# Convex isoperimetric sets in the Heisenberg group 

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#### Abstract

We characterize convex isoperimetric sets in the Heisenberg group. We first prove Sobolev regularity for a certain class of $\mathbb{R}^{2}$-valued vector fields of bounded variation in the plane related to the curvature equations. Then we show that the boundary of convex isoperimetric sets is foliated by geodesics of the Carnot-Carathéodory distance.


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## 1. Introduction

We identify the Heisenberg group $\mathbb{H}^{1}$ with $\mathbb{C} \times \mathbb{R}$ endowed with the group law

$$
(z, t)\left(z^{\prime}, t^{\prime}\right)=\left(z+z^{\prime}, t+t^{\prime}+2 \operatorname{Im}\left(z \bar{z}^{\prime}\right)\right)
$$

where $t, t^{\prime} \in \mathbb{R}$ and $z=x+i y, z^{\prime}=x^{\prime}+i y^{\prime} \in \mathbb{C}$. The Lie algebra of left-invariant vector fields is spanned by

$$
X=\frac{\partial}{\partial x}+2 y \frac{\partial}{\partial t}, \quad Y=\frac{\partial}{\partial y}-2 x \frac{\partial}{\partial t} \quad \text { and } \quad T=\frac{\partial}{\partial t}
$$

and the distribution of planes spanned by $X$ and $Y$, called horizontal distribution, generates the Lie algebra by brackets.

The natural volume in $\mathbb{H}^{1}$ is the Haar measure, which, up to a positive factor, coincides with Lebesgue measure in $\mathbb{R}^{3}=\mathbb{C} \times \mathbb{R}$. Lebesgue measure is also the Riemannian volume of the left-invariant metric for which $X, Y$ and $T$ are orthonormal. We denote by $|E|$ the volume of a (Lebesgue) measurable set $E \subset \mathbb{H}^{1}$. The horizontal perimeter (or simply perimeter) of $E$ is

$$
\begin{equation*}
P(E)=\sup \left\{\int_{E}\left(X \varphi_{1}+Y \varphi_{2}\right) d x d y d t \mid \varphi_{1}, \varphi_{2} \in C_{c}^{1}\left(\mathbb{R}^{3}\right), \varphi_{1}^{2}+\varphi_{2}^{2} \leq 1\right\} \tag{1.1}
\end{equation*}
$$

If $P(E)<+\infty$, the set $E$ is said to be of finite perimeter. Perimeter is left-invariant and 3-homogeneous with respect to the group of dilations $\delta_{\lambda}: \mathbb{H}^{1} \rightarrow \mathbb{H}^{1}, \delta_{\lambda}(z, t)=$
$\left(\lambda z, \lambda^{2} t\right), \lambda>0$, that is $P\left(\delta_{\lambda}(E)\right)=\lambda^{3} P(E)$. Definition (1.1) of perimeter is modelled on De Giorgi's notion of perimeter in Euclidean spaces which was generalized in [10] to Carnot-Carathéodory spaces. For smooth sets (e.g. of class $C^{2}$ ), perimeter coincides with the Minkowski content and with the 3-dimensional Hausdorff measure of $\partial E$ constructed by means of the standard Carnot-Carathéodory metric in $\mathbb{H}^{1}$ (see [16] and [9], respectively). Volume and perimeter are related via the isoperimetric inequality

$$
\begin{equation*}
|E| \leq C_{\text {isop }} P(E)^{4 / 3} \tag{1.2}
\end{equation*}
$$

where $C_{\text {isop }}>0$ and $E \subset \mathbb{H}^{1}$ is any measurable set with finite perimeter and volume. This inequality is proved by P. Pansu in [17] and [18] for smooth domains with the 3 -dimensional Hausdorff measure of $\partial E$ in $\mathbb{H}^{1}$ replacing $P(E)$. In the general form (1.2), the inequality is proved in [10] (see also [8]). The problem of computing the sharp constant $C_{\text {isop }}$ leads to the notion of isoperimetric set. A set $E \subset \mathbb{H}^{1}$ with $0<|E|<+\infty$ is an isoperimetric set if it minimizes the isoperimetric ratio

$$
\begin{equation*}
\operatorname{Isop}(E)=\frac{P(E)^{4 / 3}}{|E|} \tag{1.3}
\end{equation*}
$$

Isoperimetric sets do exist: this is proved in [12] by a concentration-compactness argument. Pansu notes that the boundary of a smooth isoperimetric set has "constant mean curvature" and that a smooth surface has "constant mean curvature" if and only if it is foliated by horizontal lifts of plane circles with constant radius. Then he conjectures that an isoperimetric set is obtained by rotating around the center of the group a geodesic joining two points in the center. Recently, Pansu's conjecture reappeared in [11].

The problem of determining isoperimetric sets is interesting for two reasons. On the one hand, the non commutative group law makes it difficult to prove by a rearrangement argument that the isoperimetric ratio (1.3) is minimized by rotationally symmetric sets, as it is natural to conjecture. On the other hand, there is no regularity theory for measurable sets in $\mathbb{H}^{1}$ minimizing perimeter with (or without) a volume constraint. So, new techniques and ideas are needed.

All known results assume either symmetry or regularity (or both). In fact, assuming rotational symmetry and regularity it is easy to determine the isoperimetric profile (see [13], and [22] for the general case). Actually, it suffices to assume the rotational symmetry of a certain horizontal section (see [7] and especially the calibration argument in [21]). The solution of the rotationally symmetric case with no regularity assumption is in [14]. On the other hand, isoperimetric sets can be also determined assuming only the $C^{2}$ regularity of the boundary and no symmetry (see [23]). Further evidence supporting Pansu's conjecture is provided by [15], where a 2-dimensional version of the problem is solved. We refer to the monograph [3] for a more detailed introduction to the isoperimetric problem in the Heisenberg group.

In this article, we characterize isoperimetric sets which are convex. By convex set we mean a subset of $\mathbb{H}^{1}=\mathbb{R}^{3}$ which is convex with respect to the standard
vector space structure of $\mathbb{R}^{3}$. Convexity is a left invariant property in $\mathbb{H}^{1}$ because left translations are affine mappings.

Theorem 1.1 (Convex isoperimetric sets). Up to a left translation and a dilation, any closed convex isoperimetric set in $\mathbb{H}^{1}$ coincides with

$$
\begin{equation*}
E_{\text {isop }}=\left\{(z, t) \in \mathbb{H}^{1}| | t|\leq \arccos | z\left|+|z| \sqrt{1-|z|^{2}},|z| \leq 1\right\}\right. \tag{1.4}
\end{equation*}
$$

The boundary of the set in (1.4) is foliated by Heisenberg geodesics, it is globally of class $C^{2}$, but it fails to be of class $C^{3}$ at the north and south poles $(0, \pm \pi / 2)$. At these points, the plane spanned by the vector fields $X$ and $Y$, the horizontal plane, is tangent to the boundary.

We explain the main steps in the proof of Theorem 1.1. Let us consider a bounded convex set of the form

$$
\begin{equation*}
E=\left\{(z, t) \in \mathbb{H}^{1} \mid z \in D, f(z) \leq t \leq g(z)\right\} \tag{1.5}
\end{equation*}
$$

where $D \subset \mathbb{R}^{2}$ is a compact convex set in the plane with nonempty interior, and $-g, f: D \rightarrow \mathbb{R}$ are convex functions. If $E$ is isoperimetric, then the function $f: D \rightarrow \mathbb{R}$ satisfies the partial differential equation

$$
\begin{equation*}
\operatorname{div}\left(\frac{\nabla f(z)+2 z^{\perp}}{\left|\nabla f(z)+2 z^{\perp}\right|}\right)=\frac{3 P(E)}{4|E|} \tag{1.6}
\end{equation*}
$$

in $\operatorname{int}(D) \backslash \Sigma(f)$. Here and in the following we let $z=(x, y)$ and $z^{\perp}=(-y, x)$. We call the singular set $\Sigma(f)=\left\{z \in \operatorname{int}(D) \mid-2 z^{\perp} \in \partial f(z)\right\}$ the characteristic set of $f$. This set is always contained in a line segment. The basic facts concerning $\Sigma(f)$ are discussed in the appendix. Equation (1.6) is derived in Section 2.

The curvature operator in (1.6) has been studied by several authors under different regularity assumptions (besides the previous references, see also [4-6, 19, 20]). In the smooth case, the number

$$
\begin{equation*}
H=\frac{3 P(E)}{4|E|} \tag{1.7}
\end{equation*}
$$

is known as the (horizontal) curvature of $\partial E$. In our case, equation (1.6) is to be interpreted in distributional sense. The distributional derivatives of $u(z)=\nabla f(z)+$ $2 z^{\perp}$ are measures, because $f$ is a convex function, and so the equation states that the distributional divergence of $u /|u|$ is constant.

Our goal is to give equation (1.6) a pointwise meaning along integral curves of the vector field orthogonal to $u /|u|$. The first step is to show that the equation holds in the Sobolev sense. Precisely, we prove that the distributional derivative of $u /|u|$ is a measure which is absolutely continuous with respect to Lebesgue measure. This result is a corollary of the following regularity theorem for $B V$ vector fields, which is interesting in itself. We prove a slightly more general statement in the third section.

