

## ON THE INNER CONE PROPERTY FOR CONVEX SETS IN TWO-STEP CARNOT GROUPS, WITH APPLICATIONS TO MONOTONE SETS

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**Abstract:** In the setting of two-step Carnot groups we show a “cone property” for horizontally convex sets. Namely, we prove that, given a horizontally convex set  $C$ , a pair of points  $P \in \partial C$  and  $Q \in \text{int}(C)$ , both belonging to a horizontal line  $\ell$ , then an open truncated subRiemannian cone around  $\ell$  and with vertex at  $P$  is contained in  $C$ .

We apply our result to the problem of classification of horizontally monotone sets in Carnot groups. We are able to show that monotone sets in the direct product  $\mathbb{H} \times \mathbb{R}$  of the Heisenberg group with the real line have hyperplanes as boundaries.

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### 1. Introduction and main results

The starting point of this paper is a result of Arena, Caruso, and Monti [2], where it is proved that, in subRiemannian two-step Carnot groups of Métivier type, given a horizontally convex set  $C$ , a pair of points  $P \in \partial C$  and  $Q \in \text{int}(C)$ , both belonging to a horizontal line  $\ell$ , then an open truncated subRiemannian cone around  $\ell$  and with vertex at  $P$  is contained in  $C$ . Here we prove that the same property holds in general two-step Carnot groups without assuming the restrictive Métivier condition.

We apply the construction above to the problem of classification of *precisely monotone* sets, in the sense of Cheeger and Kleiner [7]. Recall that a precisely monotone set  $E$  in a Carnot group  $\mathbb{G}$  is a set which is (horizontally) convex and such that  $\mathbb{G} \setminus E$  is convex too. In [7] it is shown that a precisely monotone set  $E$  in the three-dimensional Heisenberg group  $\mathbb{H}$  satisfies  $\Pi \subset E \subset \overline{\Pi}$  for a suitable half-space  $\Pi$ . The classification is nontrivial, see also [8].

In this paper we show the same result in  $\mathbb{H} \times \mathbb{R}$ , the direct product of  $\mathbb{H}$  with the Euclidean line. Namely, we will prove that the boundary of a (precisely) monotone set in  $\mathbb{H} \times \mathbb{R}$  is a hyperplane. See Theorem 1.2 for

the precise statement. Although our result concerns the seemingly easy model  $\mathbb{H} \times \mathbb{R}$ , the proof requires a certain amount of work. Furthermore, our techniques, based on the use of the cone property above, connected with the notion of *intrinsic graph*, in the sense of [12], or [1], could be useful to attack the analogous classification problem for general two-step Carnot groups. Let us mention that, to the author's knowledge, such problem is at the moment open.

To describe our setting, let  $(z, t) \in Z \times T = \mathbb{R}^m \times \mathbb{R}^\ell = \mathbb{G}$  be equipped with the group law

$$(1.1) \quad (z, t) \cdot (\zeta, \tau) = (z + \zeta, t + \tau + \langle z, A\zeta \rangle) = (z + \zeta, t + \tau + Q(z, \zeta)),$$

where  $A = (A^1, \dots, A^\ell)$  and  $A^\beta \in \mathbb{R}^{m \times m}$  is a skew-symmetric matrix for all  $\beta = 1, \dots, \ell$ . In other words,  $Q: Z \times Z \rightarrow T$  is bilinear and skew-symmetric. The left invariant horizontal vector fields

$$(1.2) \quad X_j := \partial_{z_j} + \sum_{\beta=1}^{\ell} \sum_{k=1}^m a_{kj}^{\beta} z_k \partial_{t_{\beta}}, \quad j = 1, \dots, m,$$

define in a standard way a subRiemannian distance  $d$ . We denote by  $B((x, t), r)$  the corresponding ball of center  $(z, t)$  and radius  $r$ . See [17] for the definition. The *horizontal plane* at a point  $(z, t)$  will be denoted by  $H_{(z, t)} = \text{span}\{X_j(z, t) : 1 \leq j \leq m\}$ . Furthermore, for any  $u \in \mathbb{R}^m$ , letting  $u \cdot X = \sum_{k \leq m} u_k X_k$ , we see that the integral curve of  $u \cdot X$  at time 1 starting from  $(z, t)$  has the form  $e^{u \cdot X}(z, t) = (z + u, t + Q(z, u))$ .

A *horizontal line*, or a *line* for short, is a set of the form

$$(1.3) \quad \ell = \{e^{su \cdot X}(z, t) = (z, t) \cdot (su, 0) = (z + su, t + sQ(z, u)) : s \in \mathbb{R}\},$$

where  $u \in \mathbb{R}^m$  and  $(z, t) \in \mathbb{G}$ . *Lines* in our setting are particular Euclidean lines. Note that the Euclidean line parametrized by the path  $\gamma(s) = (as, bs, cs)$  is a horizontal line if and only if  $c = 0$ . We always parametrize horizontal lines with constant speed.

Finally, we say that a set  $C \subset \mathbb{G}$  is *horizontally convex*, or *convex* for short, if for all  $P, Q$  (horizontally) aligned points contained in  $C$ , the whole line containing  $P$  and  $Q$  is contained in  $C$ . We say instead that  $E \subset \mathbb{G}$  is *horizontally precisely monotone*, or simply *monotone*, if both  $E$  and  $\mathbb{G} \setminus E$  are convex. It is well known that the class of horizontally convex sets includes strictly all sets which are convex in the Euclidean sense.

Convexity in Carnot groups has been introduced in [9] and [14]. Further references on properties of convex sets are [16, 5, 6]. Many authors have studied regularity properties of convex functions. Monotone sets in the Heisenberg group appear and have a prominent role in the recent papers [7, 8, 18, 10]. Observe also that the class of monotone sets is

somewhat similar to the more restricted class of *sets with constant horizontal normal*, appearing in geometric measure theory in Carnot groups. See [12, 13, 4].

Here is our first result.

**Theorem 1.1.** *Let  $C \subset \mathbb{G}$  be a convex set in a two-step Carnot group with law (1.1). Let  $P = (z, t) \in \overline{C}$  and assume that there is  $\xi \in Z = \mathbb{R}^m$  such that  $Q := (z, t) \cdot (\xi, 0) \in \text{int}(C)$ . Then there is  $\varepsilon > 0$  such that*

$$(1.4) \quad \bigcup_{0 < s \leq 1} B((z', t') \cdot (s\xi, 0), \varepsilon s) \subset C \quad \text{for all } (z', t') \in B((z, t), \varepsilon) \cap \overline{C}.$$

An analogous statement holds without assuming that  $(z, t) \cdot (\xi, 0) \in \text{int}(C)$ , but requiring that a surface  $\Sigma$  containing  $Q$  and transversal to the line containing  $P$  and  $Q$  is contained in  $C$ . Note also that the statement holds for  $P = (z, t)$  in the closure of  $C$  and becomes more significant when  $(z, t) \in \partial C$ . Even the situation  $P \in \partial C \setminus C$  is included. In this case, the mere definition of convexity does not ensure that the open horizontal segment connecting  $(z, t)$  with  $(z, t) \cdot (\xi, 0)$  is contained in  $C$ . However, the proof of the theorem shows that such segment is contained in  $\text{int}(C)$ .

The theorem above was already proved by Arena, Caruso, and Monti in the case of two-step groups of Métivier type [2, Theorem 1.4]. Recall that a group of the form (1.1) satisfies the Métivier condition if for all  $t \in T$  and for any non-zero  $z \in Z$ , there is  $\zeta \in Z$  such that  $Q(z, \zeta) = t$ . A similar, more qualitative statement was proved by Cheeger and Kleiner in the setting of the three-dimensional Heisenberg group [7, Proposition 4.6].

In the present paper we generalize the statement to general step-two groups. Our techniques are inspired by the argument in [7], but we exploit some higher-order envelopes, see Subsection 2.1. Namely, the key starting point of the classification argument of [7] relies on the analysis of the map

$$Z \times Z \ni (u_1, u_2) \mapsto \Gamma(u_1, u_2) := e^{u_2 \cdot X} e^{u_1 \cdot X}(0, 0),$$

which turns out to be a submersion at any point  $(\xi, \xi) \neq (0, 0)$  (the argument in [2] is essentially based on a similar property). This property fails to hold if the group does not satisfy the Métivier condition (see the discussion in Subsection 2.1). We are then forced to analyze “higher order” maps of the form

$$(1.5) \quad \begin{aligned} Z^p \ni (u_1, u_2, \dots, u_p) &\mapsto \Gamma(u_1, u_2, \dots, u_p) \\ &\mapsto e^{u_p \cdot X} \dots e^{u_2 \cdot X} e^{u_1 \cdot X}(0, 0), \end{aligned}$$

where  $p$  can be greater than 2. It may happen that at some  $(\xi, \xi, \dots, \xi) \neq 0 \in Z^p$  there is no  $p \in \mathbb{N}$  such that the map  $\Gamma$  is a submersion. However, we are able to show that, given any two-step Carnot group, for sufficiently large  $p$  the map  $\Gamma$  is *open* at any point of the form  $(\xi, \xi, \dots, \xi)$ . See the statement in Theorem 2.1 and see also Remark 2.2. The existence of the cone in (1.4) follows easily.

We apply our construction to the study in a model case of the classification problem for monotone subsets of a Carnot group. It is easy to see that in any Carnot group of step two, given any hyperplane  $\Sigma$ , any of the two open half spaces whose union equals  $\mathbb{G} \setminus \Sigma$  is a monotone set. The natural question is whether any monotone set  $E$  has a hyperplane as boundary. If  $\mathbb{G}$  is of Métivier type, then this turns out to be true under the mild assumption  $\partial E \subsetneq \mathbb{G}$ , which is always satisfied if  $\mathbb{G} = \mathbb{H}^n$  is the  $n$ -th Heisenberg group. The short argument in the proof of Proposition 3.1 has been kindly indicated to the author by Roberto Monti. See Proposition 3.2 for the case of  $\mathbb{H}^n$ .

The first significant example of non Métivier two-step Carnot group appears in the following situation. Consider the direct product  $\mathbb{G} := \mathbb{H}_{(x,y,t)} \times \mathbb{R}_u = \mathbb{R}^4$  with law

$$(x, y, u, t) \cdot (x', y', u', t') = (x + x', y + y', u + u', t + t' + 2(yx' - y'x))$$

and with horizontal left-invariant vector fields  $X = \partial_x + 2y\partial_t$ ,  $Y = \partial_y - 2x\partial_t$ , and  $U = \partial_u$ . We shall prove the following theorem.

**Theorem 1.2.** *Let  $E \subset \mathbb{H} \times \mathbb{R}$  be a monotone set. Then there is a three-dimensional hyperplane  $\Sigma$  such that  $\text{int}(E)$  and  $\text{int}(E^c)$  are the connected components of  $\mathbb{R}^4 \setminus \Sigma$ .*

To prove such a theorem, keeping in mind the argument of [7], we have to introduce some new ideas, in order to have a starting point which is suitable for a higher-dimensional situation. Namely, we will use our result on cones, and the consequent fact that around a point  $P$  at the boundary of  $E$ , either the boundary is contained in the horizontal plane  $H_P$ , or else it is contained in an intrinsic graph. See the discussion in Section 3. We are then able to classify, first locally and then globally, intrinsic graphs whose corresponding epigraphs are monotone sets.

Let us mention that an alternative approach to Theorem 1.2 could be carried out via a parametric analysis of the intersections of  $E$  with sets with  $u = \text{constant}$ , which are basically copies of the Heisenberg group. This alternative method however would not provide any help in the analysis of more general situations. We are choosing instead our approach using graphs with the intention of applying similar ideas to more

general two-step Carnot groups. Concerning the situation in higher step, the examples in [4] suggest that monotone sets in higher step Carnot groups could be very ugly.

The plan of the paper is the following. In Subsection 2.1 we discuss the openness property of the map  $\Gamma$  mentioned in (1.5). In Subsection 2.2 we apply such property to the construction of the cones in Theorem 1.1. In Section 3 we show that the boundary of a monotone set can be written as an intrinsic graph. Section 4 is devoted to the classification of monotone sets in the group  $\mathbb{H} \times \mathbb{R}$ .

## 2. Higher order envelopes, open maps, and cones

**2.1. Multiexponential open maps.** In this section, given  $p \in \mathbb{N}$ , we consider the map  $\Gamma: Z^p \rightarrow Z \times T$ ,

$$(2.1) \quad \begin{aligned} \Gamma(u_1, \dots, u_p) &:= e^{u_p \cdot X} \dots e^{u_2 \cdot X} e^{u_1 \cdot X}(0) \\ &= \left( \sum_{j \leq p} u_j, \sum_{1 \leq j < k \leq p} Q(u_j, u_k) \right). \end{aligned}$$

Here we use the notation  $u \cdot X = \sum_{k=1}^m \langle u, e_k \rangle X_k$  for a vector  $u = (\langle u, e_1 \rangle, \dots, \langle u, e_m \rangle) \in Z = \mathbb{R}^m$ . Our purpose is to prove that the map  $\Gamma$  is open at any point  $(\xi, \dots, \xi)$ .

