X^{\lambda}\right)=g(0 ; X), \quad x \in \mathbb{R}^{n} .
$$

We blow up the curve $\gamma$ at the corner point $0 \in \mathbb{R}^{n}$. For $\lambda>0$ and $t \in[-\lambda, \lambda]$, let $\gamma^{\lambda}(t)=\delta_{\lambda} \gamma(t / \lambda)$. Because $\dot{\gamma}_{R}(0)=X_{1}(0), \dot{\gamma}_{L}(0)=-X_{2}(0)$, we obtain the limit curve $\gamma^{\infty}=\lim _{n \rightarrow \infty} \gamma^{\lambda}$,

$$
\gamma^{\infty}(t)=\left\{\begin{array}{c}
\mathrm{e}_{1} t, t \in[0,1]  \tag{2.7}\\
-\mathrm{e}_{2} t, t \in[-1,0),
\end{array}\right.
$$

where $\mathrm{e}_{1}=(1,0, \ldots, 0)$ and $\mathrm{e}_{2}=(0,1,0, \ldots, 0)$.
Proposition 2.2. If the curve $\gamma$ is length minimizing in $(M, \mathcal{D}, g)$ then the curve $\gamma^{\infty}$ is length minimizing in $\left(\mathbb{R}^{n}, \mathcal{D}^{\infty}, g^{\infty}\right)$.

Proposition 2.2 is proved in [4], Proposition 2.1. Our goal is to prove that the corner $\gamma^{\infty}$ is not length minimizing in $\left(\mathbb{R}^{n}, \mathcal{D}^{\infty}, g^{\infty}\right)$. Thus, we can without loss of generality assume that $M=\mathbb{R}^{n}, \mathcal{D}=\mathcal{D}^{\infty}$, and $\gamma=\gamma^{\infty}$. Since $\gamma$ is contained in the orbit of the distribution $\operatorname{span}\left\{X_{1}, X_{2}\right\}$, we can also assume that $m=2$. Finally, we can pass to exponential coordinates of the second type associated with $X_{1}, \ldots, X_{n}$. Namely, we can assume that for all $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ belonging to a neighbourhood of $0 \in \mathbb{R}^{n}$ we have

$$
x=\exp \left(x_{1} X_{1}\right) \circ \ldots \circ \exp \left(x_{n} X_{n}\right)(0) .
$$

Then we can also assume that $X_{1}$ and $X_{2}$ are vector fields in $\mathbb{R}^{n}$ of the form

$$
\begin{equation*}
X_{1}=\frac{\partial}{\partial x_{1}} \quad \text { and } \quad X_{2}=\frac{\partial}{\partial x_{2}}+\sum_{i=3}^{n} p_{i}(x) \frac{\partial}{\partial x_{i}}, \tag{2.8}
\end{equation*}
$$

where $p_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}, i=3, \ldots, n$, are polynomials of the variable $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ such that $p_{i}\left(\delta_{\lambda}(x)\right)=$ $\lambda^{w_{i}-1} p_{i}(x)$.

Condition (2.6) passes to the limit as $\lambda \rightarrow \infty$. Then, assumption (1.3) at the point $x=0$ reads

$$
\begin{equation*}
\mathcal{L}_{i}(0) \neq \mathcal{L}_{i-1}(0) \quad \Rightarrow \quad \mathcal{L}_{i+1}(0)=\mathcal{L}_{i}(0), \quad \text { for all } i \geq 2 . \tag{2.9}
\end{equation*}
$$

Condition $\mathcal{L}_{i+1}(0)=\mathcal{L}_{i}(0)$ is equivalent to $\mathcal{D}_{i+1}(0) \subset \mathcal{L}_{i}(0)$. As $\mathcal{D}=\mathcal{D}^{\infty},(2.9)$ is equivalent to

$$
\begin{equation*}
\mathcal{L}_{i}(0) \neq \mathcal{L}_{i-1}(0) \quad \Rightarrow \quad \mathcal{D}_{i+1}(0)=\{0\}, \quad \text { for all } i \geq 2 \tag{2.10}
\end{equation*}
$$