Theorem 1.2 (Improved regularity). Let $\Omega \subset \mathbb{R}^{2}$ be a bounded open set and let $u \in B V\left(\Omega ; \mathbb{R}^{2}\right)$ be a vector field. Suppose that:
(i) There exists $\delta>0$ such that $|u(z)| \geq \delta$ for $\mathcal{L}^{2}$-a.e. $z \in \Omega$;
(ii) $\operatorname{div} u^{\perp} \in L^{1}(\Omega)$;
(iii) $\operatorname{div}\left(\frac{u}{|u|}\right) \in L^{1}(\Omega)$.

Then we have $u /|u| \in W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$.
The improved regularity of the boundary is the starting point for the geometric characterization of convex isoperimetric sets. The vector field $v(z)=-u^{\perp}(z)=$ $2 z-\nabla f^{\perp}(z)$ belongs to $B V_{\text {loc }}\left(\operatorname{int}(D) ; \mathbb{R}^{2}\right)$. Moreover, its distributional divergence is in $L^{\infty}$, in fact

$$
\begin{equation*}
\operatorname{div} v=4 \quad \operatorname{in} \operatorname{int}(D) \tag{1.8}
\end{equation*}
$$

Thanks to the theory on the Cauchy Problem for $B V$ vector fields recently developed by Ambrosio in [1], the bound on the divergence ensures the existence of a unique regular Lagrangian flow $\Phi: K \times[-\varrho, \varrho] \rightarrow D$ starting from a compact set $K \subset \operatorname{int}(D)$, where $\varrho>0$ is small enough. For $\mathcal{L}^{2}$-a.e. $z \in K$, the curve $s \mapsto \Phi(z, s)$ is an integral curve of $v$ passing through $z$ at time $s=0$. In Section 4, we show that $\mathcal{L}^{2}$-a.e. integral curve of $v$ is an arc of circle.

Theorem 1.3 (Foliation by circles). Let $E \subset \mathbb{H}^{1}$ be a convex isoperimetric set of the form (1.5) and let $K \subset \operatorname{int}(D) \backslash \Sigma(f)$ be a compact set. Then, for $\mathcal{L}^{2}$-a.e. $z \in K$, the curve $s \mapsto \Phi(z, s)$ is an arc of circle having radius $1 / H$, with $H>0$ as in (1.7).

The proof of Theorem 1.3 relies upon a reparameterization argument. The vector field $v$ has a regular flow $\Phi$ starting from $K \subset \operatorname{int}(D) \backslash \Sigma(f)$ but it is only in $B V_{\text {loc }}\left(\operatorname{int}(D) ; \mathbb{R}^{2}\right)$. On the other hand, $v /|v|$ is in $W_{\text {loc }}^{1,1}\left(\operatorname{int}(D) ; \mathbb{R}^{2}\right)$, but its divergence is only in $L_{\text {loc }}^{1}(\operatorname{int}(D))$ and so we have no regular flow for $v /|v|$. In order to compute the second order derivative of a generic integral curve of $v$, we introduce a suitable reparameterization $\gamma(s)=\Phi(z, \tau(s))$. We show that for any vector field $w \in W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$ defined in some open neighborhood $\Omega$ of $K$, the curve $\kappa(s)=w(\gamma(s))$ is in $W^{1,1}$ for $\mathcal{L}^{2}$-a.e. $z \in K$ and moreover

$$
\begin{equation*}
\dot{\kappa}=(\nabla w \circ \gamma) \dot{\gamma} \quad \text { in the weak sense. } \tag{1.9}
\end{equation*}
$$

Using the chain rule (1.9) with $w=v /|v|$, which is in $W^{1,1}$ by Theorem 1.2, equation (1.6) can be given a pointwise meaning along the flow, and the integral curves of $v$ turn out to be arcs of circles.

Since the horizontal lift of the flow $\Phi$ foliates the graph of the function $f$, it follows that the bottom (and upper) part of the boundary of $E$ is foliated by geodesics of the Heisenberg Carnot-Carathéodory metric. However, this is not yet
enough to finish our argument. We still need to show that the bottom and upper parts of $\partial E$ match together in the proper way. To do this, we write $E$ in the form

$$
\begin{equation*}
E=\left\{(x, y, t) \in \mathbb{H}^{1} \mid(y, t) \in F, h(y, t) \leq x \leq k(y, t)\right\}, \tag{1.10}
\end{equation*}
$$

for some compact convex set $F \subset \mathbb{R}^{2}$ and convex functions $h,-k: F \rightarrow \mathbb{R}$. The analysis carried out for $f$ can be also carried out for $h$ even though in this case computations are, unfortunately, more complicated. This provides the last piece of information needed to get the set $E_{\text {isop }}$ in (1.4).

A short overview is now in order. In Section 2, we derive the curvature equations for convex isoperimetric sets. In Section 3, we prove the generalized version of Theorem 1.2 which is needed in Section 4, where we establish the foliation by geodesics property for convex isoperimetric sets. Finally, in Section 5 we prove Theorem 1.1. The results concerning convex sets are collected in the appendix at the end of the paper.

## 2. Curvature equations for convex isoperimetric sets

We derive partial differential equations for certain vector fields built from the functions which parameterize the boundary of convex isoperimetric sets. We study graphs of the form $t=f(x, y)$ and graphs of the form $x=h(y, t)$. We use the following representation formula for perimeter. If $E \subset \mathbb{H}^{1}$ is a bounded set such that $\partial E$ is locally a Lipschitz surface, then

$$
\begin{equation*}
P(E)=\int_{\partial E} \sqrt{(X \cdot v)^{2}+(Y \cdot v)^{2}} d \mathcal{H}^{2} \tag{2.1}
\end{equation*}
$$

where $v$ is a unit normal to $\partial E$. Here and in the following, denotes the standard inner product in $\mathbb{R}^{3}$ or $\mathbb{R}^{2}$ (we think of $X$ and $Y$ as vectors in $\mathbb{R}^{3}$ ). $\mathcal{H}^{2}$ is the 2dimensional Hausdorff measure in $\mathbb{R}^{3}$ with respect to the usual Euclidean distance. For a proof of formula (2.1), see [9].

Graphs of the form $t=f(x, y)$
We denote elements of $\mathbb{H}^{1}=\mathbb{R}^{2} \times \mathbb{R}$ by $(z, t)$ with $t \in \mathbb{R}$ and $z=(x, y) \in \mathbb{R}^{2}$. We write $z^{\perp}=(-y, x)$. Let $E$ be a convex set in $\mathbb{H}^{1}$ of the form (1.5). We define the characteristic set of the convex function $f: D \rightarrow \mathbb{R}$ as

$$
\begin{equation*}
\Sigma(f)=\left\{z \in \operatorname{int}(D) \mid-2 z^{\perp} \in \partial f(z)\right\} \tag{2.2}
\end{equation*}
$$

where $\partial f(z)$ stands for the subdifferential of $f$ at $z$.
Proposition 2.1 (Curvature equation I). Let $E \subset \mathbb{H}^{1}$ be a convex isoperimetric set with perimeter $P(E)$ and volume $|E|$. Then the function $f: D \rightarrow \mathbb{R}$ satisfies in distributional sense the partial differential equation

$$
\begin{equation*}
\operatorname{div}\left(\frac{\nabla f(z)+2 z^{\perp}}{\left|\nabla f(z)+2 z^{\perp}\right|}\right)=H \tag{2.3}
\end{equation*}
$$

in $\operatorname{int}(D) \backslash \Sigma(f)$, with $H=3 P(E) / 4|E|$.

Proof. By Theorem A. 1 in the appendix, $\Sigma(f)$ is contained in the union of at most two line segments.

Let $\varphi \in C_{c}^{\infty}(\operatorname{int}(D) \backslash \Sigma(f))$ and for $\varepsilon \in \mathbb{R}$, consider the set

$$
E_{\varepsilon}=\left\{(z, t) \in \mathbb{H}^{1} \mid z \in D, f(z)+\varepsilon \varphi(z) \leq t \leq g(z)\right\} .
$$

We let $P(\varepsilon)=P\left(E_{\varepsilon}\right), V(\varepsilon)=\left|E_{\varepsilon}\right|$ and $R(\varepsilon)=P(\varepsilon)^{4} / V(\varepsilon)^{3}$. If $E$ is an isoperimetric set, then $R(\varepsilon)$ has a minimum at $\varepsilon=0$, and then

$$
\begin{equation*}
R^{\prime}(0)=\left.\frac{P^{3}}{V^{4}}\left(4 P^{\prime} V-3 P V^{\prime}\right)\right|_{\varepsilon=0}=0 \tag{2.4}
\end{equation*}
$$

There exists $\varepsilon_{0}>0$ such that

$$
\begin{equation*}
V^{\prime}(\varepsilon)=-\int_{D} \varphi(z) d z \quad \text { for }|\varepsilon|<\varepsilon_{0} \tag{2.5}
\end{equation*}
$$

Let $v_{\varepsilon}$ be the exterior unit normal to $\partial E_{\varepsilon}$ and let $S_{\varepsilon}=\left\{(z, f(z)+\varepsilon \varphi(z)) \in \mathbb{H}^{1} \mid z \in\right.$ $D\}$ be the graph of $f+\varepsilon \varphi$. By the Heisenberg area formula (2.1) and from the standard area formula for graphs of functions in Euclidean spaces, we find

$$
\begin{align*}
P^{\prime}(\varepsilon) & =\frac{d}{d \varepsilon} \int_{S_{\varepsilon}} \sqrt{\left(X \cdot v_{\varepsilon}\right)^{2}+\left(Y \cdot v_{\varepsilon}\right)^{2}} d \mathcal{H}^{2}  \tag{2.6}\\
& =\frac{d}{d \varepsilon} \int_{D}\left|\nabla f(z)+\varepsilon \nabla \varphi(z)+2 z^{\perp}\right| d z
\end{align*}
$$