**Theorem 2.1.** *In any two-step Carnot group with law (1.1) there is  $p \in \mathbb{N}$  such that for all  $\xi \in Z$  the map  $\Gamma$  in (2.1) is open at any point  $(\xi, \dots, \xi) \in Z^p$ . More precisely, for all  $\xi \in Z$  there is  $c_0 > 0$  such that  $\Gamma(B_{\text{Euc}}((\xi, \xi, \dots, \xi), r)) \supseteq B_{\text{Euc}}(\Gamma(\xi, \dots, \xi), c_0 r^2)$  for all  $r \leq c_0$ .*

Here and hereafter, if needed,  $B_{\text{Euc}}$  denotes Euclidean balls. From the proof it will be clear that the number  $p$  can be estimated in terms of the dimension  $m$  of  $Z$ . However, we do not need a sharp estimate of  $p$ .

*Remark 2.2.* If the group satisfies the M  tivier condition, then it suffices to choose  $p = 2$  and it turns out that for all  $\xi \neq 0$  the map  $\Gamma$  is a submersion (see [7]). For more general groups, we are forced to consider larger values of  $p$ . Furthermore, the map  $\Gamma$  cannot be a submersion at any point  $(\xi, \dots, \xi) \in Z^p$ , with any  $p \in \mathbb{N}$  if the curve  $\gamma(s) = (s\xi, 0)$  is an abnormal extremal for the length-minimizing subRiemannian problem (note that the M  tivier condition characterizes two-step Carnot groups where subRiemannian abnormal geodesics do not appear at all, see [15, Proposition 3.6]).

In spite of these pathologies, the “quadratic” openness proved in Theorem 2.1 holds and is sufficient for our purposes. We observe finally that our theorem holds also for  $\xi = 0$ .

In order to prove Theorem 2.1 we need to show some properties of the following map. Define for all  $q \in \mathbb{N}$  the function  $H_{2q}: Z^{2q} \rightarrow T$ ,

$$H_{2q}(z_1, \zeta_1, z_2, \zeta_2, \dots, z_q, \zeta_q) := \sum_{k=1}^q Q(z_k, \zeta_k) \in T.$$

**Lemma 2.3.** *Let  $Z = \mathbb{R}^m$  and assume the Hörmander condition. Then there is  $q \in \mathbb{N}$  and  $c > 0$  such that*

$$(2.2) \quad H_{2q}(B_{\text{Euc}}(0, r)) \supseteq B_{\text{Euc}}(0, cr^2) \quad \text{for all } r \in ]0, +\infty[.$$

From the proof it is clear that, if  $Z$  has dimension  $m$ , then any  $q \geq m-1$  ensures (2.2) for some  $c > 0$ .

*Proof:* Write  $Q(z, \zeta) = \langle z, A\zeta \rangle$ . The step two Hörmander condition ensures that  $\text{span}\{A_{jk} : j < k\} = T$ . Since  $A_{jk} = Q(e_j, e_k)$ , we conclude that for all  $t \in T$  there is a vector  $(c_{jk})_{1 \leq j < k \leq m}$  such that

$$(2.3) \quad \sum_{j < k} c_{jk} Q(e_j, e_k) = t.$$

Writing the left-hand side in the form  $\sum_{k=2}^m Q(\sum_{j=1}^{k-1} c_{jk} e_j, e_k)$ , we see that the map  $H_{2q}$  with  $q = m-1$  is onto.

Next, to finish the proof we use the following homogeneity argument. Choose  $(c_{jk})_{j < k}$  such that (2.3) holds with estimate  $\sum |c_{jk}| \leq C|t|$  (this can be done with  $C$  uniform in  $t$  by elementary linear algebra). Then for all  $k = 2, \dots, m$  we let  $w_k := \sum_{j=1}^{k-1} c_{jk} e_j$  (so that  $|w_k| \leq C|t|$ ) and we choose

$$z_k := |w_k|^{-1/2} w_k \quad \text{and} \quad \zeta_k = |w_k|^{1/2} e_k$$

(if  $w_k = 0$ , then we take  $z_k = \zeta_k = 0$ ). Thus, we get  $\sum_{k=2}^m Q(z_k, \zeta_k) = t$  with the required estimate  $|z_k|, |\zeta_k| \leq C|t|^{1/2}$  for all  $k = 2, \dots, m$ .  $\square$

*Proof of Theorem 2.1:* Let  $\xi \in Z$  and define

$$\begin{aligned} f_\xi(u_1, \dots, u_p) &:= \Gamma(\xi + u_1, \dots, \xi + u_p) - \Gamma(\xi, \dots, \xi) \\ &= \left( \sum_{j \leq p} u_j, Q \left( \sum_{i=1}^p (p-2i+1) u_i, \xi \right) + \sum_{1 \leq j < k \leq p} Q(u_j, u_k) \right). \end{aligned}$$

We will show that for each  $(z, t) \in Z \times T$  there is  $(u_1, \dots, u_p) \in Z^p$  such that

$$f_\xi(u_1, \dots, u_p) = (z, t) \quad \text{and} \quad |u_1| + \dots + |u_p| \leq C(|z| + |t|^{1/2}).$$

Let  $(z, t) \in Z \times T$ . Write the system

$$\left( \sum_{j \leq p} u_j, Q \left( \sum_{i=1}^p (p-2i+1) u_i, \xi \right) + \sum_{1 \leq j < k \leq p} Q(u_j, u_k) \right) = (z, t).$$

Since we cannot expect that the map is a submersion for all  $\xi$ , see Remark 2.2, we give a discussion independent of  $\xi$  by restricting ourselves to the vectors  $u_j$  such that  $\sum_{i=1}^p (p-2i+1) u_i = 0$  or, equivalently,

$$(2.4) \quad u_p = \sum_{j=1}^{p-1} \frac{p-2j+1}{p-1} u_j.$$

In this way, we are led to work with a new system on  $Z^{p-1}$  of the form

$$\left( \sum_{j \leq p-1} u_j + \sum_{j \leq p-1} \frac{p-2j+1}{p-1} u_j, \sum_{j < k \leq p-1} Q(u_j, u_k) + Q \left( \sum_{k \leq p-1} u_k, \sum_{j \leq p-1} \frac{p-2j+1}{p-1} u_j \right) \right) = (z, t).$$

The  $Z$ -component is  $\sum_{j \leq p-1} \frac{2(p-j)}{p-1} u_j = z$  and, after some calculation, we see that the  $T$ -component is

$$\begin{aligned} & \sum_{j < k \leq p-1} Q(u_j, u_k) + Q \left( \sum_{k \leq p-1} u_k, \sum_{j \leq p-1} \frac{p-2j+1}{p-1} u_j \right) \\ &= \sum_{j < k} \left\{ 1 + \frac{p-2k+1}{p-1} - \frac{p-2j+1}{p-1} \right\} Q(u_j, u_k) \\ &= \frac{1}{p-1} \sum_{j < k \leq p-1} (p+2j-2k-1) Q(u_j, u_k). \end{aligned}$$

Ultimately, we get the system

$$\begin{cases} \sum_{j \leq p-1} (p-j) u_j = \frac{p-1}{2} z, \\ \sum_{j < k \leq p-1} (p+2j-2k-1) Q(u_j, u_k) = (p-1) t. \end{cases}$$

From the first equation we find

$$(2.5) \quad u_{p-1} = \frac{(p-1)z}{2} - \sum_{j \leq p-2} (p-j) u_j.$$

Therefore, the second equation becomes

$$\begin{aligned} & \sum_{j < k \leq p-2} (p-1-2(k-j))Q(u_j, u_k) \\ & + \sum_{j \leq p-2} (-p+1+2j)Q\left(u_j, \frac{p-1}{2}z - \sum_{k \leq p-2} (p-k)u_k\right) = (p-1)t. \end{aligned}$$

After some simplifications we get

$$\begin{aligned} & Q\left(\sum_{j \leq p-2} (-p+1+2j)u_j, \frac{p-1}{2}z\right) \\ & + (p-1) \sum_{j < k \leq p-2} (1+k-j)Q(u_j, u_k) = (p-1)t. \end{aligned}$$

Next, we eliminate the first term by choosing  $\sum_{j \leq p-2} (-p+2j+1)u_j = 0$  or, in other words,

$$(2.6) \quad u_{p-2} = \sum_{j \leq p-3} \frac{p-2j-1}{p-3} u_j.$$

We get

$$t = \sum_{j < k \leq p-3} (1+k-j)Q(u_j, u_k) + \sum_{j \leq p-3} (p-1-j)Q\left(u_j, \sum_{k \leq p-3} \frac{p-2k-1}{p-3} u_k\right),$$

which, after some manipulations takes the simple form

$$(2.7) \quad \sum_{1 \leq j < k \leq p-3} (p-3-2(k-j))Q(u_j, u_k) = (p-3)t.$$

In order to solve system (2.7), we take  $p-3 = 2\ell+1$  a sufficiently large odd number and we use Lemma 2.4 below. This ensures that for all  $t \in T$  there are  $u_1, \dots, u_{p-3}$  which solve (2.7) and such that  $|u_1| + \dots + |u_{p-3}| \leq C|t|^{1/2}$ . Then, by (2.6) we find  $u_{p-2}$  with estimate  $|u_{p-2}| \leq C|t|^{1/2}$ . Furthermore, (2.5) gives the value of  $u_{p-1}$  and the estimate  $|u_{p-1}| \leq C(|z| + |t|^{1/2})$ . To conclude the proof, we get  $u_p$  from (2.4) again evaluating with  $|z| + |t|^{1/2}$ .  $\square$

In order to state and prove Lemma 2.4 define, for  $q \geq 3$ ,

$$(2.8) \quad P_q(u_1, \dots, u_q) := \sum_{1 \leq j < k \leq q} (q-2(k-j))Q(u_j, u_k).$$

It is easy to see that  $P_2 = 0$  and we agree to let  $P_1 = 0$ , if needed.



**Lemma 2.4** (Properties of the function (2.8)). *Let  $q = 2\ell + 1 \geq 5$  be an odd number. Then there is a linear change of variable  $T_\ell: Z^{2\ell+1} \rightarrow Z^{2\ell+1}$ , which we denote by  $(u_1, \dots, u_{2\ell+1}) \mapsto (v_1, \dots, v_{2\ell+1})$ , such that*

$$P_{2\ell+1}(u_1, \dots, u_{2\ell+1}) = \sum_{k=1}^{\ell} \frac{2\ell+1}{2k+1} Q(v_{2k}, v_{2k+1}).$$

*In particular, by Lemma 2.3, if the Hörmander condition holds, then the map  $P_{2\ell+1}$  is open at the origin in the sense of (2.2) for sufficiently large  $\ell$ .*

Note that the variable  $v_1$  does not appear after the change of variable. An analogous version (not needed for our purposes) holds for  $P_{2\ell}$  with  $2\ell \geq 4$ .

*Proof of Lemma 2.4:* Let  $q \geq 5$  be an odd number. Let us make the linear change of variable

$$(2.9) \quad v_{q-1} = \sum_{j=1}^{q-1} (2j-q)u_j \quad \text{and} \quad v_q = u_q - \sum_{j=1}^{q-2} \frac{-q+2+2j}{q-2} u_j,$$

which leaves the variables  $u_1, \dots, u_{q-2}$  unchanged. We have the following recursive relation

$$(2.10) \quad P_q(u_1, \dots, u_q) := Q(v_{q-1}, v_q) + \frac{q}{q-2} P_{q-2}(u_1, \dots, u_{q-2}).$$

To prove (2.10), write

$$P_q(u_1, \dots, u_q) = \sum_{1 \leq j < k \leq q-1} (q-2k+2j)Q(u_j, u_k) + \sum_{j \leq q-1} (2j-q)Q(u_j, u_q).$$

Letting  $v_{q-1} = \sum_{j \leq q-1} (2j-q)u_j$  or, equivalently,  $u_{q-1} = \frac{1}{q-2}v_{q-1} - \sum_{j \leq q-2} \frac{2j-q}{q-2}u_j$ , we can eliminate  $u_{q-1}$ , getting

$$\begin{aligned} P_q(u_1, u_2, \dots, u_{q-2}, u_{q-1}, u_q) &= \sum_{j < k \leq q-2} (q-2k+2j)Q(u_j, u_k) \\ &+ Q\left(\sum_{j \leq q-2} (-q+2+2j)u_j, \frac{1}{q-2}v_{q-1} - \sum_{k \leq q-2} \frac{2k-q}{q-2}u_k\right) \\ &+ Q(v_{q-1}, u_q) =: (*). \end{aligned}$$

To conclude the computations, we insert the change concerning  $v_q$  in (2.9) and we get

$$\begin{aligned}
 (*) &= Q\left(v_{q-1}, u_q - \sum_{j \leq q-2} \frac{-q+2+2j}{q-2} u_j\right) \\
 &\quad + \sum_{j < k \leq q-2} \left\{ (q-2k+2j) + \frac{(-q+2+2j)(q-2k)}{q-2} \right. \\
 &\quad \quad \quad \left. + \frac{(2j-q)(-q+2+2k)}{q-2} \right\} Q(u_j, u_k) \\
 &= Q(v_{q-1}, v_q) + \frac{q}{q-2} \sum_{j < k \leq q-2} (q-2-2(k-j)) Q(u_j, u_k),
 \end{aligned}$$

which is the desired identity (2.10).