We prove this claim. Let $X_{\beta}$ be an iterated commutator such that length $(\beta)=i+1$, i.e., $X_{\beta} \in \mathcal{D}_{i+1}$. According to Proposition 2.1, we have

$$
X_{\beta}=\sum_{j=1}^{n} p_{j}(x) \frac{\partial}{\partial x_{j}},
$$

where $p_{j}$ are polynomials satisfying $p_{j}\left(\delta_{\lambda}(x)\right)=\lambda^{w_{j}-i-1} p_{j}(x)$. Then the sum above ranges over indices $j$ such that $w_{j} \geq i+1$. On the other hand, $\mathcal{L}_{i}(0)=\operatorname{span}\left\{\frac{\partial}{\partial x_{j}}: w_{j} \leq i\right\}$ and thus if $X_{\beta}(0) \in \mathcal{L}_{i}(0)$ we conclude that $X_{\beta}(0)=0$. This proves $(2.10)$.

For given indices $j, k=1, \ldots, n$, we say that $j \prec k$ if there exists $i \geq 2$ such that $X_{j}(0) \in \mathcal{D}_{i}(0)$ but $X_{k}(0) \notin \mathcal{L}_{i}(0)$. From (2.10), we deduce that the weights $w_{1}, \ldots, w_{n}$ satisfy the following condition:

$$
\begin{equation*}
j \prec k \quad \Rightarrow \quad w_{j}+2 \leq w_{k} . \tag{2.11}
\end{equation*}
$$

We are ready to prove that the corner $\gamma=\gamma^{\infty}$ in (2.7) is not a length minimizer in $\left(\mathbb{R}^{n}, \mathcal{D}, g\right)$; this will give a contradiction and prove Theorem 1.1.

We can without loss of generality assume that $X_{1}$ and $X_{2}$ in (2.8) are orthonormal with respect to the metric $g$. This is because two different metrics are locally equivalent and the equivalence constants do not affect our estimates below. Then the length of a $\mathcal{D}$-horizontal curve $\gamma:[0,1] \rightarrow \mathbb{R}^{n}$ is

$$
\begin{equation*}
L(\gamma)=\int_{0}^{1} \sqrt{\dot{\gamma}_{1}(t)^{2}+\dot{\gamma}_{2}(t)^{2}} \mathrm{~d} t \tag{2.12}
\end{equation*}
$$

The proof is by induction on the dimension $n$ of $\mathbb{R}^{n}$. In order to fix the base of induction we distinguish two cases:
(1) We have $\mathcal{L}_{2}(0)=\mathcal{L}_{1}(0)$. In this case, the base of induction is $n=2$. On $\mathbb{R}^{2}$ we have the standard Euclidean metric and corners are not length minimizing.
(2) We have $\mathcal{L}_{2}(0) \neq \mathcal{L}_{1}(0)$. In this case, the base of induction is $n=3$. On $\mathbb{R}^{3}$ we have the Heisenberg group structure. We know that corners are not length minimizing for any sub-Riemannian metric in the Heisenberg group.

We assume that the claim holds for $n-1$ with $n \geq 3,4$ in the two cases, and we prove it for $n$.
Let $\pi: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n-1}$ be the projection $\pi\left(x_{1}, \ldots, x_{n}\right)=\left(x_{1}, \ldots, x_{n-1}\right)$, and define the vector fields $\widehat{X}_{1}=$ $\pi_{*} X_{1}$ and $\widehat{X}_{2}=\pi_{*} X_{2}$. Recall that for each $j=3, \ldots, n$ the polynomial $p_{j}$ appearing in $X_{2}$ in (2.8) satisfies $p_{j} \circ \delta_{\lambda}=\lambda^{w_{j}-1} p_{j}$ and thus it depends only on the variables $x_{1}, \ldots, x_{j-1}$. In particular, for each $j=3, \ldots, n$ the polynomial $p_{j}$ does not depend on $x_{n}$. It follows that

$$
\widehat{X}_{1}(x)=\frac{\partial}{\partial x_{1}}, \quad \widehat{X}_{2}(x)=\frac{\partial}{\partial x_{2}}+\sum_{j=3}^{n-1} p_{j}(x) \frac{\partial}{\partial x_{j}}, \quad \text { where } x=\left(x_{1}, \ldots, x_{n-1}\right) \in \mathbb{R}^{n-1}
$$

We let $\widehat{\mathcal{D}}=\operatorname{span}\left\{\widehat{X}_{1}, \widehat{X}_{2}\right\}$ and we denote by $\widehat{g}$ the metric on $\widehat{\mathcal{D}}$ that makes $\widehat{X}_{1}$ and $\widehat{X}_{2}$ orthonormal. The distribution $\widehat{\mathcal{D}}$ satisfies (1.3).