By Proposition A. 2 in the appendix, for any compact set $K \subset \operatorname{int}(D) \backslash \Sigma(f)$, there is $\delta>0$ such that $\left|\nabla f(z)+2 z^{\perp}\right| \geq \delta$ for $\mathcal{L}^{2}$-a.e. $z \in K$. Then we can interchange derivative and integral for all small $\varepsilon$. At $\varepsilon=0$ we obtain

$$
\begin{equation*}
P^{\prime}(0)=\int_{D} \frac{\nabla f(z)+2 z^{\perp}}{\left|\nabla f(z)+2 z^{\perp}\right|} \cdot \nabla \varphi(z) d z \tag{2.7}
\end{equation*}
$$

From (2.4), (2.5) and (2.7) we find

$$
4|E| \int_{D} \frac{\nabla f(z)+2 z^{\perp}}{\left|\nabla f(z)+2 z^{\perp}\right|} \cdot \nabla \varphi(z) d z=-3 P(E) \int_{D} \varphi(z) d z
$$

for arbitrary $\varphi \in C_{c}^{\infty}(\operatorname{int}(D) \backslash \Sigma(f))$. This is (2.3).

## Graphs of the form $x=h(y, t)$

We denote points in $\mathbb{H}^{1}=\mathbb{R} \times \mathbb{R}^{2}$ by $(x, \zeta)$ with $x \in \mathbb{R}$ and $\zeta=(y, t) \in \mathbb{R}^{2}$. Let $E$ be a convex set in $\mathbb{H}^{1}$ of the form (1.10). We define the characteristic set of the convex function $h: F \rightarrow \mathbb{R}$ as the set $\Sigma(h)$ of the points $\zeta \in F$ such that the horizontal plane spanned by the vector fields $X$ and $Y$ at the point $p=(h(\zeta), \zeta) \in$ $\mathbb{R} \times \mathbb{R}^{2}$ is a supporting plane for $E$ (i.e. the plane does not intersect the interior of $E)$. By Theorem A. 1 in the appendix, $\Sigma(h)$ is contained in the union of at most two line segments.

Proposition 2.2 (Curvature equation II). Let $E \subset \mathbb{H}^{1}$ be a convex isoperimetric set with perimeter $P(E)$ and volume $|E|$. Then the function $h: F \rightarrow \mathbb{R}$ satisfies in distributional sense the partial differential equation

$$
\begin{equation*}
\left(\partial_{y}-2 h \partial_{t}\right)\left(\frac{u_{1}}{|u|}\right)-2 y \partial_{t}\left(\frac{u_{2}}{|u|}\right)=H \tag{2.8}
\end{equation*}
$$

in $\operatorname{int}(F) \backslash \Sigma(h)$, with $H=3 P(E) / 4|E|$ and $u=\left(u_{1}, u_{2}\right)=\left(h_{y}-2 h h_{t}, 1-2 y h_{t}\right)$.
Proof. Notice that $u \in B V_{\text {loc }}\left(\operatorname{int}(F) ; \mathbb{R}^{2}\right)$ and moreover, by Proposition A. 3 in the appendix, for each compact set $K \subset(\operatorname{int}(F) \backslash \Sigma(h))$, there exists $\delta>0$ such that $|u| \geq \delta \mathcal{L}^{2}$-a.e. in $K$. Hence $u /|u| \in B V_{\text {loc }}\left(\operatorname{int}(F) \backslash \Sigma(h) ; \mathbb{R}^{2}\right)$.

Let $\varphi \in C_{c}^{\infty}(\operatorname{int}(F) \backslash \Sigma(h))$ and for any $\varepsilon \in \mathbb{R}$ consider the set

$$
E_{\varepsilon}=\left\{(x, \zeta) \in \mathbb{H}^{1} \mid \zeta \in F, h(\zeta)+\varepsilon \varphi(\zeta) \leq x \leq k(\zeta)\right\}
$$

We let $P(\varepsilon)=P\left(E_{\varepsilon}\right), V(\varepsilon)=\left|E_{\varepsilon}\right|$ and $R(\varepsilon)=P(\varepsilon)^{4} / V(\varepsilon)^{3}$. Denoting by $S_{\varepsilon}=\left\{(h(\zeta)+\varepsilon \varphi(\zeta), \zeta) \in \mathbb{H}^{1} \mid \zeta \in F\right\}$ the graph of $h+\varepsilon \varphi$ and by $\nu_{\varepsilon}$ the exterior unit normal to $\partial E_{\varepsilon}$, from the Heisenberg area formula (2.1) and from the standard area formula we get

$$
\begin{aligned}
P^{\prime}(\varepsilon) & =\frac{d}{d \varepsilon} \int_{S_{\varepsilon}} \sqrt{\left(X \cdot v_{\varepsilon}\right)^{2}+\left(Y \cdot v_{\varepsilon}\right)^{2}} d \mathcal{H}^{2} \\
& =\frac{d}{d \varepsilon} \int_{F}\left(\left(1-2 y\left(h_{t}+\varepsilon \varphi_{t}\right)\right)^{2}+\left(\left(h_{y}+\varepsilon \varphi_{y}\right)-2(h+\varepsilon \varphi)\left(h_{t}+\varepsilon \varphi_{t}\right)\right)^{2}\right)^{1 / 2} d \zeta
\end{aligned}
$$

Because we have $\left(h_{y}-2 h h_{t}\right)^{2}+\left(1-2 y h_{t}\right)^{2} \geq \delta^{2}>0 \mathcal{L}^{2}$-a.e. in a neighbourhood of $\operatorname{spt}(\varphi)$, we can differentiate under the integral sign and we obtain

$$
P^{\prime}(0)=\int_{F} \frac{\left(h_{y}-2 h h_{t}\right)\left(\varphi_{y}-2(\varphi h)_{t}\right)-2 y \varphi_{t}\left(1-2 y h_{t}\right)}{\sqrt{\left(h_{y}-2 h h_{t}\right)^{2}+\left(1-2 y h_{t}\right)^{2}}} d \zeta
$$

at $\varepsilon=0$. An integration by parts yields

$$
P^{\prime}(0)=-\int_{F} \varphi d\left(\left(\partial_{y}-2 h \partial_{t}\right)\left(\frac{u_{1}}{|u|}\right)-2 y \partial_{t}\left(\frac{u_{2}}{|u|}\right)\right)
$$

where the derivatives of $u /|u|$ are measures. If $E$ minimizes the isoperimetric ratio then, as in the proof of Proposition 2.1, we get

$$
4|E| \int_{F} \varphi d\left(\left(\partial_{y}-2 h \partial_{t}\right)\left(\frac{u_{1}}{|u|}\right)-2 y \partial_{t}\left(\frac{u_{2}}{|u|}\right)\right)=3 P(E) \int_{F} \varphi d \zeta
$$

for arbitrary $\varphi \in C_{c}^{\infty}(\operatorname{int}(F) \backslash \Sigma(h))$. This is (2.8).

## 3. Improved regularity of the boundary

In this section, we prove a regularity result for vector fields with bounded variation in the plane arising from the parameterization of the boundary of convex isoperimetric sets.

Let $\Omega \subset \mathbb{R}^{2}$ be an open set and let $a, b \in C(\Omega)$ be continuous functions. We consider the differential operator $\mathcal{M}$ acting on $B V\left(\Omega ; \mathbb{R}^{2}\right)$ defined by

$$
\begin{equation*}
\mathcal{M} u=\left(\partial_{1}-a \partial_{2}\right) u_{1}+b \partial_{2} u_{2} \tag{3.1}
\end{equation*}
$$

In general, $\mathcal{M} u$ is a measure in $\Omega$. We have $\mathcal{M}=\operatorname{div}$ when $a=0$ and $b=1$. When $a=2 h$ and $b=-2 y, \mathcal{M}$ is the operator appearing in the left hand side of (2.8).