Next, we iterate formula (2.10) starting from  $q = 2\ell + 1$  and we get

$$\begin{aligned}
 P_{2\ell+1}(u_1, \dots, u_{2\ell+1}) &= Q(v_{2\ell}, v_{2\ell+1}) + \frac{2\ell+1}{2\ell-1} P_{2\ell-1}(u_1, \dots, u_{2\ell-1}) \\
 &= Q(v_{2\ell}, v_{2\ell+1}) + \frac{2\ell+1}{2\ell-1} Q(v_{2\ell-2}, v_{2\ell-1}) \\
 &\quad + \frac{2\ell+1}{2\ell-3} P_{2\ell-3}(u_1, \dots, u_{2\ell-3}) \\
 &\quad \vdots
 \end{aligned}$$

When we encounter  $P_q$  with  $q = 3$ , the iteration stops, because

$$\begin{aligned}
 P_3(u_1, u_2, u_3) &:= Q(u_1, u_2) + Q(u_2, u_3) - Q(u_1, u_3) \\
 &= Q(u_1, u_2) + Q(u_2 - u_1, u_3) \\
 &= Q(u_1, u_1 + v_2) + Q(v_2, u_3) \\
 &= Q(v_2, u_3 - u_1) = Q(v_2, v_3),
 \end{aligned}$$

where we applied the change of variable  $v_2 = u_2 - u_1$  and  $v_3 = u_3 - u_1$ , i.e. (2.9) with  $q = 3$ .

The proof of the lemma is concluded.  $\square$

**2.2. Construction of the inner cone.** Using the higher-order envelopes of Theorem 2.1, we can prove Theorem 1.1.

Let  $C \subset \mathbb{G}$  be a convex set in a two-step Carnot group  $\mathbb{G}$  with group law (1.1). Let us consider a point  $(z, t) \in C$ . Let  $p \in \mathbb{N}$  be such that

Theorem 2.1 holds true. Assume that there is  $\xi \in Z$  such that  $(z, t) \cdot (p\xi, 0) \in \text{int}(C)$ . Define the map

$$\Gamma_{(\zeta, \tau)}(u_1, \dots, u_p) := e^{u_p \cdot X} \dots e^{u_2 \cdot X} \cdot e^{u_1 \cdot X}(\zeta, \tau),$$

for all  $u_1, \dots, u_p, \zeta \in Z$  and  $\tau \in T$ .

Note that  $\Gamma_{(\zeta, \tau)}(u_1, \dots, u_p) = (\zeta, \tau) \cdot \Gamma(u_1, \dots, u_p)$ , where  $\Gamma = \Gamma_{(0,0)}$  is the function introduced in (2.1). In this notation, if the hypotheses of Theorem 1.1 are satisfied, we have  $(\zeta, \tau) := \Gamma_{(z,t)}(\xi, \dots, \xi) \in \text{int}(C)$ . Let  $\varrho > 0$  so that  $B((\zeta, \tau), \varrho) \subset \text{int}(C)$ , where  $B$  denotes the subRiemannian ball. By continuity, there is  $\varepsilon > 0$  and  $\varrho > 0$  so that, if  $(z', t') \in \mathbb{G}$  and  $u_1, \dots, u_p \in Z$  satisfy

$$(2.11) \quad \max\{d((z, t), (z', t')), |u_j - \xi| : j = 1, \dots, p\} < \varepsilon,$$

then

$$(2.12) \quad \Gamma_{(z', t')}(u_1, \dots, u_p) \in B((\zeta, \tau), \varrho) \subset \text{int}(C).$$

We are now ready to give the proof of Theorem 1.1.

*Proof of Theorem 1.1: Step 1:* If  $u_j$  and  $(z', t')$  satisfy (2.11), then for all nonnegative  $\lambda_1, \dots, \lambda_p$  satisfying  $\sum \lambda_j = p$ , we have

$$\Gamma_{(z', t')}(\lambda_1 u_1, \dots, \lambda_p u_p) \in B((\zeta, \tau), \varrho) \subset \text{int}(C).$$

To prove Step 1, denote by  $d$  the subRiemannian distance and write first

$$\begin{aligned} d(\Gamma_{(z', t')}(\lambda_1 u_1, \dots, \lambda_p u_p), \Gamma_{(z, t)}(\xi, \dots, \xi)) \\ \leq d(\Gamma_{(z', t')}(\lambda_1 u_1, \dots, \lambda_p u_p), \Gamma_{(z', t')}(\xi, \dots, \xi)) \\ + d(\Gamma_{(z', t')}(\xi, \dots, \xi), \Gamma_{(z, t)}(\xi, \dots, \xi)). \end{aligned}$$

The second term can be estimated with  $C(|z' - z|^{1/2} + |t' - t|^{1/2})$ , where the constant  $C$  depends on the vector  $\xi$  appearing in the hypotheses.

By the explicit form of  $\Gamma$  described in (2.1), the first term can be evaluated by

$$\begin{aligned} d((z', t') \cdot \Gamma_{(0,0)}(\lambda_1 u_1, \dots, \lambda_p u_p), (z', t') \cdot \Gamma_{(0,0)}(\lambda_1 \xi, \dots, \lambda_p \xi)) \\ \leq C \left( \left| \sum_{j \leq p} \lambda_j u_j - p\xi \right| + \left| \sum_{j < k \leq p} Q(\lambda_j u_j, \lambda_k u_k) + Q\left(\sum_{j \leq p} \lambda_j u_j, p\xi\right) \right|^{1/2} \right). \end{aligned}$$

The first term is easily evaluated:  $|\sum_j \lambda_j u_j - p\xi| = |\sum_j \lambda_j (u_j - \xi)| \leq C\varepsilon$ . To estimate the term under the square root, note that for all  $j, k$  we have by homogeneity and anti-symmetry,

$$|Q(\lambda_j u_j, \lambda_k u_k)| = \lambda_j \lambda_k |Q(u_j - u_k, u_k)| \leq C|\xi| \cdot |u_j - u_k| \leq C\varepsilon.$$

An analogous argument can be used to see that  $|Q(\sum_{j \leq p} \lambda_j u_j, p\xi)| \leq C\varepsilon$ , where again the constant  $C$  depends on  $|\xi|$ , which is fixed. This ends Step 1.

*Step 2:* We show that for all  $\lambda \in ]0, 1[$  and for all  $(z', t') \in C$  with  $d((z, t), (z', t')) < \varepsilon$  we have

$$\{\Gamma_{(z', t')}(\lambda u_1, \lambda u_2, \dots, \lambda u_p) : |u_j - \xi| < \varepsilon \text{ for all } j \in \{1, \dots, p\}\} \subset C.$$

To accomplish Step 2 observe first that if  $|u_1 - \xi| < \varepsilon$ , then by applying Step 1 with  $\lambda_1 = p$  and  $\lambda_j = 0$  for  $j \geq 2$ , we obtain  $\Gamma_{(z', t')}(pu_1, 0, 0, \dots, 0) \in C$ . Furthermore, the points  $(z', t')$  and  $\Gamma_{(z', t')}(pu_1, 0, 0, \dots, 0)$  are aligned. Therefore, we have  $\Gamma_{(z', t')}(\lambda u_1, 0, 0, \dots, 0) \in C$  for all  $\lambda \in [0, p]$ .

Next, consider the aligned points

$$\Gamma_{(z', t')}(\lambda u_1, 0, \dots, 0) \in C \quad \text{and} \quad \Gamma_{(z', t')}(\lambda u_1, (p - \lambda)u_2, 0, \dots, 0).$$

The second point belongs to  $C$  by Step 1 and the points are aligned. Then

$$\Gamma_{(z', t')}(\lambda u_1, su_2, 0, \dots, 0) \in C \quad \text{for all } s \in [0, p - \lambda].$$

Taking  $s = \lambda$  (which is less than  $p - \lambda$ ), we discover that

$$\Gamma_{(z', t')}(\lambda u_1, \lambda u_2, 0, \dots, 0) \in C.$$

Iterating the procedure, we easily conclude that  $\Gamma_{(z', t')}(\lambda u_1, \dots, \lambda u_p) \in C$  and Step 2 is proved.

*Step 3:* We finalize the construction of the inner cone. Let  $(z', t') \in C$  with  $d((z, t), (z', t')) < \varepsilon$ . Then, by the previous steps, denoting by  $\ell_{(z', t')}$  the left translation,

$$\begin{aligned} C &\supseteq \{\Gamma_{(z', t')}(\lambda u_1, \lambda u_2, \dots, \lambda u_p) : |u_j - \xi| < \varepsilon \forall j, \lambda \in [0, 1]\} \\ &= \ell_{(z', t')} \{\Gamma(\lambda u_1, \dots, \lambda u_p) : |u_j - \xi| < \varepsilon \forall j, \lambda \in [0, 1]\} \\ &= \ell_{(z', t')} \bigcup_{\lambda \in [0, 1]} \delta_\lambda \{\Gamma(u_1, \dots, u_p) : |u_j - \xi| < \varepsilon\} \quad (\text{by Theorem 2.1}) \\ &\supseteq \ell_{(z', t')} \bigcup_{\lambda \in [0, 1]} \delta_\lambda B_{\text{Euc}}(\Gamma(\xi, \dots, \xi), c_0 \varepsilon^2). \end{aligned}$$

Let  $\delta_0$  be such that  $B_{\text{Euc}}(\Gamma(\xi, \dots, \xi), c_0 \varepsilon^2) \supseteq B(\Gamma(\xi, \dots, \xi), \delta_0)$ . Therefore

$$C \supseteq \ell_{(z', t')} \bigcup_{\lambda \in [0, 1]} B(\Gamma(\lambda \xi, \dots, \lambda \xi), \delta_0 \lambda) = \bigcup_{\lambda \in [0, 1]} B((z', t') \cdot (\lambda p \xi, 0), \delta_0 \lambda),$$

which is the required open cone. We used the known properties of the natural dilations  $\delta_\lambda(z, t) := (\lambda z, \lambda^2 t)$  of the group. Namely,  $\delta_\lambda B((z, t), r) = B(\delta_\lambda(z, t), \lambda r)$ , for all  $\lambda, r > 0$  and  $(z, t) \in \mathbb{G}$ .

*Step 4:* In Step 3, we worked with  $(z', t') \in C$ . To conclude, it suffices to approximate any desired  $(z', t') \in \partial C$  with  $d((z', t'), (z, t)) < \varepsilon$  with a sequence  $(z'_n, t'_n) \in C$  and the construction is finished.  $\square$

### 3. Monotone sets and intrinsic epigraphs

We start this section giving the details of the proof of the classification of monotone sets in Métivier groups. The idea of the proof has been suggested to the author by R. Monti. All we need to know about Métivier groups is that for all  $P, Q$  in such kind of group  $\mathbb{G}$  there is a Heisenberg subgroup  $H_0$  such that the left coset  $P \cdot H_0 =: H$  contains  $P$  and  $Q$ . A Heisenberg subgroup of a two-step group (1.1) is a three-dimensional subgroup of the form  $\{\text{span}\{(z, 0), (\zeta, 0), (0, Q(z, \zeta))\}\}$  for suitable independent  $z, \zeta \in Z$  with  $Q(z, \zeta) \neq 0$ . See [2, 15] and the references therein.