The projection of the curve $\gamma$ in (2.7) to $\mathbb{R}^{n-1}$, the curve $\widehat{\gamma}=\pi(\gamma)=\left(\gamma_{1}, \ldots, \gamma_{n-1}\right)$, is a corner at $0 \in \mathbb{R}^{n-1}$. By the inductive assumption, this curve is not length minimizing in $\left(\mathbb{R}^{n-1}, \widehat{\mathcal{D}}, \widehat{g}\right)$. Then there exists a $\widehat{\mathcal{D}}$-horizontal curve $\widehat{\sigma}=\left(\widehat{\sigma}_{1}, \ldots, \widehat{\sigma}_{n-1}\right)$ in $\mathbb{R}^{n-1}$ joining the point $\mathrm{e}_{2} \in \mathbb{R}^{n-1}$ to the point $\mathrm{e}_{1} \in \mathbb{R}^{n-1}$ and satisfying

$$
k:=L(\widehat{\gamma})-L(\widehat{\sigma})>0
$$

Let $\sigma:=\left(\sigma_{1}, \ldots, \sigma_{n}\right)$ be the $\mathcal{D}$-horizontal lift to $\mathbb{R}^{n}$ of the plane curve ( $\widehat{\sigma}_{1}, \widehat{\sigma}_{2}$ ) starting from the initial point $\mathrm{e}_{2}$. Clearly, we have $\sigma_{i}=\widehat{\sigma}_{i}$ for $i=1, \ldots, n-1$ and

$$
L(\sigma)=\int_{0}^{1} \sqrt{\dot{\sigma}_{1}(t)^{2}+\dot{\sigma}_{2}(t)^{2}} \mathrm{~d} t=L(\widehat{\sigma})
$$

Finally, the end-point of $\sigma$ is of the form $\mathrm{e}_{1}+h \mathrm{e}_{n} \in \mathbb{R}^{n}$, for some $h \in \mathbb{R}$.

By our choice of the basis $X_{1}, \ldots, X_{n}$, there exists a multi-index $\beta=\left(\beta_{1}, \ldots, \beta_{i}\right) \in \mathcal{I}, i \geq 3$, such that $X_{n}=X_{\beta}$. Since we are in exponential coordinates and also using Proposition 2.1, we deduce that $X_{n}=\partial / \partial x_{n}$. Thus, we have

$$
\begin{equation*}
\frac{\partial}{\partial x_{n}}=X_{\beta}=\left[X_{\beta_{i}},\left[X_{\beta_{i-1}}, \ldots\left[X_{\beta_{2}}, X_{\beta_{1}}\right] \ldots\right]\right] . \tag{2.13}
\end{equation*}
$$

The integer $w_{n}=i$ is the length of the multi-index. We define the multi-index $\widehat{\beta}=\left(\beta_{1}, \ldots, \beta_{i-1}\right)$, that has length $i-1=w_{n}-1$, and we define the corresponding iterated commutator

$$
Z=X_{\widehat{\beta}}=\left[X_{\beta_{i-1}}, \ldots\left[X_{\beta_{2}}, X_{\beta_{1}}\right] \ldots\right]=\sum_{j=1}^{n} b_{j}(x) \frac{\partial}{\partial x_{j}},
$$

where $b_{j} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ are suitable functions, and, in fact, polynomials. By Proposition 2.1, these polynomials are homogeneous:

$$
b_{j}\left(\delta_{\lambda}(x)\right)=\lambda^{w_{j}-w_{n}+1} b_{j}(x), \quad x \in \mathbb{R}^{n} .
$$

Thus, when $w_{j}-w_{n}+1<0$ the polynomial $b_{j}$ vanishes identically, $b_{j}=0$, and the vector field $Z$ has the form

$$
Z=\sum_{w_{j} \geq w_{n}-1} b_{j}(x) \frac{\partial}{\partial x_{j}} .
$$

If $w_{j}=w_{n}-1$ then $b_{j}(x)$ has homogeneous degree 0 and thus it is constant. On the other hand, we have $\partial / \partial x_{n} \in \mathcal{D}_{i}(0)$ and thus $\mathcal{D}_{i}(0) \neq\{0\}$. From (2.10) it follows that $\mathcal{L}_{i-1}(0)=\mathcal{L}_{i-2}(0)$, that is $\mathcal{D}_{i-1}(0)=\{0\}$. Because we have $Z \in \mathcal{D}_{i-1}$, then $Z(0)=0$ and we conclude that $b_{j}=0$ when $w_{j}=w_{n}-1$ and $Z$ is, in fact, of the form