We recall that a vector field $u \in L^{1}\left(\Omega ; \mathbb{R}^{2}\right)$ has approximate limit $\bar{u}(q) \in \mathbb{R}^{2}$ at the point $q \in \Omega$ if

$$
\begin{equation*}
\lim _{r \downarrow 0} f_{B(q, r)}|u-\bar{u}(q)| d \mathcal{L}^{2}=0 \tag{3.2}
\end{equation*}
$$

Here, $B(q, r)$ denotes the open Euclidean ball of radius $r$ centered at $q$. The approximate discontinuity set of $u$ is the set $S_{u}$ of points in $\Omega$ at which $u$ has no approximate limit. The jump set of $u$ is the set $J_{u}$ of points $q \in \Omega$ for which there exist $u^{+}(q), u^{-}(q) \in \mathbb{R}^{2}$ and $v_{u}(q) \in S^{1}$ such that $u^{+}(q) \neq u^{-}(q)$ and

$$
\begin{equation*}
\lim _{r \downarrow 0} f_{B^{ \pm}(q, r)}\left|u-u^{ \pm}(q)\right| d \mathcal{L}^{2}=0 \tag{3.3}
\end{equation*}
$$

where $B^{ \pm}(q, r)=\left\{q^{\prime} \in B(q, r) \mid \pm\left(q^{\prime}-q\right) \cdot v_{u}(q)>0\right\}$. Finally, the precise representative $u^{*}: \Omega \rightarrow \mathbb{R}^{2}$ of $u$ is:

$$
u^{*}(q)= \begin{cases}\bar{u}(q) & q \in \Omega \backslash S_{u}  \tag{3.4}\\ \frac{1}{2}\left(u^{+}(q)+u^{-}(q)\right) & q \in J_{u} \\ 0 & q \in S_{u} \backslash J_{u}\end{cases}
$$

By the Lebesgue density theorem, it is $u=u^{*} \mathcal{L}^{2}$-a.e. in $\Omega$.
Theorem 3.1 (Improved regularity). Let $\Omega \subset \mathbb{R}^{2}$ be a bounded open set, let $u=$ $\left(u_{1}, u_{2}\right) \in B V\left(\Omega ; \mathbb{R}^{2}\right)$, and let $a, b \in C(\Omega)$ be continuous functions such that $b \neq 0$ in $\Omega$. Assume that:
(i) There exists $\delta>0$ such that $|u| \geq \delta \mathcal{L}^{2}$-a.e. in $\Omega$;
(ii) $\mathcal{M} u^{\perp} \in L^{1}(\Omega)$, where $u^{\perp}=\left(-u_{2}, u_{1}\right)$;
(iii) $\mathcal{M}\left(\frac{u}{|u|}\right) \in L^{1}(\Omega)$.

Then $u /|u| \in W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$ and there exists a function $\mu: J_{u} \rightarrow(0,+\infty)$, such that $u^{-}=\mu u^{+}$on $J_{u}$.

Proof. If $u \in B V\left(\Omega ; \mathbb{R}^{2}\right)$ the measure $D u$ has the decomposition $D u=D^{a} u+$ $D^{s} u$ with $D^{a} u \ll \mathcal{L}^{2}$ and $D^{s} u \perp \mathcal{L}^{2} . D^{a} u=\nabla u \mathcal{L}^{2}$ is the absolutely continuous part of the measure, with $\nabla u \in L^{1}\left(\Omega ; M^{2 \times 2}\right)$ and $M^{2 \times 2}$ denotes the space of $2 \times 2$ matrices with real entries. Moreover, it is $D^{s} u=D^{c} u+D^{j} u$, where $D^{c} u=$ $D^{s} u \downharpoonright \Omega \backslash S_{u}$ is the Cantor part and $D^{j} u=D^{s} u\left\llcorner S_{u}\right.$ is the jump part. By the representation theorem of Federer-Vol'pert and by Alberti's rank one theorem, the measures $D^{c} u$ and $D^{j} u$ admit the representations

$$
\begin{equation*}
D^{j} u=\left(u^{+}-u^{-}\right) \otimes v_{u} \mathcal{H}^{1}\left\llcorner J_{u} \quad \text { and } \quad D^{c} u=\eta \otimes \xi\left|D^{c} u\right|\right. \tag{3.5}
\end{equation*}
$$

where $\left|D^{c} u\right|$ is the total variation of $D^{c} u$ and $\eta, \xi: \Omega \rightarrow S^{1}$ are suitable Borel maps. The measure $\left|D^{c} u\right|$ is absolutely continuous with respect to $\mathcal{H}^{1}$. The set $J_{u}$ is an $\mathcal{H}^{1}$-rectifiable Borel subset of $S_{u}$ with $\mathcal{H}^{1}\left(S_{u} \backslash J_{u}\right)=0$. For these and related results on $B V$ vector fields we refer the reader to the monograph [2].

We claim that both the jump and the Cantor parts of $D(u /|u|)$ vanish. Let $F: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$ be a smooth mapping such that

$$
\begin{equation*}
F(0)=0, \quad F(q)=\frac{q}{|q|} \text { for }|q| \geq \frac{\delta}{2} \quad \text { and } \quad|\nabla F| \in L^{\infty}\left(\mathbb{R}^{2}\right) \tag{3.6}
\end{equation*}
$$

If $|q|>\delta / 2$, the derivative of $F$ is

$$
\begin{equation*}
\nabla F(q)=\frac{1}{|q|^{3}} q^{\perp} \otimes q^{\perp} \tag{3.7}
\end{equation*}
$$

By (i), the vector field $v=F \circ u$ is in $B V\left(\Omega ; \mathbb{R}^{2}\right)$ (this fact is implicitly assumed in (iii)).

By the chain rule for $B V$ vector fields, we have

$$
\begin{equation*}
D^{a} v=\nabla F\left(u^{*}\right) \nabla u \mathcal{L}^{2} \quad \text { and } \quad D^{c} v=\left(\nabla F\left(u^{*}\right) \eta\right) \otimes \xi\left|D^{c} u\right| . \tag{3.8}
\end{equation*}
$$

A point $q \in \Omega$ belongs to $J_{v}$ if and only if $q \in J_{u}$ and $F\left(u^{+}(q)\right) \neq F\left(u^{-}(q)\right)$. Moreover, it is $D^{j} v=\left(F\left(u^{+}\right)-F\left(u^{-}\right)\right) \otimes v_{u} \mathcal{H}^{1}\left\llcorner J_{u}\right.$. Notice that $\left|u^{ \pm}(q)\right| \geq \delta$ for all $q \in J_{u}$, by (i). Then, the jump set of $v$ is

$$
\begin{equation*}
J_{v}=\left\{q \in J_{u}\left|u^{+}(q) /\left|u^{+}(q)\right| \neq u^{-}(q) /\left|u^{-}(q)\right|\right\}\right. \tag{3.9}
\end{equation*}
$$

and the jump part of $D v$ is

$$
\begin{equation*}
D^{j} v=\left(\frac{u^{+}}{\left|u^{+}\right|}-\frac{u^{-}}{\left|u^{-}\right|}\right) \otimes v_{u} \mathcal{H}^{1}\left\llcorner J_{u}\right. \tag{3.10}
\end{equation*}
$$

By assumption (ii), the part of the measure $-\partial_{1} u_{2}+b \partial_{2} u_{1}+a \partial_{2} u_{2}$ concentrated on $J_{u}$ vanishes. From the formula for $D^{j} u$ in (3.5), we can compute $D_{k}^{j} u_{l}$ for $k, l=1,2$, and we obtain

$$
\begin{equation*}
\left(\left(u_{2}^{-}-u_{2}^{+}\right), b\left(u_{1}^{+}-u_{1}^{-}\right)+a\left(u_{2}^{+}-u_{2}^{-}\right)\right) \cdot v_{u}=0 \tag{3.11}
\end{equation*}
$$

$\mathcal{H}^{1}$-a.e. on $J_{u}$.
By assumption (iii), the part of the measure $\left(\partial_{1}-a \partial_{2}\right)\left(u_{1} /|u|\right)+b \partial_{2}\left(u_{2} /|u|\right)$ concentrated on $J_{v}$ vanishes. Thus, using (3.10), we can compute $D_{k}^{j}\left(u_{l} /|u|\right)$ for $k, l=1,2$, and we get

$$
\begin{equation*}
\left(\frac{u_{1}^{+}}{\left|u^{+}\right|}-\frac{u_{1}^{-}}{\left|u^{-}\right|}, \frac{a u_{1}^{-}-b u_{2}^{-}}{\left|u^{-}\right|}-\frac{a u_{1}^{+}-b u_{2}^{+}}{\left|u^{+}\right|}\right) \cdot v_{u}=0 \tag{3.12}
\end{equation*}
$$

$\mathcal{H}^{1}$-a.e. on $J_{v}$.
From (3.12) and (3.11), we deduce that there exists $\lambda \in \mathbb{R}$ such that the following system of equations is satisfied $\mathcal{H}^{1}$-a.e. on $J_{v}$ :

$$
\left\{\begin{array}{l}
\frac{u_{1}^{+}}{\left|u^{+}\right|}-\frac{u_{1}^{-}}{\left|u^{-}\right|}=-\lambda\left(u_{2}^{+}-u_{2}^{-}\right) \\
\frac{b u_{2}^{+}-a u_{1}^{+}}{\left|u^{+}\right|}-\frac{b u_{2}^{-}-a u_{1}^{-}}{\left|u^{-}\right|}=\lambda\left(b\left(u_{1}^{+}-u_{1}^{-}\right)+a\left(u_{2}^{+}-u_{2}^{-}\right)\right)
\end{array}\right.
$$

By elementary linear algebra, using $b \neq 0$, we obtain the equivalent system

$$
\left\{\begin{array}{l}
\frac{u_{1}^{+}}{\left|u^{+}\right|}-\frac{u_{1}^{-}}{\left|u^{-}\right|}=-\lambda\left(u_{2}^{+}-u_{2}^{-}\right) \\
\frac{u_{2}^{+}}{\left|u^{+}\right|}-\frac{u_{2}^{-}}{\left|u^{-}\right|}=\lambda\left(u_{1}^{+}-u_{1}^{-}\right)
\end{array} \quad \Leftrightarrow \quad\left(\frac{u^{+}}{\left|u^{+}\right|}-\frac{u^{-}}{\left|u^{-}\right|}\right) \cdot\left(u^{+}-u^{-}\right)=0\right.
$$

that is $u^{-} \cdot u^{+}=\left|u^{-} \| u^{+}\right|, \mathcal{H}^{1}$-a.e. on $J_{v}$. It follows that $u^{-}=\mu u^{+} \mathcal{H}^{1}$-a.e. on $J_{v}$ for some $\mu: J_{v} \rightarrow(0,+\infty)$, and then $u^{+} /\left|u^{+}\right|=u^{-} /\left|u^{-}\right| \mathcal{H}^{1}$-a.e. on $J_{v}$. This proves that $\mathcal{H}^{1}\left(J_{v}\right)=0$ and thus $D^{j} v=0$.