**Proposition 3.1.** *If  $E \subset \mathbb{G}$  is a precisely monotone set in a Métivier group, and if  $\partial E \subsetneq \mathbb{G}$ , then either  $E \in \{\emptyset, \mathbb{G}\}$ , or there is an open half-space  $\Omega \subset \mathbb{G}$  such that  $\Omega \subset E \subset \overline{\Omega}$ .*

Unfortunately we cannot claim that horizontally monotone sets in  $\mathbb{H}$  are Euclidean convex. Otherwise, the proof below would be much shorter. As an example, consider the set

$$E = \{(x, y, t) \in \mathbb{H} : y \geq 0\} \setminus \{(x, 0, 0) : x \in \mathbb{R}\}.$$

*Proof:* Since  $\partial E \neq \mathbb{G}$ , we have  $\text{int}(E) \cup \text{int}(E^c) \neq \emptyset$ . Assume that  $\text{int}(E) \neq \emptyset$ . Write then  $\mathbb{G}$  as a disjoint union of  $\text{int}(E)$  and  $(\text{int}(E))^c = \overline{E^c}$ . To prove the proposition, it suffices to show that the nonempty set  $\text{int}(E)$  is monotone in the Euclidean sense, i.e. that  $\text{int}(E)$  and  $\overline{E^c}$  are convex in the usual sense. Then the statement follows from [7, Lemma 4.2].

Let  $P$  and  $Q \in \text{int}(E)$  and let  $H := P \cdot H_0$  be the left coset of a Heisenberg subgroup of  $\mathbb{G}$  containing  $P$  and  $Q$  (which always exists, by the Métivier condition). Note that  $E \cap H$  is monotone in  $H$  and also  $P, Q \in \text{int}_H(E \cap H)$ , which, by [7], is either an open halfspace in  $H$  or the whole  $H$  (we denote by  $\text{int}_H$  the interior in the induced topology on  $H$ ). Thus, the segment  $[P, Q]$  which connects  $P$  and  $Q$  is contained in  $\text{int}_H(E \cap H)$  which is a subset of  $E \cap H$ . Finally,  $[P, Q] \subset E$ . Applying the same argument to all pairs  $P', Q'$  with  $P'$  close to  $P$  and  $Q' = P' + (Q - P)$  close to  $Q$ , we conclude that  $[P, Q] \subset \text{int}(E)$ .

In order to show that the set  $\overline{E^c}$  is Euclidean convex, assume that it is nonempty, consider  $P, Q \in E^c$  and let  $H = P \cdot H_0$  be the left coset of a Heisenberg subgroup  $H_0$  such that  $P, Q \in H$ . Since  $E^c \cap H$  is monotone in  $H$ , either  $E^c \cap H = H$  and then  $[P, Q] \subset E^c$ , or else there is an open three-dimensional half-space  $\Sigma_0 \subset H$  such that  $\Sigma_0 \subset E^c \cap H \subset \Sigma$ , where  $\Sigma \subset H$  is the closure of  $\Sigma_0$ . Taking the closure in  $\mathbb{G}$ , we see that  $\overline{E^c} \cap H = \Sigma$  and thus  $[P, Q] \subset \Sigma = \overline{E^c} \cap H \subset \overline{E^c}$ .

Finally, approximating points  $P, Q \in \overline{E^c}$  with points in  $E^c$ , we get the Euclidean convexity of  $\overline{E^c}$ .  $\square$

Next we prove that if  $\mathbb{G} = \mathbb{H}^n$  is the  $n$ -th Heisenberg group, then the assumption  $\partial E \subsetneq \mathbb{H}^n$  can be removed. Recall that  $\mathbb{H}^n = \mathbb{C}^n \times \mathbb{R} \ni (z, t) = (x + iy, t)$  is equipped with the law

$$(z, t) \cdot (\zeta, \tau) := (z + \zeta, t + \tau + 2\Im(z \cdot \bar{\zeta})) \quad \text{with } z \cdot \bar{\zeta} = \sum_{j=1}^n z_j \bar{\zeta}_j.$$

**Proposition 3.2.** *If  $E \subset \mathbb{H}^n$  is a precisely monotone set and  $\emptyset \neq E \neq \mathbb{H}^n$ , then there is an open half-space  $\Omega \subset \mathbb{H}^n$  such that  $\Omega \subset E \subset \overline{\Omega}$ .*

Unfortunately we have not been able to get a generalization of Proposition 3.2 to the setting of Métivier groups.

*Proof: Step 1:* We start by observing that given two open half-spaces  $\Omega_1$  and  $\Omega_2 \subset \mathbb{H}^n$  there are two aligned points  $P_1 \in \Omega_1$  and  $P_2 \in \Omega_2$ . To check this remark, if  $\partial\Omega_1$  is a characteristic plane, then, by left-translating its characteristic point to the origin, we must have  $\Omega_1 = \{t > 0\}$  or  $\Omega_1 = \{t < 0\}$ ; if instead  $\Omega_1$  has noncharacteristic boundary, up to left translations it must be defined by the inequality  $\sum_{j=1}^n (a_j x_j + b_j y_j) > 0$  for some  $a, b \in \mathbb{R}^n$ .

In the first case, starting from the fact that any halfspace  $\Omega_2$  must contain points  $(\zeta, \tau)$  with  $\zeta \neq 0$ , it is easy to check that any such point  $(\zeta, \tau)$  is aligned with a point with  $t > 0$ . This can be seen by considering the parametrized line  $\gamma(s) := (\zeta, \tau) \cdot (-is\zeta, 0)$  with large  $s$ . The case  $\Omega_1 = \{t < 0\}$  is analogous (look at  $\gamma(s)$  with  $s \rightarrow -\infty$ ).

In the second case, given  $(\zeta, \tau) \in \mathbb{H}$ , all points  $\gamma(s) = (\xi + i\eta, \tau) \cdot (s(0 + ib), 0)$  belong to  $\Omega_1$  for large  $s$ .

*Step 2:* We show the proposition by induction. Assume that the statement is true for  $\mathbb{H}^n$ . Let  $E \subset \mathbb{H}^{n+1}$  be a monotone set with  $\emptyset \neq E \neq \mathbb{H}^{n+1}$ . By Proposition 3.2, it suffices to show that  $\text{int}(E) \cup \text{int}(E^c) \neq \emptyset$ .

For any  $z_{n+1} \in \mathbb{C}$  define  $H_{z_{n+1}} := \{(z, z_{n+1}, t) : (z, t) \in \mathbb{H}^n\}$ . Assume without loss of generality that  $H_0 \cap E \neq \emptyset$ . Then by the inductive hypothesis  $H_0 \cap E \supset \Omega_1$ , where  $\Omega_1$  is an open  $(2n+1)$ -dimensional halfspace contained in  $H_0$ .

There are two cases. Either for all  $\zeta_{n+1} \in \mathbb{C} \setminus \{0\}$  we have  $H_{\zeta_{n+1}} \subset E^c$ , or there is  $\zeta_{n+1} \in \mathbb{C} \setminus \{0\}$  such that the set  $E \cap H_{\zeta_{n+1}}$  is nonempty and then it contains an open half space  $\Omega_2 \subset H_{\zeta_{n+1}}$ . The first case cannot occur, because it would contradict the convexity of  $E^c$ . Then we are always in the second case. By Step 1 we can find two points  $P = (z, 0, t) \in \Omega_1$  and  $Q = (z, 0, t) \cdot (\zeta, \zeta_{n+1}, 0) \in \Omega_2$ , where both  $\Omega_1$  and  $\Omega_2$  are contained in  $E$ .

Let  $\gamma(s) = (z, 0, t) \cdot (s\zeta, s\zeta_{n+1}, 0)$  be a line connecting  $\gamma(0) = P$  and  $\gamma(1) = Q$ . By convexity,  $\gamma[0, 1] \subset E$ . Taking points  $(z', t')$  in a small open neighborhood  $G$  of  $(z, t) \in \mathbb{H}_n$ , we can consider the set of points

$$\Sigma := \{(z', 0, t') \cdot (s\zeta, s\zeta_{n+1}, 0) : s \in ]0, 1[, (z', t') \in G\}.$$

It is clear that  $\Sigma$  is a  $(2n+2)$ -dimensional surface contained in  $E$ . Take any point  $P \in \mathbb{H}^{n+1}$  aligned with some point  $Q \in \Sigma$  and such that the segment  $[P, Q]$  is transversal to  $\Sigma$  at  $Q$ . If  $P \notin \partial E$  we immediately get  $\partial E \neq \mathbb{H}^{n+1}$  and Step 2 is concluded. If instead  $P \in \partial E$ , then we can apply Proposition 3.3 to conclude that  $\text{int}(E) \neq \emptyset$ .  $\square$

Next, we describe some properties of monotone sets in a Carnot group of step two which follow from the results of the previous section.

First of all, as a corollary of Theorem 2.1 we get the generalization of [7, Proposition 4.6] to general two-step Carnot groups.

**Proposition 3.3.** *Let  $E \subset \mathbb{G}$  be a monotone set in a two-step Carnot group  $G$  with law (1.1). If a point  $(z, t)$  belongs to a line  $\ell$  and a surface  $\Sigma \subset E$  intersects transversally the line  $\ell$  at a second different point  $(\zeta, \tau) \in \ell$ , then:*

- (i) *If  $(z, t) \in \overline{E}$ , then the open segment  $] (z, t), (\zeta, \tau) [ \subset \ell$  is contained in  $\text{int}(E)$ .*
- (ii) *If  $(z, t) \in E^c$ , then the connected component of  $\ell \setminus \{(\zeta, \tau)\}$  not containing  $(z, t)$  is contained in  $\text{int}(E)$ .*

The first statement (i) is a qualitative version of Theorem 1.1 (it can also be proved by following the argument in [7, Proposition 4.6] after replacing the map  $Z \times Z \ni (u_1, u_2) \mapsto \Gamma_P(u_1, u_2)$  with the map  $Z^p \ni (u_1, \dots, u_p) \mapsto \Gamma_P(u_1, \dots, u_p)$  and using Theorem 2.1). As observed in [7], statement (i) holds for a merely convex set  $E$ . The second statement can be proved easily arguing as in [7].

Then we have the generalization of the following lemmas to Step 2 Carnot groups  $(\mathbb{G}, \cdot)$ .

**Lemma 3.4** (compare with [7, Lemma 4.8]). *Let  $(\mathbb{G}, \cdot)$  be a two-step Carnot group with law (1.1). If  $(z, t)$  and  $(\zeta, \tau) = (z, t) \cdot (u, 0) \in \mathbb{G}$  are aligned points with  $u \neq 0$  and both belong to  $\partial E$ , then the whole line connecting them is contained in  $\partial E$ .*

*Proof:* The same as in [7].  $\square$

The following lemma should be compared with its Heisenberg version, Lemma 4.9 in [7].

**Lemma 3.5.** *If  $E \subset \mathbb{G}$  is monotone, then for all  $P \in \partial E$  and for any two-dimensional subspace  $V \subset H_P$ , there is a line  $\ell$  satisfying  $P \in \ell \subset V$  and completely contained in  $\partial E$ .*

*Proof:* Assume without loss of generality that  $(0, 0) \in \partial E$ . Let  $u, v \in Z$  be a pair of orthonormal vectors. Consider the circle  $S := \{(u \cos \theta + v \sin \theta, 0) : \theta \in \mathbb{R}\}$ . We claim that there is a point  $(z, 0) \in S \cap \partial E$ . Indeed, if this would not be true, then either  $S \subset \text{int}(E)$  or else  $S \subset \text{int}(E^c)$ . Assume the former and let  $(z, 0)$  and  $(-z, 0)$  be two points of  $S$ . Since they are aligned, their intermediate point  $(0, 0)$  would be in  $\text{int}(E)$ , giving a contradiction.  $\square$

Next we recall the definition of *intrinsic graph*.

**Definition 3.6** (Intrinsic graph). Let  $\mathbb{G} = Z \times T$  be the two-step Carnot group with law (1.1). Let  $\xi \in Z$  be a unit vector and let  $\xi^\perp \subset Z$  be its orthogonal space. Let  $W \subset \xi^\perp \times T$  be an open set in  $\xi^\perp \times T$  containing the origin. Given a function  $\psi: W \rightarrow \mathbb{R}$ , the  $\xi$ -graph (and the corresponding epigraph and ipograph) associated with  $\psi$  are

$$\begin{aligned} \text{graph}(\psi) &:= \{(\eta, \tau) \cdot (\psi(\eta, \tau)\xi, 0) : (\eta, \tau) \in W\}, \\ \text{epi}(\psi) &:= \{(\eta, \tau) \cdot (s\xi, 0) : (\eta, \tau) \in W \text{ and } \psi(\eta, \tau) < s < \infty\}, \text{ and} \\ \text{ipo}(\psi) &:= \{(\eta, \tau) \cdot (s\xi, 0) : (\eta, \tau) \in W \text{ and } -\infty < s < \psi(\eta, \tau)\}. \end{aligned}$$

Intrinsic graphs appear in Geometric Measure Theory in Carnot groups. See [12, 1, 11] and the references therein.

Here we show the following theorem which relates monotonicity and intrinsic graphs.