$$
Z=\sum_{w_{j}=w_{n}} b_{j}(x) \frac{\partial}{\partial x_{j}},
$$

with $b_{j}\left(\delta_{\lambda}(x)\right)=\lambda b_{j}(x)$. Therefore we have

$$
\begin{equation*}
b_{j}(x)=c_{j 1} x_{1}+c_{j 2} x_{2} \tag{2.14}
\end{equation*}
$$

for all $j$ such that $w_{j}=w_{n}$, and for suitable constants $c_{j 1}, c_{j 2}$. Since the coefficients of $X_{2}\left(\right.$ and $\left.X_{1}\right)$ in (2.8) do not contain the variables $x_{j}$ such that $w_{j}=w_{n}$, we infer that

$$
\frac{\partial}{\partial x_{n}}=\left[X_{\beta_{i}}, Z\right]=\sum_{w_{j}=w_{n}} \partial_{\beta_{i}} b_{j}(x) \frac{\partial}{\partial x_{j}},
$$

and this implies that

$$
\begin{equation*}
\frac{\partial}{\partial x_{\beta_{i}}} b_{n}(x)=1 \quad \text { and } \quad \frac{\partial}{\partial x_{\beta_{i}}} b_{j}(x)=0, \quad j \neq n, \tag{2.15}
\end{equation*}
$$

where either $\beta_{i}=1$ or $\beta_{i}=2$. We conclude that either $c_{n 1}=1$ or $c_{n 2}=1$ (or both).
Assume that $c_{n 1}=1$. The proof in the case $c_{n 2}=1$ is analogous. By our choice of the basis $X_{1}, \ldots, X_{n}$, for any $j=3, \ldots, n-1$ there exists a multi-index $\beta^{j} \in \mathcal{I}$ such that

$$
X_{j}=X_{\beta^{j}}=\frac{\partial}{\partial x_{j}}+\sum_{k=j+1}^{n} p_{j k}(x) \frac{\partial}{\partial x_{k}},
$$

for suitable polynomials $p_{j k}$. Thus, at the point $x=\mathrm{e}_{1} \in \mathbb{R}^{n}$, the vectors

$$
X_{1}(x), X_{2}(x), \ldots, X_{n-1}(x), Z(x)
$$

are linearly independent, i.e., they form a basis of $T_{x} \mathbb{R}^{n}$. In particular, the vector field $Z$ is an iterated commutator of $X_{1}$ and $X_{2}$ with length $w_{n}-1$. By the Nagel-Stein-Wainger estimate for the Carnot-Carathèodory distance (see [9] and, in particular, Thm. 4), there exist a neighbourhood $U$ of $x=\mathrm{e}_{1}$ and a constant $C>0$ such that

$$
\begin{equation*}
d(x, \exp (t Z)(x)) \leq C t^{\frac{1}{w_{n}-1}} \quad \text { for all } \exp (t Z)(x) \in U \tag{2.16}
\end{equation*}
$$