By assumption (ii), the Cantor part of the measure $-\partial_{1} u_{2}+b \partial_{2} u_{1}+a \partial_{2} u_{2}$ vanishes. From $D^{c} u=\eta \otimes \xi\left|D^{c} u\right|$, we can compute $D_{k}^{c} u_{l}$ for $k, l=1,2$, and we find

$$
\begin{equation*}
\left(-\eta_{2}, b \eta_{1}+a \eta_{2}\right) \cdot \xi=0 \tag{3.13}
\end{equation*}
$$

$\left|D^{c} u\right|$-a.e. on $\Omega$. By assumption (iii), the Cantor part of the measure

$$
\left(\partial_{1}-a \partial_{2}\right)\left(u_{1} /|u|\right)+b \partial_{2}\left(u_{2} /|u|\right)
$$

vanishes, too. Thus, letting $\vartheta=\left(\vartheta_{1}, \vartheta_{2}\right)=\left(\left(u^{*}\right)^{\perp} \otimes\left(u^{*}\right)^{\perp}\right) \eta$, we can use (3.7) and (3.8) to compute $D_{k}^{c}\left(u_{l} /|u|\right)$ for $k, l=1,2$, and we find

$$
\begin{equation*}
\left(\vartheta_{1}, b \vartheta_{2}-a \vartheta_{1}\right) \cdot \xi=0 \tag{3.14}
\end{equation*}
$$

$\left|D^{c} u\right|$-a.e. on $\Omega$. From (3.13) and (3.14), we deduce that there exists $\lambda \in \mathbb{R}$ such that

$$
\left\{\begin{array} { l } 
{ \vartheta _ { 1 } = - \lambda \eta _ { 2 } } \\
{ b \vartheta _ { 2 } - a \vartheta _ { 1 } = \lambda ( b \eta _ { 1 } + a \eta _ { 2 } ) }
\end{array} \Leftrightarrow \left\{\begin{array}{l}
\vartheta_{1}=-\lambda \eta_{2} \\
\vartheta_{2}=\lambda \eta_{1}
\end{array}\right.\right.
$$

Here we used $b \neq 0$. This, in turn, is equivalent with

$$
\vartheta \cdot \eta=\left(\left(\left(u^{*}\right)^{\perp} \otimes\left(u^{*}\right)^{\perp}\right) \eta\right) \cdot \eta=0
$$

$\left|D^{c} u\right|$-a.e. on $\Omega$. Using the identity $\left(\left(\left(u^{*}\right)^{\perp} \otimes\left(u^{*}\right)^{\perp}\right) \eta\right) \cdot \eta=\left(\left(u^{*}\right)^{\perp} \cdot \eta\right)^{2}$, we deduce that $\left(u^{*}\right)^{\perp} \cdot \eta=0$, and thus $\left(\left(u^{*}\right)^{\perp} \otimes\left(u^{*}\right)^{\perp}\right) \eta=0$ as well. By (3.7), this proves that $\left(\nabla F\left(u^{*}\right) \eta\right) \otimes \xi=0\left|D^{c} u\right|$-a.e. on $\Omega$, whence $D^{c} v=0$ by (3.8).

Theorem 3.1 has the following corollaries.
Corollary 3.2. Let $E \subset \mathbb{H}^{1}$ be a convex isoperimetric set and let $f: D \rightarrow \mathbb{R}$ be the function in (1.5). Then we have

$$
\begin{equation*}
\frac{\nabla f(z)+2 z^{\perp}}{\left|\nabla f(z)+2 z^{\perp}\right|} \in W_{\mathrm{loc}}^{1,1}\left(\operatorname{int}(D) \backslash \Sigma(f) ; \mathbb{R}^{2}\right) \tag{3.15}
\end{equation*}
$$

Proof. The vector field $u(z)=\left(u_{1}(z), u_{2}(z)\right)=\nabla f(z)+2 z^{\perp}$ satisfies div $u^{\perp}=$ $-4 \operatorname{in} \operatorname{int}(D)$. By Proposition A. 2 in the appendix, $|u| \geq \delta>0$ on a compact set in $\operatorname{int}(D) \backslash \Sigma(f)$. Moreover, by Proposition 2.1, we have $\operatorname{div}(u /|u|)=3 P(E) / 4|E|$ in $\operatorname{int}(D) \backslash \Sigma(f)$. The claim follows from Theorem 3.1 with $a=0$ and $b=1$.

Corollary 3.3. Let $E \subset \mathbb{H}^{1}$ be a convex isoperimetric set and let $h: F \rightarrow \mathbb{R}$ be the function in (1.10). Then we have

$$
\begin{equation*}
\frac{u}{|u|} \in W_{\operatorname{loc}}^{1,1}\left(\operatorname{int}(F) \backslash(\{y=0\} \cup \Sigma(h)) ; \mathbb{R}^{2}\right), \tag{3.16}
\end{equation*}
$$

with $u=\left(u_{1}, u_{2}\right)=\left(h_{y}-2 h h_{t}, 1-2 y h_{t}\right)$.
Proof. We use Theorem 3.1 with $a=2 h$ and $b=-2 y$. The vector field $u$ is in $B V_{\text {loc }}\left(\operatorname{int}(F) ; \mathbb{R}^{2}\right)$ and satisfies
$\mathcal{M} u^{\perp}=\left(\partial_{y}-2 h \partial_{t}\right)\left(2 y h_{t}-1\right)-2 y \partial_{t}\left(h_{y}-2 h h_{t}\right)=2 h_{t}\left(1+2 y h_{t}\right) \in L_{\text {loc }}^{\infty}(\operatorname{int}(F))$.
From Proposition A. 3 in the appendix, it follows that $u /|u| \in B V_{\text {loc }}(\operatorname{int}(F) \backslash \Sigma(h))$, and moreover, by Proposition 2.2, $u /|u|$ satisfies

$$
\mathcal{M}\left(\frac{u}{|u|}\right)=\left(\partial_{y}-2 h \partial_{t}\right)\left(\frac{u_{1}}{|u|}\right)-2 y \partial_{t}\left(\frac{u_{2}}{|u|}\right)=\frac{3 P(E)}{4|E|}
$$

in $\operatorname{int}(F) \backslash \Sigma(h)$. Now the claim follows from Theorem 3.1.

## 4. Integration of the curvature equation

In this section we solve equations (2.3) and (2.8). These equations can be integrated along a regular Lagrangian flow and the solutions are suitable arcs of circle.

We recall the definition of regular flow in the plane. Let $\Omega \subset \mathbb{R}^{2}$ be an open set and let $u \in B V\left(\Omega ; \mathbb{R}^{2}\right) \cap L^{\infty}\left(\Omega ; \mathbb{R}^{2}\right)$. For $q \in \Omega$ and $\varrho>0$ small enough define

$$
\begin{aligned}
& C_{q}([-\varrho, \varrho] ; \Omega) \\
& =\left\{\gamma \in C([-\varrho, \varrho] ; \Omega) \mid \gamma(s)=q+\int_{0}^{s} u(\gamma(\sigma)) d \sigma, s \in[-\varrho, \varrho]\right\} .
\end{aligned}
$$

Let $K \subset \Omega$ be a compact set. An $\mathcal{L}^{2}$-measurable map $\Phi: K \rightarrow C([-\varrho, \varrho] ; \Omega)$ is a Lagrangian flow starting from $K$ relative to $u$ if $\Phi(q) \in C_{q}([-\varrho, \varrho] ; \Omega)$ for $\mathcal{L}^{2}$ a.e. $q \in K$. With abuse of notation, $\Phi$ is identified with the map $\Phi: K \times[-\varrho, \varrho] \rightarrow$ $\Omega, \Phi(q, s)=\Phi(q)(s)$. The flow is said to be regular if there exists a constant $m \geq 1$ such that

$$
\begin{equation*}
\frac{1}{m} \mathcal{L}^{2}(A) \leq \mathcal{L}^{2}(\Phi(A, s)) \leq m \mathcal{L}^{2}(A) \tag{4.1}
\end{equation*}
$$

for all $\mathcal{L}^{2}$-measurable sets $A \subset K$ and for all $s \in[-\varrho, \varrho]$.
Theorem 4.1 (Foliation by circles I). Let $E \subset \mathbb{H}^{1}$ be a convex isoperimetric set and let $f: D \rightarrow \mathbb{R}$ be the function in (1.5). Let $K \subset \operatorname{int}(D) \backslash \Sigma(f)$ be a compact set and $\Omega \subset \operatorname{int}(D) \backslash \Sigma(f)$ an open neighborhood of $K$.