**Theorem 3.7.** *Let  $E \subset \mathbb{G}$  be a monotone set in a Carnot group of step two. Assume that  $(0, 0) \in \partial E$ . Then, either the whole horizontal plane  $H_{(0,0)}$  is contained in  $\partial E$ , or else there is a vector  $\xi \in Z$ , an open neighborhood  $W$  of the origin in  $\xi^\perp \times T$ , and a function  $\psi: W \rightarrow \mathbb{R}$  such that*

$$(3.1) \quad \text{epi}(\psi) \subset \text{int}(E), \quad \text{ipo}(\psi) \subset \text{int}(E^c), \quad \text{and} \quad \text{graph}(\psi) \subset \partial E.$$

*Furthermore, the function  $\psi$  is continuous.*



*Remark 3.8.* Observe the following facts.

- (a) The inner cone property is not global. For example, let us consider the standard Heisenberg group  $\mathbb{H} = \{(x, y, t) \in \mathbb{R}^3\}$  with group law  $(x, y, t) \cdot (\xi, \eta, \tau) := (x + \xi, y + \eta, t + \tau + 2(y\xi - x\eta))$ . Consider the monotone set  $E = \{x > t\} \subset \mathbb{H}$ , where  $\mathbb{H}$  denotes the standard Heisenberg group. There is no  $\lambda > 0$  such that the cone  $\cup_{0 < s < \infty} B((s, 0, 0), \lambda s)$  is contained in  $E$ . However, the truncated cone of the form  $\cup_{0 < s \leq s_0} B((s, 0, 0), \lambda s)$  is contained in  $E$  for suitable  $s_0$ .
- (b) The same example as in (a) shows that the set  $W$  appearing in the statement of Theorem 3.7 can be a strict subset of the whole  $\xi^\perp \times T$ . Indeed,  $(0, 1/2, 0) \in \partial E$  becomes a characteristic point, but intrinsic graphs inherently cannot have characteristic points.
- (c) The continuity of  $\psi$  can be strengthened by saying that the function  $\psi$  is *intrinsic Lipschitz-continuous* (see [3]).
- (d) A version of Theorem 3.7 holds assuming that  $E$  is convex, but not necessarily monotone.

*Proof of Theorem 3.7:* Assume that  $(0, 0) \in \partial E$  and  $(s_0 e_1, 0) \in \text{int}(E)$  for some  $s_0 > 0$ . Split  $(x, t) = (x_1, \hat{x}_1, t)$ , where  $\hat{x}_1 = (x_2, x_3, \dots, x_m) \in \mathbb{R}^{m-1}$ . We know by Theorem 1.1 that there is a positive number  $\lambda$  such that

$$\bigcup_{0 < s \leq s_0} B((s, 0, 0), \lambda s) \subset \text{int}(E) \quad \text{and} \quad \bigcup_{0 < s \leq s_0} B((-s, 0, 0), \lambda s) \subset \text{int}(E^c).$$

*Step 1:* To show the existence of the function  $\psi$ , we prove that there is a neighborhood  $W \subset \mathbb{R}^{m-1} \times \mathbb{R}^k$  such that, for all  $(\hat{x}_1, t) \in W$ , we have

$$(3.2) \quad (0, \hat{x}_1, t) \cdot (-s_0, 0, 0) \in \text{int}(E^c) \quad \text{and} \quad (0, \hat{x}_1, t) \cdot (s_0, 0, 0) \in \text{int}(E).$$

To prove the second inclusion, since  $(0, \hat{x}_1, t) \cdot (s_0, 0, 0) = (s_0, \hat{x}_1, t + Q((0, \hat{x}_1), (s_0, 0)))$ , it suffices to see that there is  $W$  such that for all  $(\hat{x}_1, t) \in W$ , we have

$$\|(-s_0, 0, 0) \cdot (s_0, \hat{x}_1, t + Q((0, \hat{x}_1), (s_0, 0)))\| < \lambda s_0,$$

which, using the group law and the skew-symmetry of  $Q$ , becomes

$$\|(0, \hat{x}_1, t + 2Q((0, \hat{x}_1), (s_0, 0)))\| < \lambda s_0,$$

which clearly holds true provided that  $|\hat{x}_1|$  and  $|t|$  are small. We use the estimate  $C|\hat{x}_1| + C|t|^{1/2} + C|\hat{x}_1|^{1/2}s_0^{1/2} < \lambda s_0$  if  $\hat{x}_1$  and  $|t|$  are small enough.

The proof of the first inclusion in (3.2) is analogous.

*Step 2:* By the properties of monotone sets we can conclude that, for all  $(\hat{x}_1, t) \in W$ , there is a number  $\psi(\hat{x}_1, t)$  such that the point  $(0, \hat{x}_1, t) \cdot (s, 0, 0)$  belongs to  $\text{int}(E)$  if  $s \in ]\psi(\hat{x}_1, t), +\infty[$ , to  $\text{int}(E^c)$  if  $s \in ]-\infty, \psi(\hat{x}_1, t)[$ , and to  $\partial E$  if  $s = \psi(\hat{x}_1, t)$ . The function  $\psi$  is also bounded by  $s_0$  on  $W$ .

*Step 3:* We show the continuity of  $\psi$  at the origin (at other points the same argument works). Namely, we prove that

$$\limsup_{(\hat{x}_1, t) \rightarrow (0, 0) \in \mathbb{R}^{m-1} \times \mathbb{R}^k} \psi(\hat{x}_1, t) \leq 0.$$

The proof of  $\liminf \geq 0$  is analogous. Assume that  $\lim_{n \rightarrow \infty} \psi(\hat{x}_1^n, t^n) = L > 0$  for some sequence  $(\hat{x}_1^n, t^n) \rightarrow (0, 0)$ . Then  $(0, \hat{x}_1^n, t^n) \cdot (\psi(\hat{x}_1^n, t^n), 0, 0) \in \partial E$  for all  $n \in \mathbb{N}$ . The sequence converges to  $(0, 0, L)$  and this contradicts the fact that the half line  $\{(s, 0, 0) : 0 < s < \infty\}$  is contained in  $\text{int}(E)$ .  $\square$

#### 4. Monotone sets in $\mathbb{H} \times \mathbb{R}$

Consider the direct product  $\mathbb{H}_{(x,y,t)} \times \mathbb{R}_u = \mathbb{R}^4$  with law

$$(x, y, u, t) \cdot (x', y', u', t') = (x + x', y + y', u + u', t + t' + 2(yx' - y'x))$$

and with the subRiemannian distance defined by the vector fields  $X = \partial_x + 2y\partial_t$ ,  $Y = \partial_y - 2x\partial_t$ , and  $U = \partial_u$ .

In this section we show the following result, which gives immediately the proof of Theorem 1.2:

**Theorem 4.1.** *Let  $E \subset \mathbb{H} \times \mathbb{R}$  be a monotone set, with  $\emptyset \neq E \neq \mathbb{R}^4$ . Then  $\partial E$  is a hyperplane.*

It is reasonable to think that a similar statement can be proved in the setting  $\mathbb{H} \times \mathbb{R}^k$ , with  $k > 1$ . However, this would require some generalizations of the proof which are not completely trivial.

##### 4.1. Preparatory lemmas on monotone sets in $\mathbb{H} \times \mathbb{R}$ .

**Lemma 4.2.**  *$\text{int}(E) \cup \text{int}(E^c)$  is a nonempty (open) dense set or, equivalently,  $\text{int}(\partial E) = \emptyset$ .*

*Proof:* Assume that  $\partial E \supset \Omega := \{(x, y, u, t) \in \mathbb{R}^4 : |x|, |y|, |u|, |t| < \varepsilon\}$  for some  $\varepsilon > 0$ . We consider the set  $\Sigma := \{(x, y, 0, t) : (x, y, t) \in \mathbb{R}^3 \simeq \mathbb{H}\}$ . Since  $E \cap \Sigma$  is monotone in  $\mathbb{H}$ , by [7], there is  $(z, 0, t) \in \Omega$  and there is  $\delta > 0$  such that  $B((z, 0, t), \delta) \cap \Sigma \subset E \cap \Sigma$ , or  $B((z, 0, t), \delta) \cap \Sigma \subset E^c \cap \Sigma$ . Assume the former and consider the point  $(z, \frac{\varepsilon}{2}, t)$  which belongs to  $\partial E$  by assumption. Then, by Proposition 3.3 we see that  $(z, u, t) \in \text{int}(E)$  for all  $u \in ]-\infty, \frac{\varepsilon}{2}]$ . In particular,  $(z, 0, t) \in \text{int}(E)$ . We have a contradiction and the lemma is proved.  $\square$

**Lemma 4.3.** *Let  $P \in \mathbb{H} \times \mathbb{R}$  be such that the horizontal plane  $H_P$  is contained in  $\partial E$ . Then  $\partial E = H_P$ .*

*Proof:* Let  $H_0 = \{(\xi, \eta, \omega, 0) : (\xi, \eta, \omega) \in \mathbb{R}^3\} \subset \partial E$  and assume by contradiction that a point  $P_0 = (x_0, y_0, u_0, t_0)$  with  $t_0 \neq 0$  belongs to  $\partial E$ . Without loss of generality assume  $P_0 = (x_0, y_0, 0, 1)$ .

*Step 1:* We claim that we can assume  $z_0 = (x_0, y_0) \neq (0, 0)$ . Indeed, if  $P_0 = (0, 0, 0, 1) \in \partial E$ , we consider the plane  $S = \{(\zeta, 0, 1) : \zeta \in \mathbb{C}\}$ . By Lemma 3.5, there is a line  $\{(sz_0, 0, 1) : s \in \mathbb{R}\} \subset \partial E$ . Up to a rotation we may assume  $z_0 = (x_0, 0)$ .

*Step 2:* Next we show that  $\partial E = \mathbb{G}$ , contradicting Lemma 4.2. To see the claim, let  $P_0 = (x_0, 0, 0, 1) \in \partial E$  be the point found in Step 1. Consider the family of lines

$$\gamma_{\xi, \eta, \sigma}(\lambda) = (x_0, 0, 0, 1) \cdot (\lambda\xi, \lambda\eta, \lambda\sigma, 0) = (x_0 + \lambda\xi, \lambda\eta, \lambda\sigma, 1 - 2x_0\lambda\eta),$$

where  $(\xi, \eta, \sigma) \in \mathbb{R}^3$  is a nonzero vector. It is easy to see that for any  $\eta \neq 0$  the line  $\gamma_{\xi, \eta, \sigma}$  touches the plane  $t = 0$  at a time  $\lambda = \frac{1}{2x_0\eta} \neq 0$ . This implies that the whole line is contained in  $\partial E$ , by Lemma 3.4. Therefore, if we consider the horizontal plane at  $P_0$

$$H_{P_0} = \{(x_0 + \xi, \eta, \sigma, 1 - 2x_0\eta) : \xi, \eta, \sigma \in \mathbb{R}\},$$

all its points with  $\eta \neq 0$  belong to  $\partial E$ . Since  $\partial E$  is closed, we conclude that  $H_{P_0} \subset \partial E$ . Note that  $H_{P_0}$  is the plane of equation  $t = 1 - 2x_0y$ . Now, consider the family of lines

$$\gamma_{x, y, u}(\lambda) = (x, y, u, 0) \cdot (\lambda, 0, 0, 0) = (x + \lambda, y, u, 2\lambda y),$$

where  $(x, y, u) \in \mathbb{R}^3$ . It is easy to see that if  $\frac{1}{2x_0} \neq y \neq 0$ , the line  $\gamma_{x, y, u}$  touches the plane  $H_{P_0}$  at the nonzero time  $\lambda = \frac{1-2x_0y}{2y} \neq 0$ . By Lemma 3.4, this implies that for all  $(x, y, u)$  with  $0 \neq y \neq \frac{1}{2x_0}$ , the line  $\gamma_{x, y, u}$  is fully contained into  $\partial E$ . Thus

$$\begin{aligned} \partial E &\supset \left\{ (x + \lambda, y, u, 2\lambda y) : x, u, \lambda \in \mathbb{R} \quad \text{and} \quad 0 \neq y \neq \frac{1}{2x_0} \right\} \\ &= \left\{ (\xi, \eta, \omega, \tau) : \xi, \omega, \tau \in \mathbb{R} \quad \text{and} \quad \frac{1}{2x_0} \neq \eta \neq 0 \right\}. \end{aligned}$$

Taking the closure of the latter set, we conclude that  $\partial E = \mathbb{R}^4$ . This contradicts Lemma 4.2 and ends the proof.  $\square$

**Lemma 4.4.** *Let  $E$  and  $E^c$  be both nonempty. Let  $0 \in \partial E$ . Then one, and only one, of the following three items holds.*

- (i)  $\partial E = \{(x, y, u, 0) : (x, y, u) \in \mathbb{R}^3\}$ .
- (ii) The point  $(0, 0, 1, 0)$  belongs to  $\text{int}(E) \cup \text{int}(E^c)$ .
- (iii) There is  $(a, b) \neq (0, 0)$  such that  $(a, b, 0, 0) \in \text{int}(E) \cup \text{int}(E^c)$ .