Let us fix a positive parameter $\epsilon>0$ and let $\left(\gamma_{1}^{\epsilon}, \gamma_{2}^{\epsilon}\right)$ be the planar curve obtained by the concatenation of the following three curves: the line segment from $(0,1)$ to $(0, \epsilon)$, the curve $\left(\epsilon \sigma_{1}, \epsilon \sigma_{2}\right)$, and the line segment from $(\epsilon, 0)$ to $(1,0)$. When $a=0$ and $b=1$ we consider the same curve but starting from $(1,0)$. Let $\gamma^{\epsilon}=\left(\gamma_{1}^{\epsilon}, \ldots, \gamma_{n}^{\epsilon}\right)$ be the $\mathcal{D}$-horizontal lift of this curve to $\mathbb{R}^{n}$, starting from the point $\mathrm{e}_{2}$ (starting from $\mathrm{e}_{1}$, when $a=0$ and $b=1)$. Notice that the $\mathcal{D}$-horizontal lift of $\left(\epsilon \sigma_{1}, \epsilon \sigma_{2}\right)$ is the curve $\delta_{\epsilon} \circ \sigma=\left(\epsilon^{w_{1}} \sigma_{1}, \ldots, \epsilon^{w_{n}} \sigma_{n}\right)$, by (2.2), and hence the end-point of $\delta_{\epsilon} \circ \sigma$ is the point $\epsilon \mathrm{e}_{1}+\epsilon^{w_{n}} h \mathrm{e}_{n}$. Moving along the vector field $X_{1}$ does not change the $n$th coordinate $x_{n}$, hence we conclude that the final point of $\gamma^{\epsilon}$ is $x^{\epsilon}=\mathrm{e}_{1}+\epsilon^{w_{n}} h \mathrm{e}_{n}$. By (2.14) and (2.15) we have $\mathrm{e}_{1}+\epsilon^{w_{n}} h \mathrm{e}_{n}=\exp \left(\epsilon^{w_{n}} h Z\right)\left(\mathrm{e}_{1}\right)$. Since $x^{\epsilon} \rightarrow x=\mathrm{e}_{1}$ as $\epsilon \rightarrow 0$, by (2.16) and for $\epsilon$ small enough we have

$$
\begin{equation*}
d\left(x, x^{\epsilon}\right) \leq C h^{\frac{1}{w_{n}-1}} \epsilon^{\frac{w_{n}}{w_{n}-1}} . \tag{2.17}
\end{equation*}
$$

The sub-Riemannian length in $\left(\mathbb{R}^{n}, \mathcal{D}, g\right)$ of $\gamma^{\epsilon}$ is

$$
\begin{align*}
L\left(\gamma^{\epsilon}\right) & =(1-\epsilon) L(\gamma)+L\left(\delta_{\epsilon} \circ \sigma\right) \\
& =(1-\epsilon) L(\gamma)+\epsilon L(\sigma) \\
& =L(\gamma)-\epsilon(L(\gamma)-L(\sigma)) \\
& =L(\gamma)-\epsilon(L(\widehat{\gamma})-L(\widehat{\sigma})) \\
& =L(\gamma)-\epsilon k \tag{2.18}
\end{align*}
$$

Thus, from (2.17) and (2.18) we obtain (below we let $y=\mathrm{e}_{2}$ )

$$
\begin{aligned}
d(y, x) & \leq d\left(y, x^{\epsilon}\right)+d\left(x^{\epsilon}, x\right) \\
& \leq L\left(\gamma^{\epsilon}\right)+C \epsilon^{\frac{w_{n}}{w_{n}-1}} h^{\frac{1}{w_{n}-1}} \\
& =L(\gamma)-\epsilon k+C \epsilon^{\frac{w_{n}}{w_{n}-1}} h^{\frac{1}{w_{n}-1}}
\end{aligned}
$$

Since $k>0$, there exists an $\epsilon>0$ such that $C h^{\frac{1}{w_{n}-1}} \epsilon^{\frac{1}{w_{n}-1}}<k / 2$ and hence

$$
d(x, y)<L(\gamma)-\epsilon k / 2<L(\gamma)
$$

This proves that $\gamma$ is not length minimizing in $\left(\mathbb{R}^{n}, \mathcal{D}, g\right)$. This also concludes the proof by induction of Theorem 1.1.

## 3. Proof of Theorem 1.2

Let $\Delta=\operatorname{span}\left\{X_{1}, X_{2}\right\}$ be the distribution of planes in $\mathbb{R}^{4}$ spanned by the vector fields $X_{1}$ and $X_{2}$ in (1.4). We fix the metric $g$ on $\Delta$ making $X_{1}, X_{2}$ an orthonormal frame for $\Delta$. Length minimizers for the sub-Riemannian distance are extremals in the sense of Geometric Control Theory, i.e., they satisfy certain necessary conditions given by Pontryagin Maximum Principle. Extremals may be either normal or abnormal. Normal extremals are always smooth. The following proposition classifies abnormal nonsmooth extremals.

Proposition 3.1. In the structure $\left(\mathbb{R}^{4}, \Delta\right)$, the only nonsmooth abnormal extremals are the curves

$$
\gamma(t)= \begin{cases}\left(-t x_{1},-t x_{2}, 0, a\right) & \text { if } t \in[-1,0]  \tag{3.1}\\ \left(t y_{1}, t y_{2}, 0, a\right) & \text { if } t \in(0,1]\end{cases}
$$

where $a \in \mathbb{R}$ and $\left(x_{1}, x_{2}\right),\left(y_{1}, y_{2}\right) \in \mathbb{R}^{2}$ are linearly independent.