For a sufficiently small $\varrho>0$, there exists a (unique) regular Lagrangian flow $\Phi: K \times[-\varrho, \varrho] \rightarrow \Omega$ of the vector field $v(z)=2 z-\nabla f^{\perp}(z)$. Moreover, for $\mathcal{L}^{2}$-a.e. $z \in K$, the integral curve $s \mapsto \Phi(z, s)$ is an arc of circle with radius $4|E| / 3 P(E)$ oriented clockwise.

For a convex isoperimetric set of the form (1.10) there is an analogous statement.

Theorem 4.2 (Foliation by circles II). Let $E \subset \mathbb{H}^{1}$ be a convex isoperimetric set and let $h: F \rightarrow \mathbb{R}$ be the function in (1.10). Let $K \subset \operatorname{int}(F) \backslash(\{y=0\} \cup \Sigma(h))$ be a compact set and $\Omega \subset \operatorname{int}(F) \backslash(\{y=0\} \cup \Sigma(h))$ be an open neighborhood of $K$.

For a sufficiently small $\varrho>0$, there exists a (unique) regular Lagrangian flow $\Phi: K \times[-\varrho, \varrho] \rightarrow \Omega$ of the vector field $v(\zeta)=\left(1-2 y h_{t}, 2 y h_{y}-2 h\right)$, $\zeta=(y, t) \in F$. Moreover, for $\mathcal{L}^{2}$-a.e. $\zeta \in K$, the projection onto the xy-plane of the curve

$$
\begin{equation*}
s \mapsto(h(\Phi(\zeta, s)), \Phi(\zeta, s)) \in \mathbb{H}^{1}, \quad s \in[-\varrho, \varrho] \tag{4.2}
\end{equation*}
$$

is an arc of circle with radius $4|E| / 3 P(E)$ oriented clockwise.
The existence of a (unique) regular Lagrangian flow stated in Theorems 4.1 and 4.2 follows from Ambrosio's theory on the Cauchy problem for $B V$ vector fields.

Theorem 4.3 (Ambrosio). Let $\Omega \subset \mathbb{R}^{2}$ be an open set and assume that:
(i) $u \in B V\left(\Omega ; \mathbb{R}^{2}\right) \cap L^{\infty}\left(\Omega ; \mathbb{R}^{2}\right)$;
(ii) $\operatorname{div} u \in L^{\infty}(\Omega)$.

Then, for any compact set $K \subset \Omega$ and for a small enough $\varrho>0$, there exists $a$ (unique) regular Lagrangian flow $\Phi: K \times[-\varrho, \varrho] \rightarrow \Omega$ starting from $K$ relative to $u$.

The existence statement is Theorem 6.2 and the uniqueness statement is in [1, Theorem 6.4]. We do not need the uniqueness of the flow in our argument. Ambrosio's theory holds more generally for non autonomous vector fields in any space dimension.

The characterization of the flow in Theorems 4.1 and 4.2 is obtained on proving that a suitable reparameterization $\gamma$ of a generic integral curve of the flow has a second order derivative which satisfies the equation $\ddot{\gamma}=-H \dot{\gamma}^{\perp}$ in a weak sense with $H=3 P(E) / 4|E|$.

In order to make this reparameterization argument precise, let us consider an open set $\Omega \subset \mathbb{R}^{2}$ and a vector field $v$ in $\Omega$ satisfying the assumptions of Theorem 4.3, that is

$$
\begin{equation*}
v \in B V\left(\Omega ; \mathbb{R}^{2}\right) \cap L^{\infty}\left(\Omega ; \mathbb{R}^{2}\right) \quad \text { and } \quad \operatorname{div} v \in L^{\infty}(\Omega) \tag{4.3}
\end{equation*}
$$

The function $v$ is defined pointwise, i.e. we choose a representative in the equivalence class of $v$. Our results hold independently of this choice. However, there is an exceptional set of points which may a priori depend on the representative.

Given a compact set $K \subset \Omega$ and a sufficiently small $\varrho>0$, there exists a unique regular Lagrangian flow $\Phi: K \times[-\varrho, \varrho] \rightarrow \Omega$ starting from $K$ relative to $v$. Let $\lambda: \Omega \rightarrow \mathbb{R}$ be a measurable function such that

$$
\begin{equation*}
0<c_{1} \leq \lambda \leq c_{2} \quad \mathcal{L}^{2} \text {-a.e. in } \Omega \tag{4.4}
\end{equation*}
$$

Then, for $\mathcal{L}^{2}$-a.e. $q \in K$, the curve $s \mapsto \lambda(\Phi(q, s))$ is measurable and

$$
c_{1} \leq \lambda(\Phi(q, s)) \leq c_{2} \quad \text { for a.e. } s \in[-\varrho, \varrho] .
$$

This follows from (4.1) by a Fubini-type argument. In fact, if $N \subset \Omega$ is an $\mathcal{L}^{2}$ negligible set, then $\Phi^{-1}(N) \subset K \times[-\varrho, \varrho]$ is also negligible. Thus, for $\mathcal{L}^{2}$ a.e. $q \in K$, the change of parameter $\sigma_{q}:[-\varrho, \varrho] \rightarrow\left[\sigma_{q}(-\varrho), \sigma_{q}(\varrho)\right]$ defined by

$$
\sigma_{q}(s)=\int_{0}^{s} \lambda(\Phi(q, \xi)) d \xi
$$

is bi-Lipschitz, strictly increasing and admits therefore a bi-Lipschitz, strictly increasing inverse $\tau_{q}:\left[\sigma_{q}(-\varrho), \sigma_{q}(\varrho)\right] \rightarrow[-\varrho, \varrho]$, which satisfies

$$
\begin{equation*}
\dot{\tau}_{q}(s)=\frac{1}{\dot{\sigma}_{q}\left(\tau_{q}(s)\right)}=\frac{1}{\lambda\left(\Phi\left(q, \tau_{q}(s)\right)\right)} \quad \text { for a.e. } s \in\left[\sigma_{q}(-\varrho), \sigma_{q}(\varrho)\right] \tag{4.5}
\end{equation*}
$$

Consequently, for $\mathcal{L}^{2}$-a.e. $q \in K$, the curve $\gamma_{q}:\left[\sigma_{q}(-\varrho), \sigma_{q}(\varrho)\right] \rightarrow \Omega$ defined by

$$
\begin{equation*}
\gamma_{q}(s)=\Phi\left(q, \tau_{q}(s)\right) \tag{4.6}
\end{equation*}
$$

is absolutely continuous and satisfies

$$
\begin{equation*}
\dot{\gamma}_{q}(s)=\frac{v\left(\gamma_{q}(s)\right)}{\lambda\left(\gamma_{q}(s)\right)} \quad \text { for a.e. } s \in\left[\sigma_{q}(-\varrho), \sigma_{q}(\varrho)\right], \tag{4.7}
\end{equation*}
$$

i.e. $\gamma_{q}$ is an integral curve of the vector field $v / \lambda$.

Theorem 4.4 (Chain rule for integral curves). Let $\Omega \subset \mathbb{R}^{2}$ be an open set, $K \subset$ $\Omega$ be a compact subset, $w \in W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$ and $\varrho>0$ small enough. For $\mathcal{L}^{2}$-a.e. $q \in$ $K$, the curve $w \circ \gamma_{q}$ with $\gamma_{q}$ defined in (4.6) belongs to $W^{1,1}\left(\left(\sigma_{q}(-\varrho), \sigma_{q}(\varrho)\right) ; \mathbb{R}^{2}\right)$ and its weak derivative is $\left(\nabla w \circ \gamma_{q}\right) \dot{\gamma}_{q}$.

Proof. Let $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ be a sequence of smooth vector fields $w_{k} \in W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$ converging to $w$ in $W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$. The map $(q, s) \mapsto w(\Phi(q, s))$ is measurable and integrable on $K \times[-\varrho, \varrho]$. By Fubini's theorem and by the bounded volume distortion property (4.1) of the flow, we have

$$
\int_{K} \int_{-\varrho}^{\varrho}\left|w_{k}(\Phi(q, s))-w(\Phi(q, s))\right| d s d q \leq m \int_{-\varrho}^{\varrho} \int_{\Phi(K, s)}\left|w_{k}(q)-w(q)\right| d q d s
$$

and

$$
\begin{aligned}
& \int_{K} \int_{-\varrho}^{\varrho}\left|\left(\nabla w_{k}(\Phi(q, s))-\nabla w(\Phi(q, s))\right) v(\Phi(q, s))\right| d s d q \\
& \leq\|v\|_{\infty} \int_{-\varrho}^{\varrho} \int_{K}\left|\nabla w_{k}(\Phi(q, s))-\nabla w(\Phi(q, s))\right| d q d s \\
& \leq m\|v\|_{\infty} \int_{-\varrho}^{\varrho} \int_{\Phi(K, s)}\left|\nabla w_{k}(q)-\nabla w(q)\right| d q d s .
\end{aligned}
$$