*Proof:* Assume that  $\{t = 0\} \subset \partial E$ . Then by Lemma 4.3 the boundary  $\partial E$  is the hyperplane of equation  $t = 0$ .

If (i) does not hold, in order to show that at least one among (ii) and (iii) must hold, assume that both such conditions fail. This means that all points  $(x, y, 0, 0)$  and  $(0, 0, u, 0)$  belong to  $\partial E$  for all  $(x, y, u) \in \mathbb{R}^3$ . Therefore,  $\gamma(\lambda) = (\lambda x, \lambda y, (1 - \lambda)u, 0) \in \partial E$  for all  $x, y, u, \lambda \in \mathbb{R}$ . Taking  $\lambda = \frac{1}{2}$  and since  $x, y$ , and  $u$  are arbitrary, we find that (i) holds, getting a contradiction.  $\square$

**4.2. Classification of monotone  $(pX + qY)$ -epigraph.** In this section we analyze case (iii) of Lemma 4.4.

**Proposition 4.5** (Monotone  $(pX + qY)$ -epigraphs (local statement)). *Assume that the origin  $(0, 0, 0, 0) \in \partial E$  and  $(p, q, 0, 0) \in \text{int}(E)$  for some  $(p, q) \neq (0, 0)$ . Then there is a neighborhood  $\Omega$  of 0 in  $\mathbb{R}^4$  and a plane  $\Sigma$  of equation  $rx + sy + cu + bt = 0$ , not containing  $(p, q, 0, 0)$  such that  $\partial E \cap \Omega = \Sigma \cap \Omega$ .*

If we assume without loss of generality that  $(p, q) = (1, 0)$ , then the plane has equation  $x = ay + cu + bt$ .

*Proof of Proposition 4.5:* Without loss of generality, we prove the statement in the case  $(p, q) = (1, 0)$ . By Theorem 3.7, there is a neighborhood  $U$  of  $(0, 0, 0) \in \mathbb{R}^3$  and a continuous function  $\psi(y, u, t)$  on  $U$  such that

$$(4.1) \quad \begin{aligned} \text{epi}(\psi) &:= \{(s, y, u, t + 2ys) : (y, u, t) \in U, s > \psi(y, u, t)\} \subset \text{int}(E), \\ \text{graph}(\psi) &:= \{(\psi(y, u, t), y, u, t + 2y\psi(y, u, t)) : (y, u, t) \in U\} \subset \partial E, \end{aligned}$$

while  $\text{ipo}(\psi) \subset \text{int}(E^c)$ .

*Step 1:* We show that for all  $\bar{u}$  close to 0, there are  $a(\bar{u})$ ,  $b(\bar{u})$ , and  $c(\bar{u}) \in \mathbb{R}$  so that

$$(4.2) \quad \begin{aligned} \{(x, y, \bar{u}, t) \in \mathbb{R}^4 : x > a(\bar{u})y + b(\bar{u})t + c(\bar{u})\} &\subset (E \cap \{u = \bar{u}\}) \\ &\subset \{(x, y, \bar{u}, t) \in \mathbb{R}^4 : x \geq a(\bar{u})y + b(\bar{u})t + c(\bar{u})\}. \end{aligned}$$

To accomplish the step, note that for all  $\bar{u} \in \mathbb{R}$ , the intersection  $E \cap \{u = \bar{u}\}$  is a monotone set in  $\mathbb{G}_{\bar{u}} := \{(x, y, \bar{u}, t) : (x, y, t) \in \mathbb{C} \times \mathbb{R}\}$ , which

is isomorphic to the Heisenberg group. Thus, by [7], there are two open half-spaces  $\Sigma^+, \Sigma^- \subset \mathbb{G}_{\bar{u}}$  with a Euclidean plane as a common boundary and such that  $E \cap \mathbb{G}_{\bar{u}} \supset \Sigma^+$  and  $E^c \cap \mathbb{G}_{\bar{u}} \supset \Sigma^-$ . It suffices to show that such half-spaces have the form (4.2). Let  $m(\bar{u})x + a(\bar{u})y + b(\bar{u})t + c(\bar{u}) < 0$  be the inequality defining  $\Sigma^+$ , where  $m(\bar{u})$ ,  $a(\bar{u})$ ,  $b(\bar{u})$ , and  $c(\bar{u})$  are real numbers. By (4.1), for small  $\bar{u}$ , we know that  $(\psi(0, \bar{u}, 0) + 1, 0, \bar{u}, 0) \in \text{int}(E)$  and  $(\psi(0, \bar{u}, 0) - 1, 0, \bar{u}, 0) \in \text{int}(E^c)$ . Therefore, we may assume  $m(\bar{u}) = -1$  and Step 1 follows.

*Step 2:* After some easy computations (omitted) involving comparison between (4.2) and (4.1), we get

$$(4.3) \quad \psi(y, u, t) = \frac{a(u)y + b(u)t + c(u)}{1 - 2b(u)y} \quad \text{for all } y, u, t \text{ close to } 0.$$

*Step 3:* The function  $\psi$  is affine in  $u$ , i.e. has the form

$$(4.4) \quad \psi(y, u, t) = A(y, t) + B(y, t)u \quad \text{for all } u, y, t \text{ close to the origin,}$$

where  $A(y, t)$  and  $B(y, t)$  are suitable functions.

To show this claim, let us consider a pair of aligned points

$$P = (\psi(y, u, t), y, u, t + 2y\psi(y, u, t)) \quad \text{and}$$

$$Q = (\psi(y, v, t), y, v, t + 2y\psi(y, v, t)) \in \partial E,$$

and also the following line containing  $P$  and  $Q$ :

$$\begin{aligned} \gamma(\sigma) &= (\psi(y, u, t), y, u, t + 2y\psi(y, u, t)) \\ &\quad \cdot (\sigma[\psi(y, v, t) - \psi(y, u, t)], 0, \sigma(v - u), 0) \\ &= (\psi(y, u, t) + \sigma[\psi(y, v, t) - \psi(y, u, t)], y, u + \sigma(v - u), \\ &\quad t + 2y\psi(y, u, t) + 2\sigma y[\psi(y, v, t) - \psi(y, u, t)]). \end{aligned}$$

Since  $\gamma(0)$  and  $\gamma(1)$  belong to  $\partial E$ , we have  $\gamma(\sigma) \in \partial E$  for all  $\sigma \in \mathbb{R}$ , by Lemma 3.4. If  $\sigma \in [0, 1]$  we are close to the origin, and the point  $\gamma(\sigma)$  should belong to the graph of  $\psi$ . Then it must have the form  $(\psi(y', u', t'), y', u', t' + 2y'\psi(y', u', t'))$ . Comparing the four coordinates, we find that  $y' = y$ ,  $u' = u + \sigma(v - u)$ ,  $t' = t$ , and ultimately

$$\psi(y, u + \sigma(v - u), t) = \psi(y, u, t) + \sigma(\psi(y, v, t) - \psi(y, u, t)),$$

for all  $y, u, v, t$  close to the origin and  $\sigma \in [0, 1]$ . Thus  $\psi$  is affine in  $u$  for all fixed  $y, t$ .

*Step 4:* We show that for suitable  $a_0, a_1, b, c \in \mathbb{R}$ , we have

$$(4.5) \quad \psi(y, u, t) = \frac{(a_0 + a_1 u)y + bt + cu}{1 - 2by}.$$

To show such statement, start from the identity

$$(4.6) \quad \psi(y, u, t) = \frac{a(u)y + b(u)t + c(u)}{1 - 2b(u)y} = A(y, t) + B(y, t)u$$

for all  $y, t, u$  close to the origin. By linearity in  $t$ , we can write the right-hand side as  $f_0(y) + f_1(y)t + [g_0(y) + g_1(y)t]u$ , where the form of  $f_0, f_1$  can be obtained letting  $u = 0$  in (4.6). Ultimately, (4.6) can be written in the form

$$(4.7) \quad \psi(y, u, t) = \frac{a(u)y + b(u)t + c(u)}{1 - 2b(u)y} = \frac{a_0y + b_0t}{1 - 2b_0y} + [g_0(y) + g_1(y)t]u,$$

where  $a_0 := a(0)$ ,  $b_0 := b(0)$ , and  $c_0 := c(0) = \psi(0, 0, 0) = 0$ . Letting  $y = t = 0$  we get that  $c(u) = g_0(0)u =: c_1u$ . Evaluating also at  $y = 0$ , we find

$$b(u)t + c_1u = b_0t + [c_1 + g_1(0)t]u =: b_0t + c_1u + b_1tu.$$

Therefore,

$$(4.8) \quad \begin{aligned} \psi(y, u, t) &= \frac{a(u)y + (b_0 + b_1u)t + c_1u}{1 - 2(b_0 + b_1u)y} \\ &= \frac{a_0y + b_0t}{1 - 2b_0y} + [c_1 + \widehat{g}_0(y) + (b_1 + \widehat{g}_1(y))t]u, \end{aligned}$$

where  $\widehat{g}_0$  and  $\widehat{g}_1$  are suitable functions vanishing at 0. Evaluating the identity (4.8) at  $t = 0$  and comparing the left-hand side and the right-hand side, we see that  $\widehat{g}_0$  is rational and smooth at 0. Expanding at the first order in  $y$  both sides with  $t = 0$ , we see that

$$(c_1u + a(u)y)(1 + 2(b_0 + b_1u)y + O(y^2)) = a_0y + O(y^2) + c_1u + \widehat{g}_0'(0)uy,$$

which gives the identity  $a(u) + 2c_1(b_0 + b_1u)u = a_0 + \widehat{g}_0'(0)u$ . Therefore, we conclude that  $a(u) =: a_0 + a_1u - 2b_1c_1u^2$ . Thus the left-hand side of (4.8) becomes

$$\psi(y, u, t) = \frac{(a_0 + a_1u - 2b_1c_1u^2)y + (b_0 + b_1u)t + c_1u}{1 - 2(b_0 + b_1u)y}.$$

But by Step 3 (see the right-hand side of (4.8)),  $\psi$  should be affine in  $u$  for all fixed  $y, t$ . This forces  $b_1 = 0$ , completing the proof of Step 4.

*Step 5:* We show that the coefficient  $a_1$  in (4.5) must vanish. To see that, we first observe (omitted computations) that locally the  $X$ -graph of  $\psi$  agrees with the Euclidean graph

$$\Sigma := \{((a_0 + a_1u)y + cu + bt, y, u, t) : y, u, t \text{ close to the origin}\}.$$

Let  $y \neq 0$  be a number close to 0 and consider a pair of points  $P = ((a_0 + c)y + a_1 y^2, y, y, 0)$  and  $Q = (-(a_0 + c)y + a_1 y^2 + bt, -y, -y, t)$ . Both  $P$  and  $Q$  belong to  $\Sigma$ . A computation shows that the choice  $t = \frac{4a_1 y^3}{1 - 2by}$  ensures that  $P$  and

$$Q = \left( -(a_0 + c)y + a_1 y^2 \left( 1 + \frac{4by}{1 - 2by} \right), -y, -y, \frac{4a_1 y^3}{1 - 2by} \right)$$

are aligned and both in  $\Sigma$ . Then the whole segment connecting  $P$  and  $Q$  belongs to  $\Sigma$ . In particular

$$\frac{P + Q}{2} = \left( \frac{a_1 y^2}{1 - 2by}, 0, 0, \frac{2a_1 y^3}{1 - 2by} \right) \in \Sigma.$$

Then  $\frac{a_1 y^2}{1 - 2by} = b \cdot \frac{2a_1 y^3}{1 - 2by}$ , which forces  $a_1 y^2 = 0$  for all  $y$  and ultimately  $a_1 = 0$ .  $\square$

**Proposition 4.6** (Monotone  $(pX + qY)$ -epigraphs (global statement)). *Let the hypotheses of Proposition 4.5 with  $(p, q) = (1, 0)$  be satisfied and let  $x = ay + bt + cu$  be the equation of the plane  $\Sigma$  appearing in Proposition 4.5. Then we have  $\Sigma = \partial E$ .*

Before the proof, recall the following geometric property concerning a given pair of lines  $\ell_1$  and  $\ell_2$  in the Heisenberg group parametrized by  $\Gamma_1(s) = (z_1, t_1) \cdot (s\zeta_1, 0)$  and  $\Gamma_2(s) = (z_2, t_2) \cdot (s\zeta_2, 0)$ . See the discussion in [7, p. 343]. If  $\ell_1$  and  $\ell_2$  are *skew* or *parallel with distinct projection*, then for all  $\sigma \in \mathbb{R}$ , except at least a singular value  $\sigma_0$ , there is  $s_\sigma \in \mathbb{R}$  such that  $\Gamma_1(\sigma)$  and  $\Gamma_2(s_\sigma)$  are aligned.