Proof. Let $\gamma:[0,1] \rightarrow \mathbb{R}^{4}$ be an abnormal extremal of the distribution $\Delta$. By Pontryagin Maximum Principle, there exists an absolutely continuous curve $\xi:[0,1] \rightarrow \mathbb{R}^{4}$ solving almost everywhere the system of differential equations

$$
\begin{equation*}
\dot{\xi}=\left(2 \dot{\gamma}_{2} \xi_{3},-2 \dot{\gamma}_{1} \xi_{3},-2 \gamma_{3} \dot{\gamma}_{1} \xi_{4}, 0\right) \tag{3.2}
\end{equation*}
$$

See, e.g. ([2], Thm. 2.1) for a formulation of Pontryagin Maximum Principle. Moreover, we have $\left\langle X_{1}(\gamma), \xi\right\rangle=$ $\left\langle X_{2}(\gamma), \xi\right\rangle=\left\langle\left[X_{1}, X_{2}\right], \xi\right\rangle=0$, where $\langle\cdot, \cdot\rangle$ is the standard scalar product of $\mathbb{R}^{4}$. The last equation $\left\langle\left[X_{1}, X_{2}\right], \xi\right\rangle=$ 0 is Goh condition, that holds automatically true in the rank 2 case (see, e.g., [10]). Namely, the curve $\xi$ also solves the system of equations

$$
\begin{align*}
& \xi_{1}+2 \gamma_{2} \xi_{3}+\gamma_{3}^{2} \xi_{4}=0 \\
& \xi_{2}-2 \gamma_{1} \xi_{3}=0 \\
& \xi_{3}-\gamma_{1} \gamma_{3} \xi_{4}=0 \tag{3.3}
\end{align*}
$$

From (3.2), we see that that $\xi_{4}$ is constant. This constant is nonzero, otherwise (3.3) would imply $\xi=0$, and this is not possible for abnormal extremals. By linearity we can assume that $\xi_{4}=1$, and thus (3.3) trasforms into the system

$$
\begin{equation*}
\xi=\left(-2 \gamma_{1} \gamma_{2} \gamma_{3}-\gamma_{3}^{2}, 2 \gamma_{1}^{2} \gamma_{3}, \gamma_{1} \gamma_{3}, 1\right) \tag{3.4}
\end{equation*}
$$

and the system (3.2) becomes

$$
\begin{equation*}
\dot{\xi}=\left(2 \gamma_{1} \gamma_{3} \dot{\gamma}_{2},-2 \gamma_{1} \gamma_{3} \dot{\gamma}_{1},-2 \gamma_{3} \dot{\gamma}_{1}, 0\right) \tag{3.5}
\end{equation*}
$$

Differentiating the second equation in (3.4), we find $\dot{\xi}_{2}=4 \gamma_{1} \gamma_{3} \dot{\gamma}_{1}+2 \gamma_{1}^{2} \dot{\gamma}_{3}$, and comparing with the second equation in (3.5), we deduce that $\gamma_{1}^{2}\left(3 \dot{\gamma}_{1} \gamma_{3}+\gamma_{1} \dot{\gamma}_{3}\right)=0$. This in turn implies that the function $\phi(t)=\gamma_{1}(t)^{3} \gamma_{1}(t)$ is a constant $c \in \mathbb{R}$.

Now there are two cases.
First case. $c=0$. In this case, the equation $\gamma_{1}^{3} \gamma_{3}=0$ implies that $\gamma$ is either a line or a corner of the form (3.1).
Second case. $c \neq 0$. In this case, by differentiating the identity $\gamma_{3}=c / \gamma_{1}^{3}$ and using the horizontality condition $\dot{\gamma}_{3}=2 \dot{\gamma}_{1} \gamma_{2}-2 \dot{\gamma}_{2} \gamma_{1}$, we deduce that

$$
\left\langle\left(3 c+2 \gamma_{1}^{4} \gamma_{2},-2 \dot{\gamma}_{1}^{5}\right),\left(\dot{\gamma}_{1}, \dot{\gamma}_{2}\right)\right\rangle=0
$$

where $\langle\cdot, \cdot\rangle$ is the standard scalar product of $\mathbb{R}^{2}$. In other words, the planar curve $\left(\gamma_{1}, \gamma_{2}\right)$ is, up to reparameterization, an integral curve of the vector field in the plane

$$
\frac{\partial}{\partial x_{1}}+\frac{3 c+2 x_{1}^{4} x_{2}}{2 x_{1}^{5}} \frac{\partial}{\partial x_{2}}, \quad x_{1} \neq 0
$$