The curve $s \mapsto w_{k}(\Phi(q, s))$ is Lipschitz, for all $k \in \mathbb{N}$ and for $\mathcal{L}^{2}$-a.e. $q \in K$, with derivative $s \mapsto \nabla w_{k}(\Phi(q, s)) v(\Phi(q, s))$. Passing to a subsequence and relabelling if necessary, we conclude that

$$
\lim _{k \rightarrow \infty} w_{k}(\Phi(q, \cdot))=w(\Phi(q, \cdot)) \quad \text { in } W^{1,1}\left((-\varrho, \varrho) ; \mathbb{R}^{2}\right)
$$

for $\mathcal{L}^{2}$-a.e. $q \in K$, and that the weak derivative of $w(\Phi(q, \cdot))$ is $\nabla w(\Phi(q, \cdot)) v(\Phi(q, \cdot))$. Combining this observation, (4.5) and (4.7), for any $\varphi \in C_{c}^{\infty}\left(\left(\sigma_{q}(-\varrho), \sigma_{q}(\varrho)\right) ; \mathbb{R}^{2}\right)$
we get,

$$
\begin{aligned}
\int_{\sigma_{q}(-\varrho)}^{\sigma_{q}(\varrho)} w\left(\gamma_{q}(s)\right) \dot{\varphi}(s) d s & =\int_{-\varrho}^{\varrho} w(\Phi(q, s)) \dot{\varphi}\left(\sigma_{q}(s)\right) \dot{\sigma}_{q}(s) d s \\
& =-\int_{-\varrho}^{\varrho} \nabla w(\Phi(q, s)) v(\Phi(q, s)) \varphi\left(\sigma_{q}(s)\right) d s \\
& =-\int_{\sigma_{q}(-\varrho)}^{\sigma_{q}(\varrho)} \nabla w\left(\gamma_{q}(s)\right) \dot{\gamma}_{q}(s) \varphi(s) d s
\end{aligned}
$$

Hence $w \circ \gamma_{q} \in W^{1,1}\left(\left(\sigma_{q}(-\varrho), \sigma_{q}(\varrho)\right) ; \mathbb{R}^{2}\right)$, and its weak derivative is $(\nabla w \circ$ $\left.\gamma_{q}\right) \dot{\gamma_{q}}$.

Proof of Theorem 4.1. Let $f: D \rightarrow \mathbb{R}$ be the convex function in (1.5). We consider the vector fields $v(z)=2 z-\nabla f^{\perp}(z)$ and $u(z)=v^{\perp}(z)=\nabla f(z)+2 z^{\perp}$ which are both in $B V_{\text {loc }}\left(\operatorname{int}(D) ; \mathbb{R}^{2}\right) \cap L_{\text {loc }}^{\infty}\left(\operatorname{int}(D) ; \mathbb{R}^{2}\right)$. We have

$$
\begin{equation*}
\operatorname{div} v=4+f_{y x}-f_{x y}=4 \tag{4.8}
\end{equation*}
$$

If $K \subset \operatorname{int}(D) \backslash \Sigma(f)$ is a compact set and $\Omega \Subset \operatorname{int}(D) \backslash \Sigma(f)$ is a suitable open neighborhood of $K$, the existence of a regular Lagrangian flow $\Phi: K \times[-\varrho, \varrho] \rightarrow$ $\Omega$ for some $\varrho>0$ is guaranteed by Theorem 4.3.

The function $\lambda: \Omega \rightarrow \mathbb{R}, \lambda=|v|$, satisfies $0<c_{1} \leq \lambda \leq c_{2} \mathcal{L}^{2}$-a.e. in $\Omega$ by Proposition A. 2 in the appendix. By Corollary 3.2, the vector field $w=v / \lambda$ belongs to $W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$, and $\operatorname{div} w^{\perp}=H \mathcal{L}^{2}$-a.e. in $\Omega$ by (2.3), with $H=3 P(E) / 4|E|$.

Let $\gamma_{z}:\left(\sigma_{z}(-\varrho), \sigma_{z}(\varrho)\right) \rightarrow \Omega$ be the integral curve of the vector field $w$ defined in (4.6) and satisfying (4.7). We claim that $\gamma_{z}$ parameterizes an arc of circle with curvature $H$. Since $\Phi(z, \cdot)$ is a reparameterization of $\gamma_{z}$, proving the claim concludes the proof of Theorem 4.1. The following identities hold a.e. in $\left(\sigma_{z}(-\varrho), \sigma_{z}(\varrho)\right)$ for $\mathcal{L}^{2}$-a.e. $z \in K$ :
(i) $\left|w \circ \gamma_{z}\right|=1$;
(ii) $\left(w \cdot \partial_{x} w\right) \circ \gamma_{z}=\left(w \cdot \partial_{y} w\right) \circ \gamma_{z}=0$;
(iii) $\left(\operatorname{div} w^{\perp}\right) \circ \gamma_{z}=H$.

Here, we are using the fact that if $F, G: \Omega \rightarrow \mathbb{R}$ are $\mathcal{L}^{2}$-measurable functions such that $F=G \mathcal{L}^{2}$-a.e. in $\Omega$, then it is $F \circ \Phi=G \circ \Phi \mathcal{L}^{3}$-a.e. in $K \times[-\varrho, \varrho]$, which is a consequence of (4.1).

By Theorem 4.4, it is $\gamma_{z} \in W^{2,1}\left(\left(\sigma_{z}(-\varrho), \sigma_{z}(\varrho)\right) ; \mathbb{R}^{2}\right)$ with $\ddot{\gamma}_{z}=\left(\nabla w \circ \gamma_{z}\right) \dot{\gamma}_{z}$. Using (ii) and (iii), we compute

$$
\ddot{\gamma}_{z} \cdot \dot{\gamma}_{z}^{\perp}=\left(-\operatorname{div} w^{\perp}\right) \circ \gamma_{z}=-H
$$

Moreover, by (i) we also have $\ddot{\gamma}_{z} \cdot \dot{\gamma}_{z}=0$ a.e. in $\left(\sigma_{z}(-\varrho), \sigma_{z}(\varrho)\right)$. Then $\ddot{\gamma}_{z}=-H \dot{\gamma}_{z}^{\perp}$ a.e. and this implies that $\gamma_{z}$ is of class $C^{\infty}$. The claim follows.

Proof of Theorem 4.2. Let $h: F \rightarrow \mathbb{R}$ be the convex function in (1.10). We use the variables $\zeta=(y, t) \in F$. The vector field $v: F \rightarrow \mathbb{R}^{2}$

$$
\begin{equation*}
v(\zeta)=\left(1-2 y h_{t}, 2 y h_{y}-2 h\right) \tag{4.9}
\end{equation*}
$$

is in $B V_{\mathrm{loc}}\left(\operatorname{int}(F) ; \mathbb{R}^{2}\right) \cap L_{\text {loc }}^{\infty}\left(\operatorname{int}(F) ; \mathbb{R}^{2}\right)$, and

$$
\begin{equation*}
\operatorname{div} v=-4 h_{t}-2 y h_{t y}+2 y h_{y t}=-4 h_{t} \in L_{\mathrm{loc}}^{\infty}(\operatorname{int}(F)) \tag{4.10}
\end{equation*}
$$

The vector field $v$ is the projection onto the $y t$-plane of the vector field

$$
\begin{equation*}
(y, t) \mapsto\left(h_{y}-2 h h_{t}\right) \partial_{x}+\left(1-2 y h_{t}\right) \partial_{y}+\left(2 y h_{y}-2 h\right) \partial_{t}, \tag{4.11}
\end{equation*}
$$

which is both horizontal and tangent to the graph of $h$ at $\mathcal{H}^{2}$-a.e. point. We denote by $u=\left(h_{y}-2 h h_{t}, 1-2 y h_{t}\right)$ the projection of the vector field (4.11) onto the $x y$-plane. The relation between $v$ and $u$ is

$$
v=\left(\begin{array}{cc}
0 & 1  \tag{4.12}\\
2 y & -2 h
\end{array}\right) u
$$

In (4.12), we think of $v$ and $u$ as column vectors. The function $\lambda: \Omega \rightarrow \mathbb{R}, \lambda=|u|$, satisfies $0<c_{1} \leq \lambda \leq c_{2} \mathcal{L}^{2}$-a.e. in $\Omega$ by Proposition A. 3 in the appendix.

Let $K$ and $\Omega$ be as in the statement of Theorem 4.2. The existence of a regular Lagrangian flow $\Phi: K \times[-\varrho, \varrho] \rightarrow \Omega$ of $v$ follows from Theorem 4.3.

Let $F \in C^{\infty}\left(\mathbb{R}^{2} ; \mathbb{R}^{2}\right)$ be a mapping which satisfies (3.6) (with $c_{1}$ in place of $\delta$ ). We consider the vector fields $v / \lambda$ and $w=F \circ u$ in $\Omega$. Then $w=u /|u|$ a.e. in $\Omega$. We have $w \in W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$ by Corollary 3.3. Observing that $2 y h_{y}-2 h=$ $2 y\left(h_{y}-2 h h_{t}\right)+2 h\left(2 y h_{t}-1\right)$, we also get that $v / \lambda \in W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$. We claim that

$$
\begin{equation*}
\nabla w \frac{v}{\lambda}=-H w^{\perp} \quad \mathcal{L}^{2} \text {-a.e. in } \Omega \tag{4.13}
\end{equation*}
$$

with $H=3 P(E) / 4|E|$. Indeed, using (3.7), we compute

$$
\begin{aligned}
\nabla w \frac{v}{\lambda} & =(\nabla F \circ u) \nabla u \frac{v}{\lambda}=\frac{1}{|u|^{4}}\left(u^{\perp} \otimes u^{\perp}\right) \nabla u v \\
& =\frac{1}{|u|^{3}}\left(u^{\perp} \cdot(\nabla u v)\right) w^{\perp}
\end{aligned}
$$

On the other hand, by (2.8) we have $\mathcal{M} w=H$ in $\Omega$, where $\mathcal{M}$ is the differential operator appearing in (3.1) with $a=2 h$ and $b=-2 y$. Let $B$ be the $2 \times 2$ matrix

$$
B=\left(\begin{array}{cc}
1 & 0 \\
-2 h & -2 y
\end{array}\right)
$$

Using (3.7) we find

$$
\mathcal{M} w=\operatorname{tr}((\nabla F \circ u) \nabla u B)=\frac{1}{|u|^{3}} \operatorname{tr}\left(\left(u^{\perp} \otimes u^{\perp}\right) \nabla u B\right)=\frac{1}{|u|^{3}} u^{\perp} \cdot\left(u^{\perp} \nabla u B\right),
$$

where, by a short computation based on (4.12), we have $u^{\perp} \cdot\left(u^{\perp} \nabla u B\right)=-u^{\perp}$. $(\nabla u v)$. This ends the proof of (4.13).