For the sake of completeness, recall that two lines  $\ell_1$  and  $\ell_2$  parametrized by  $\Gamma_1(s) = (z_1, t_1) \cdot (s\zeta_1, 0)$  and  $\Gamma_2(s) = (z_2, t_2) \cdot (s\zeta_2, 0)$  are *parallel* if  $\zeta_1$  and  $\zeta_2$  are linearly dependent and nonzero. They are *parallel with distinct projection* if they are parallel and their projections on the plane  $x, y$  are (parallel and) different lines in the plane. They are *skew* if they are not parallel and have empty intersection.

*Proof of Proposition 4.6:* Assume without loss of generality that the equation of  $\Sigma$  is  $x = cu + bt$ . This can be obtained after a rotation in the variables  $(x, y)$ . We know by Proposition 4.5 that there is  $\varepsilon > 0$  such that

$$(4.9) \quad \begin{aligned} A := \{ & (cu + bt, y, u, t) : |y|, |u|, |t| < \varepsilon \} \subset \partial E \quad \text{and} \\ & \{ (cu + bt, y, u, t) \cdot (s, 0, 0, 0) : |y|, |u|, |t| < \varepsilon, 0 < s < +\infty \} \\ & \subset \text{int}(E), \end{aligned}$$

while the same set with  $s < 0$  is contained in  $\text{int}(E^c)$ .

*Step 1:* Consider points of the form  $(bt, 0, 0, t) \in \Sigma$ . These points are in  $\partial E$  as soon as  $|t| < \varepsilon$ . Let  $\eta, \omega \in \mathbb{R}$  with  $|\eta| < \frac{1}{2b}$  and take the line

$$(4.10) \quad \begin{aligned} \gamma_1(\sigma) &= (bt, 0, 0, t) \cdot ([c\omega - 2b^2t\eta]\sigma, \sigma\eta, \sigma\omega, 0) \\ &= (bt(1 - 2b\eta\sigma) + c\sigma\omega, \sigma\eta, \sigma\omega, t(1 - 2b\eta\sigma)) \quad \text{with } \sigma \in \mathbb{R}. \end{aligned}$$

If  $\eta$  and  $\omega$  are close to zero, we have  $\gamma_1(0)$  and  $\gamma_1(1) \in \partial E$ . Therefore,  $\gamma_1(\mathbb{R}) \subset \partial E$ . Observe also that, by (4.9), if  $t, \eta, \omega$ , and  $\sigma$  are sufficiently close to 0, the point  $\gamma_1(\sigma)$  belongs to the small surface  $A$  appearing in (4.9).

*Step 2:* We claim that  $\partial E$  is contained in  $\Sigma$ .

Assume by contradiction that a point  $P = (\bar{x}, \bar{y}, \bar{u}, \bar{t})$  belongs to  $\partial E \setminus \Sigma$ . By Lemma 3.5, there is  $(p, q) \neq (0, 0)$  such that

$$(4.11) \quad \gamma_2(s) := (\bar{x}, \bar{y}, \bar{u}, \bar{t}) \cdot (ps, qs, 0, 0) \in \partial E \quad \text{for all } s \in \mathbb{R}.$$

We discuss first the case  $(b, c) \neq (0, 0)$ . In this case, it is easy to see that we can find  $t \in ]\varepsilon, \varepsilon[$ ,  $\eta, \omega \in \mathbb{R}$  such that the vectors  $(p, q)$  and  $(c\omega - 2b^2t\eta, \eta)$  appearing in (4.11) and in (4.10), respectively, are independent. After this choice, if we indicate by  $\pi_{\mathbb{H}}(x, y, u, t) = (x, y, t)$ , the lines  $\pi_{\mathbb{H}}\gamma_1$  and  $\pi_{\mathbb{H}}\gamma_2$  are not parallel in  $\mathbb{H}$ . Taking if needed a small modification of  $t$ , we may assume that they are also skew. Therefore, by the property discussed before the proof, we can find  $\sigma$  as close as we wish to the origin and a corresponding  $s_\sigma \in \mathbb{R}$  such that  $\pi_{\mathbb{H}}\gamma_1(\sigma)$  and  $\pi_{\mathbb{H}}\gamma_2(s_\sigma)$  are aligned. Denote by  $\gamma_3$  the line connecting  $\gamma_1(\sigma) = \gamma_3(0)$  and  $\gamma_2(s_\sigma) = \gamma_3(1)$ . Our line  $\gamma_3$  is contained in  $\partial E$ , is transversal to  $\Sigma$ , and touches  $\Sigma$  at a point very close to the origin. This contradicts the second line of (4.9) and concludes Step 2, at least in the case  $(b, c) \neq (0, 0)$ .

In order to discuss Step 2 for  $b = c = 0$ , note that in this case  $\Sigma$  is the plane  $x = 0$  and we can start from the inclusion  $\{(0, y, u, t) : |t| < \varepsilon, (y, u) \in \mathbb{R}^2\} \subset \partial E$ . Let us consider the curve  $\gamma_1(\sigma) = (0, \sigma, 0, t)$ , take  $(\bar{x}, \bar{y}, \bar{u}, \bar{t}) \in \partial E$  with  $\bar{x} \neq 0$  and the corresponding line  $\gamma_2$  of the form (4.11). If  $p = 0$ , then, since  $\bar{x} \neq 0$ ,  $\pi_{\mathbb{H}}\gamma_1$  and  $\pi_{\mathbb{H}}\gamma_2$  are parallel with distinct projection and we can conclude as in the case  $(b, c) \neq (0, 0)$  above. If instead,  $p \neq 0$ , the lines  $\pi_{\mathbb{H}}\gamma_1$  and  $\pi_{\mathbb{H}}\gamma_2$  are not parallel. Then, changing if needed the choice of  $t$  we may assume that they are skew. In either case, we get the same contradiction as in the case  $(b, c) \neq (0, 0)$ .

*Step 3:* We have  $\Sigma = \partial E$ . To show this claim, note that we already know that  $\partial E \subset \Sigma$ . If  $\Sigma \neq \partial E$ , then  $(\partial E)^c$  is open and connected. Therefore, it is contained either in  $\text{int}(E)$  or in  $\text{int}(E^c)$ . Assuming the former and considering the line  $\gamma(s) = (s, 0, 0, 0)$ , we see that both  $\gamma(1)$  and  $\gamma(-1)$  belong to  $\text{int}(E)$ , getting a contradiction.  $\square$



**4.3. Classification of monotone  $U$ -epigraphs.** In this section we discuss case (ii) of Lemma 4.4.

**Proposition 4.7** (Monotone  $U$ -epigraphs (local statement)). *Let  $(0, 0, 0, 0) \in \partial E$  and let  $(0, 0, 1, 0) \in \text{int } E$ . Then there are  $\varepsilon > 0$ , a neighborhood  $W$  of the origin in  $\mathbb{R}^3$ , and a linear function  $\psi: W \rightarrow \mathbb{R}$  with  $\psi(0, 0, 0) = 0$  such that  $\text{graph}(\psi) \subset \partial E$ , while  $\text{epi}(\psi) \subset \text{int}(E)$  and  $\text{ipo}(\psi) \subset \text{int}(E^c)$ .*

By the direct product structure, Euclidean and intrinsic  $U$ -graphs are the same. As a by-product of the following proof, we get a classification of all “horizontally affine” functions  $\psi: \mathbb{H} \rightarrow \mathbb{R}$ , i.e. all functions such that  $\psi((z, t) \cdot (\lambda \zeta, 0)) - \psi(z, t) = \lambda[\psi((z, t) \cdot (\zeta, 0)) - \psi(z, t)]$  for all  $\lambda, z, t, \zeta$ . This classification could have some independent interest.

*Proof:* By Theorem 3.7 the boundary of  $E$  is locally an  $U$ -graph. Namely, there is a neighborhood  $W$  of the origin in  $\mathbb{R}^3$  such that

$$(4.12) \quad \begin{aligned} &\{(x, y, \psi(x, y, t), t) : (x, y, t) \in W\} \subset \partial E, \\ &\{(x, y, u, t) : (x, y, t) \in W, +\infty > u > \psi(x, y, t)\} \subset \text{int}(E), \text{ and} \\ &\{(x, y, u, t) : (x, y, t) \in W, -\infty < u < \psi(x, y, t)\} \subset \text{int}(E^c). \end{aligned}$$

Let  $\Omega = W \times \mathbb{R}$ . Clearly, any point  $(x, y, u, t) \in \Omega$  belongs to one and only one of the three sets in the left-hand sides above.

The proof is articulated into three steps.

*Step 1:* We show that

$$\psi(z, t) = m(z)t + q(z),$$

for suitable functions  $m, q$  defined in a neighborhood of the origin in  $\mathbb{R}^2$ .

Let  $P = (x, y, \psi(x, y, t), t) = (z, \psi(z, t), t) \in \partial E \cap \Omega$ . Let us consider the horizontal plane

$$H_P = \{(x, y, \psi(x, y, t), t) \cdot (\xi, \eta, u, 0) = P \cdot (\xi, \eta, u, 0) : (\xi, \eta, u) \in \mathbb{R}^3\}.$$

Inside this plane, given an angle  $\theta$  consider the two-dimensional subspace

$$A_\theta := \{P \cdot (\lambda e^{i\theta}, u, 0) : (\lambda, u) \in \mathbb{R}^2\}.$$

We claim that  $A_\theta \cap \partial E$  consists exactly of a line. More precisely, there is  $\alpha(x, y, t, \theta) \in \mathbb{R}$  such that

$$A_\theta \cap \partial E = \{(x, y, \psi(x, y, t), t) \cdot (s \cos \theta, s \sin \theta, s\alpha(x, y, t, \theta), 0) : s \in \mathbb{R}\}.$$

To show the claim, note that Lemma 3.5 ensures that  $A_\theta \cap \partial E$  contains at least a line passing by  $P$ . We shall show that, assuming there are two different lines both contained in  $A_\theta$  and both containing  $P$ , we get a contradiction. Let  $\ell_1$  and  $\ell_2$  be such lines. None of them can have the

form  $\{P \cdot (0, 0, s, 0) : s \in \mathbb{R}\}$ , because we know that  $P \cdot (0, 0, s, 0) \in \text{int}(E)$  for all  $s > 0$ . Then we can write

$$\ell_j = \{\gamma_j(s) := P \cdot (se^{i\theta}, \alpha_j s, 0) : s \in \mathbb{R}\} \quad \text{for } j = 1, 2,$$

where  $\alpha_1 \neq \alpha_2$ . Since  $\alpha_1 \neq \alpha_2$ , we find for all  $s \neq 0$  a couple of distinct aligned points  $\gamma_1(s)$  and  $\gamma_2(s)$ . Then

$$\ell_s = \{(z, \psi(z, t), t) \cdot (se^{i\theta}, u, 0) : u \in \mathbb{R}\},$$

the line connecting them, is entirely contained in  $\partial E$ , by Lemma 3.4. Taking  $u = 1$  and  $s \neq 0$  very close to 0, we find points in  $\partial E$  as close as we wish to the point  $(z, \psi(z, t) + 1, t) \in \text{int}(E)$ . This gives a contradiction and proves the claim.