Thus the curve $\gamma$ is

$$
\begin{equation*}
\gamma(t)=\left(t, b t-\frac{3}{10} c t^{-4}, c t^{-3},-\frac{1}{5} c^{2} t^{-5}+d\right), \quad \text { with } t \neq 0 \tag{3.6}
\end{equation*}
$$

for some $b, c, d \in \mathbb{R}$. All such curves are $C^{\infty}$.
We conclude that the nonsmooth abnormal extremals in $\left(\mathbb{R}^{4}, \Delta\right)$ are precisely the corners $(3.1)$.
Remark 3.2. In the proof of Proposition 3.1, we have the formula (3.4) for the dual curve $\xi$ of an abnormal extremal $\gamma$. The coordinates of $\xi$ are polynomial functions of the coordinates of $\gamma$. This is analogous to the results obtained in $[2,3]$ for stratified nilpotent groups. In such groups, dual curves can be reconstructed using a special family of polynomials, called extremal polynomials, and abnormal extremals are always contained in the level sets of extremal polynomials.

We conclude with the Proof of Theorem 1.2.

Proof of Theorem 1.2. Thanks to Proposition 3.1, it is enough to prove the non-minimality of corners in $\left(\mathbb{R}^{4}, \Delta\right)$ at $x=0$. Since the distribution $\Delta$ satisfies the assumption (1.3) at $x=0$, we can use Theorem 1.1 and obtain the desired conclusion.

## References

[1] A. Agrachev, Some open problems, Geometric Control Theory and sub-Riemannian Geometry. Edited by G. Stefani, U. Boscain, J.-P. Gauthier, A. Sarychev, M. Sigalotti. Vol. 5 of Springer INdAM Series (2014) 1-14.
[2] E. Le Donne, G.P. Leonardi, R. Monti and D. Vittone, Extremal Curves in Nilpotent Lie Groups. Geom. Funct. Anal. 23 (2013) 1371-1401.
[3] E. Le Donne, G.P. Leonardi, R. Monti and D. Vittone, Extremal polynomials in stratified groups. Preprint ArXiv:1307.5235 (2013).
[4] G.P. Leonardi and R. Monti, End-point equations and regularity of sub-Riemannian geodesics. Geom. Funct. Anal. 18 (2008) 552-582.
[5] G.A. Margulis and G.D. Mostow, Some remarks on the definition of tangent cones in a Carnot-Carathéodory space. J. Anal. Math. 80 (2000) 299-317.
[6] R. Monti, A family of nonminimizing abnormal curves. Ann. Mat. Pura Appl. (2013).
[7] R. Monti, The regularity problem for sub-Riemannian geodesics, Geometric Control Theory and sub-Riemannian Geometry. Edited by G. Stefani, U. Boscain, J.-P. Gauthier, A. Sarychev, M. Sigalotti. Vol. 5 of Springer INdAM Series (2014) $313-332$.
[8] R. Monti, Regularity results for sub-Riemannian geodesics. Calc. Var. Partial Differ. Eqs. 49 (2014) 549-582.
[9] A. Nagel, E.M. Stein and S. Wainger, Balls and metrics defined by vector fields. I. Basic properties. Acta Math. 155 (1985) 103-147.
[10] D. Vittone, The regularity problem for sub-Riemannian geodesics. Preprint (2013). Available at http://cvgmt.sns.it/.


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    1 University of Jyväskylä, Department of Mathematics and Statistics, P.O. Box 35, 40014 Jyväskylä, Finland. ledonne@msri. org
    2 Università di Modena e Reggio Emilia, Dipartimento di Scienze Fisiche, Informatiche e Matematiche, via Campi 213/b, 41100 Modena, Italy. gianpaolo.leonardi@unimore.it
    3 Università di Padova, Dipartimento di Matematica, via Trieste 63, 35121 Padova, Italy.
    monti@math.unipd.it; vittone@math.unipd.it