To accomplish Step 1, fix  $(x, y) = z$  close to the origin and consider  $\theta \neq \varphi \pmod{\pi}$ , and  $\tau \neq t$ . The lines provided by the previous construction are

$$\{\gamma_3(s) := (z, \psi(z, t), t) \cdot (se^{i\theta}, s\alpha_z(t, \theta), 0) : s \in \mathbb{R}\} \subset \partial E \quad \text{and}$$

$$\{\gamma_4(\sigma) := (z, \psi(z, \tau), \tau) \cdot (\sigma e^{i\varphi}, \sigma\alpha_z(\tau, \varphi), 0) : \sigma \in \mathbb{R}\} \subset \partial E,$$

where we wrote  $\alpha_z(t, \theta) = \alpha(x, y, t, \theta)$ . Following the ideas of (the proof of) [19, Theorem 1.2] and the similar argument in [7, Lemma 4.10], we look for pairs of aligned points  $\gamma_3(s)$  and  $\gamma_4(\sigma)$ . A short calculation shows that  $\gamma_3(s)$  is aligned with  $\gamma_4(-\frac{\eta}{s})$  for all  $s \neq 0$ , where  $\eta := \frac{\tau - t}{2 \sin(\varphi - \theta)}$ . Fixed any  $s > 0$  (also  $s = 1$  does the job), the line  $\gamma_5$  connecting  $\gamma_3(s)$  and  $\gamma_4(-\eta/s)$  is

(4.13)

$$\begin{aligned} \gamma_5(\lambda) &= (z, 0, 0) \cdot (0, \psi(z, t), t) \cdot (se^{i\theta}, s\alpha_z(t, \theta), 0) \\ &\cdot \left( -\lambda \left( \frac{\eta}{s} e^{i\varphi} + se^{i\theta} \right), \lambda \left( \psi(z, \tau) - \psi(z, t) - \frac{\eta}{s} \alpha_z(\tau, \varphi) - s\alpha_z(t, \theta) \right), 0 \right) \\ &=: (z, 0, 0) \cdot (\widehat{z}_5(\lambda), \psi(z, t) + \lambda[\psi(z, \tau) - \psi(z, t)] + \widehat{u}_5(\lambda), t + \lambda(\tau - t)), \end{aligned}$$

where

$$\begin{aligned} \widehat{z}_5(\lambda) &= se^{i\theta} - \lambda \left( \frac{\eta}{s} e^{i\varphi} + se^{i\theta} \right), \\ \widehat{u}_5(\lambda) &= s\alpha_z(t, \theta) - \lambda \left( \frac{\eta}{s} \alpha_z(\tau, \varphi) + s\alpha_z(t, \theta) \right), \end{aligned} \quad (4.14)$$

while the form of the last coordinate follows from the choice of  $\eta$ . Note that both  $\widehat{z}_5(\lambda)$  and  $\widehat{u}_5(\lambda)$  change sign, if we change  $s = 1$  with  $s = -1$ . Therefore, if we denote by  $\gamma_6$  the curve in (4.13) and (4.14) obtained

with  $s = -1$ , the points  $\gamma_5(\lambda)$  and  $\gamma_6(\lambda) \in \partial E$  are aligned for all  $\lambda \in \mathbb{R}$ . Then their intermediate point belongs to  $\partial E$  for all  $\lambda$ . Ultimately

$$(z, \psi(z, t) + \lambda[\psi(z, \tau) - \psi(z, t)], t + \lambda(\tau - t)) \in \partial E$$

for all  $\lambda \in \mathbb{R}$  and  $t, \tau$  close to 0. In particular, choosing  $\lambda = -t/(\tau - t)$  for any pair  $t \neq \tau$ , we see that

$$\left(z, \psi(z, t) - t \frac{\psi(z, \tau) - \psi(z, t)}{\tau - t}, 0\right) \in \partial E$$

for all  $t \neq \tau$  with  $t, \tau$  close to 0. Since for all  $z$  close to the origin, the line  $\{(z, u, 0) : u \in \mathbb{R}\}$  contains only one point  $u = \psi(z, 0) \in \partial E$ , this implies that  $\psi(z, t) = m(z)t + q(z)$  for suitable functions  $m, q$  and for all  $(z, t)$  close to the origin.

*Step 2:* We show that

- (i)  $q(\lambda z) = \lambda q(z)$  for all small  $|z|$  and  $|\lambda| \leq 1$ ;
- (ii)  $m(\lambda z) - m(0) = \lambda[m(z) - m(0)]$ , for all small  $|z|$  and  $|\lambda| \leq 1$ .

If  $q(0)$  would not vanish, the function  $q$  would satisfy (i), as  $m$ .

To check (i), consider the points  $(z, \psi(z, 0), 0) = (z, q(z), 0) \in \partial E$  for all  $z$  close to the origin. Each of these points is aligned with  $(0, q(0), 0) = (0, 0, 0)$ . Then, by Lemma 3.4, we have  $(\lambda z, \lambda q(z), 0) \in \partial E$  for all  $\lambda \in \mathbb{R}$ . Thus, from Theorem 3.7, for small  $|\lambda|$  and  $|z|$ , we get  $q(\lambda z) = \lambda q(z)$ .

To show property (ii) for  $m$ , it suffices to consider the pair of points  $P = (z, m(z)t + q(z), t)$  and  $Q = (0, m(0)t, t)$ , both in the boundary of  $E$ . The points  $P, Q$  are aligned and can be connected by

$$\gamma(\lambda) = (\lambda z, m(0)t + \lambda[m(z)t + q(z) - m(0)t], t) \quad \text{with } \lambda \in \mathbb{R}.$$

Thus, we get the identity  $m(\lambda z)t + q(\lambda z) = m(0)t + \lambda[m(z)t + q(z) - m(0)t]$ , whose linear part in  $t$  gives (ii).

Note that, in spite of the fact that Theorem 3.7 is local, the functions  $m$  and  $q$  can be defined globally, by their homogeneity property.

*Step 3:* We show that  $q(z) = q(x, y) = ax + by$  is linear and  $m(z)$  is constant.

To accomplish Step 3, observe that for all  $(z, t)$  and  $(\zeta, 0)$  close to the origin, the line

$$\begin{aligned} \gamma(\lambda) &= (z, \psi(z, t), t) \cdot (\lambda \zeta, \lambda\{\psi((z, t) \cdot (\zeta, 0)) - \psi(z, t)\}, 0) \\ &= (z + \lambda \zeta, m(z)t + q(z) + \lambda\{m(z + \zeta)(t + 2\Im(z\bar{\zeta})) \\ &\quad + q(z + \zeta) - m(z)t - q(z)\}, t + 2\lambda\Im(z\bar{\zeta})) \end{aligned}$$

satisfies  $\gamma(0) \in \partial E$  and  $\gamma(1) \in \partial E$ . Then  $\gamma(\lambda) \in \partial E$  for all  $\lambda \in [0, 1]$ , and therefore

$$(4.15) \quad m(z + \lambda\zeta)(t + 2\lambda\Im(z\bar{\zeta})) + q(z + \lambda\zeta) \\ = m(z)t + q(z) + \lambda[m(z + \zeta)(t + 2\Im(z\bar{\zeta})) + q(z + \zeta) - m(z)t - q(z)].$$

Equating linear terms in  $t$  gives

$$(4.16) \quad m(z + \lambda\zeta) = m(z) + \lambda[m(z + \zeta) - m(z)],$$

for all  $z, \zeta$  close to the origin and  $\lambda \in [0, 1]$ , which means that  $m$  is affine, i.e.  $m(x, y) = c + kx + hy$  for suitable constants  $c, k, h \in \mathbb{R}$ . To see this fact, fix  $\varepsilon_0$  small, take  $|x|, |y| < \varepsilon_0$  and write

$$\begin{aligned} m(x, y) &= m(0, 0) + m(x, y) - m(x, 0) + m(x, 0) - m(0, 0) \\ &= m(0, 0) + y \frac{m(x, \varepsilon_0) - m(x, 0)}{\varepsilon_0} + x \frac{m(\varepsilon_0, 0) - m(0, 0)}{\varepsilon_0} \\ &=: c + yk(x) + hx. \end{aligned}$$

Then, the homogeneity of  $z \mapsto m(z) - m(0)$  forces  $k(x) = \text{constant}$ .

Next we consider (4.15) with  $t = 0$  using the form (4.16) of  $m(z + \lambda\zeta)$ . This gives

$$\begin{aligned} 2\lambda\Im(z\bar{\zeta})\{m(z) + \lambda[m(z + \zeta) - m(z)]\} + q(z + \lambda\zeta) \\ = q(z) + 2\lambda\Im(z\bar{\zeta})m(z + \zeta) + \lambda[q(z + \zeta) - q(z)]. \end{aligned}$$

This implies that

$$(4.17) \quad q(z + \lambda\zeta) - q(z) = \lambda(q(z + \zeta) - q(z)) \\ + 2\lambda(1 - \lambda)\Im(z\bar{\zeta})[m(z + \zeta) - m(z)].$$

Taking two derivatives in  $\lambda$ , we conclude that  $m(z + \zeta) - m(z)$  must vanish for all  $z, \zeta$  with  $\Im(z\bar{\zeta}) \neq 0$  and by continuity of  $m$ , for all  $z, \zeta \in \mathbb{R}^2$ . Then  $m$  is constant and ultimately (4.17) shows that  $q$  is affine.

This completes the proof of Proposition 4.7.  $\square$

Next, we show the following global version of Proposition 4.7.

**Proposition 4.8.** *Let the hypotheses of Proposition 4.7 be satisfied and let  $u = ax + by + ct$  be the equation of the three-dimensional plane  $\Sigma$  coming from Proposition 4.7. Then  $\Sigma = \partial E$ .*

*Proof:* Let  $W = \{(x, y, t) \in \mathbb{R}^3 : |x|, |y|, |t| < \varepsilon\}$  be a neighborhood of the origin such that the statement of Proposition 4.7 holds true.

*Step 1:* We first show that  $\partial E \supset \Sigma \cap \{(x, y, ax + by + ct, t) : |t| < \varepsilon \text{ and } (x, y) \in \mathbb{R}^2\}$ .

To show this claim, given  $t \in ]-\varepsilon, \varepsilon[$  it suffices to consider the points  $(0, 0, ct, t)$  and  $(x, y, ax + by + ct, t)$ , where  $|x|, |y| < \varepsilon$ . Both points belong to  $\partial E$  and they are aligned. Then,  $\gamma(s) = (sx, sy, s(ax + by) + ct, t) \in \partial E$ , for all  $s \in \mathbb{R}$ . Step 1 follows immediately.

*Step 2:* We show that  $\Sigma \subset \partial E$  or, in other words, for all  $(\xi, \eta, \tau) = (\zeta, \tau) \in \mathbb{R}^3$  we have

$$(4.18) \quad (\xi, \eta, a\xi + b\eta + c\tau, \tau) \in \partial E.$$

It suffices to show this statement for  $\zeta \neq 0$ . Let  $\tau \in \mathbb{R}$  and consider the curve

$$\begin{aligned} \gamma(\mu) &:= \left( \xi - \frac{\tau\eta}{2|\zeta|^2}, \eta + \frac{\tau\xi}{2|\zeta|^2}, a\left(\xi - \frac{\tau\eta}{2|\zeta|^2}\right) + b\left(\eta + \frac{\tau\xi}{2|\zeta|^2}\right), 0 \right) \\ &\quad \cdot \left( \mu \frac{\tau\eta}{2|\zeta|^2}, -\mu \frac{\tau\xi}{2|\zeta|^2}, \mu \left[ a \frac{\tau\eta}{2|\zeta|^2} - b \frac{\tau\xi}{2|\zeta|^2} + c\tau \right], 0 \right) \\ &= \left( \xi - (1 - \mu) \frac{\tau\eta}{2|\zeta|^2}, \eta + (1 - \mu) \frac{\tau\xi}{2|\zeta|^2}, \right. \\ &\quad \left. a \left[ \xi - (1 - \mu) \frac{\tau\eta}{2|\zeta|^2} \right] + b \left[ \eta + (1 - \mu) \frac{\tau\xi}{2|\zeta|^2} \right] + c\mu\tau, \mu\tau \right). \end{aligned}$$

We have  $\gamma(0) \in \partial E$ . Furthermore, if  $\mu > 0$  is so small that  $\mu\tau < \varepsilon$ , then by Step 1 we have  $\gamma(\mu) \in \partial E$ . Therefore,  $\gamma(\mu) \in \partial E$  for all  $\mu \in \mathbb{R}$ . Taking  $\mu = 1$ , we see that (4.18) holds.

*Step 3:* We show that  $\Sigma := \{(x, y, ax + by + ct, t) : (x, y, t) \in \mathbb{R}^3\} = \partial E$ .

By Step 2, we know that  $\Sigma \subset \partial E$ . Assume by contradiction that there is  $(\xi, \eta, u, \tau) \in \mathbb{R}^4 \setminus \Sigma$  belonging to  $\partial E$ . Then the whole line  $\{(\xi, \eta, u, \tau) : u \in \mathbb{R}\}$  is contained in  $\partial E$ . Up to a translation, we may assume that

$$\begin{aligned} \Sigma &= \{(x, y, \alpha x + \beta y + \gamma t, t) : (x, y, t) \in \mathbb{R}^3\} \subset \partial E \quad \text{and} \\ \ell &= \{(0, 0, u, 0) : u \in \mathbb{R}\} \subset \partial E. \end{aligned}$$

In particular, all points of the form  $P = (x, y, \alpha x + \beta y, 0)$  and  $Q = (0, 0, u, 0)$  belong to  $\partial E$  for all  $x, y, u$ . Since  $P$  and  $Q$  are aligned, we conclude that the intermediate point  $(\frac{x}{2}, \frac{y}{2}, \frac{\alpha x + \beta y + u}{2}, 0)$  belongs to  $\partial E$  for arbitrary  $x, y, u \in \mathbb{R}$ . Then the plane  $t = 0$  is fully contained in  $\partial E$  and Lemma 4.3 gives a contradiction.  $\square$

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