Denote by $\gamma_{\zeta}:\left[\sigma_{\zeta}(-\varrho), \sigma_{\zeta}(\varrho)\right] \rightarrow \Omega$ the integral curve of $v / \lambda$ defined in (4.6). Then the curve $s \mapsto w\left(\gamma_{\zeta}(s)\right)$ belongs to $W^{1,1}\left(\left(\sigma_{\zeta}(-\varrho), \sigma_{\zeta}(\varrho)\right) ; \mathbb{R}^{2}\right)$ for $\mathcal{L}^{2}$-a.e. $\zeta \in K$, and its weak derivative is equal to $\left(\nabla w \circ \gamma_{\zeta}\right) \dot{\gamma}_{\zeta}$, by Theorem 4.4. Identity (4.13) holds along $\mathcal{L}^{2}$-a.e. curve $\gamma_{\zeta}$. Hence, by Theorem 4.4 applied to $w$, we have

$$
\begin{equation*}
\frac{d}{d s}\left(w \circ \gamma_{\zeta}\right)=-H\left(w^{\perp} \circ \gamma_{\zeta}\right) \tag{4.14}
\end{equation*}
$$

in the sense of weak derivatives.
The projection of the vector field $\zeta=(y, t) \mapsto(h(\zeta), \zeta)$ onto the $x y$-plane belongs to $W^{1,1}\left(\Omega ; \mathbb{R}^{2}\right)$. Denote by $\kappa_{\zeta}:\left[\sigma_{\zeta}(-\varrho), \sigma_{\zeta}(\varrho)\right] \rightarrow \mathbb{R}^{2}$ the projection of the curve $\left(h\left(\gamma_{\zeta}\right), \gamma_{\zeta}\right)$ onto the $x y$-plane. This projection is a reparameterization of the curve in (4.2). We claim that $\kappa_{\zeta}$ parameterizes an arc of circle with curvature $H$ oriented clockwise. By Theorem 4.4, for $\mathcal{L}^{2}$-a.e. $\zeta \in K$ we have $\kappa_{\zeta} \in W^{1,1}\left(\left(\sigma_{\zeta}(-\varrho), \sigma_{\zeta}(\varrho)\right) ; \mathbb{R}^{2}\right)$, and a short computation shows that the weak derivative of $\kappa_{\zeta}$ is $\dot{\kappa}_{\zeta}=w \circ \gamma_{\zeta}$. From (4.14), we deduce that $\kappa_{\zeta} \in$ $W^{2,1}\left(\left(\sigma_{\zeta}(-\varrho), \sigma_{\zeta}(\varrho)\right) ; \mathbb{R}^{2}\right)$ and $\ddot{\kappa}_{\zeta}=-H \dot{\kappa}_{\zeta}^{\perp}$. This implies that $\kappa_{\zeta}$ is of class $C^{\infty}$ and the claim follows.

## 5. Characterization of convex isoperimetric sets

In this section, we use Theorems 4.1 and 4.2 to characterize convex isoperimetric sets. We also use Theorem A. 1 in the appendix on the structure of the characteristic set of bounded convex sets. Here, we recall some known facts about CarnotCarathéodory geodesics in the Heisenberg group.

An absolutely continuous path $\gamma:[0,1] \rightarrow \mathbb{H}^{1}$ is said to be horizontal if $\dot{\gamma}(s) \in \operatorname{span}_{\mathbb{R}}\{X(\gamma(s)), Y(\gamma(s))\}$ for a.e. $s \in[0,1]$. The curve $\gamma=\left(\gamma_{1}, \gamma_{2}, \gamma_{3}\right)$ is horizontal if and only if

$$
\begin{equation*}
\gamma_{3}(s)=\gamma_{3}(0)+2 \int_{0}^{s}\left(\dot{\gamma}_{1} \gamma_{2}-\dot{\gamma}_{2} \gamma_{1}\right) d \sigma \tag{5.1}
\end{equation*}
$$

The plane curve $\kappa=\left(\gamma_{1}, \gamma_{2}\right)$ is the horizontal projection of $\gamma$. If $\kappa=\left(\gamma_{1}, \gamma_{2}\right)$ is a given plane absolutely continuous curve, then a curve $\gamma=\left(\kappa, \gamma_{3}\right)$ with $\gamma_{3}$ given by (5.1) for some $\gamma_{3}(0)$ is called a horizontal lift of $\kappa$. The standard sub-Riemannian length of the horizontal curve $\gamma$ is

$$
\begin{equation*}
\operatorname{Length}(\gamma)=\int_{0}^{1}|\dot{\kappa}(s)| d s \tag{5.2}
\end{equation*}
$$

where $\kappa$ is the horizontal projection of $\gamma$ and $|\dot{\kappa}|$ is the Euclidean length of $\dot{\kappa} \in \mathbb{R}^{2}$. The distance $d: \mathbb{H}^{1} \times \mathbb{H}^{1} \rightarrow[0,+\infty)$ is defined as the minimum of Length $(\gamma)$
over all horizontal curves connecting given points. Geodesic curves (i.e. lenght minimizing curves) can be computed explicitly. A geodesic from $0 \in \mathbb{H}^{1}$ to $p=$ $(z, t) \in \mathbb{H}^{1}$ is of the following form:

1) If $t=0$ the geodesic is the line segment $\gamma(s)=(s z, 0), s \in[0,1]$.
2) If $t>0$ and $z \neq 0$, the geodesic connecting 0 to $p=(z, t)$ is the horizontal lift starting at 0 of the arc of circle from 0 to $z$ in the $x y$-plane oriented clockwise, such that the plane region bounded by the arc and the segment joining 0 to $z$ has area equal to $t / 4$. This geodesic is unique.
3) If $t>0$ and $z=0$, the geodesic from 0 to $p=(0, t)$ is not unique. Take any full circle oriented clockwise passing through 0 and with area equal to $t / 4$. The horizontal lift of the circle starting from 0 is a geodesic.

The case $t<0$ is similar. All other geodesics are obtained by left translation. If the arc of circle in 2) and 3) has radius $0<R<+\infty$, we say that the geodesic has curvature $H=1 / R$.

The geodesics joining $(0,-\pi / 2)$ to $(0, \pi / 2)$ are of the type described in 3$)$ above (up to a vertical translation). Their union bounds the isoperimetric set conjectured by Pansu. The horizontal lift of the plane circle $\kappa(s)=\frac{1}{2}(1+\cos s,-\sin s)$, $s \in[-\pi, \pi]$, passing through the point $(1,0,0) \in \mathbb{H}^{1}$ at time $s=0$ is the curve $\gamma:[-\pi, \pi] \rightarrow \mathbb{H}^{1}$

$$
\begin{equation*}
\gamma(s)=\frac{1}{2}(1+\cos s,-\sin s, s+\sin s) \tag{5.3}
\end{equation*}
$$

The third coordinate can be computed using formula (5.1). The curve $\gamma$ is a geodesic with curvature $H=2$, starting from $\gamma(-\pi)=(0,0,-\pi / 2)$ and reaching $\gamma(\pi)=$ $(0,0, \pi / 2)$. If $(z, t)=\gamma(s) \in \gamma([-\pi, \pi])$ is a point on the curve, then we have

$$
|z|=\left(\frac{1+\cos s}{2}\right)^{1 / 2} \quad \text { and } \quad t=\frac{1}{2}(s+\sin s)
$$

and we obtain the relation $|t|=\arccos |z|+|z| \sqrt{1-|z|^{2}}$. The region $|t| \leq$ $\arccos |z|+|z| \sqrt{1-|z|^{2}}$ is the set $E_{\text {isop }}$ in (1.4).

Corollary 5.1 (Foliation by geodesics I). Under the assumptions of Theorem 4.1, for all $z \in K$, there is a geodesic $\gamma_{z}:[-\varrho, \varrho] \rightarrow \partial E$ with curvature $H=$ $3 P(E) / 4|E|$, such that $\gamma_{z}(0)=(z, f(z))$. Moreover, the length of $\gamma_{z}$ is bounded from below by a positive constant depending on $K$.
Proof. The vector field $v(z)=2 z-\nabla f^{\perp}(z)$ is the projection onto the $x y$-plane of the vector field

$$
(x, y) \mapsto\left(f_{y}+2 x\right) \partial_{x}+\left(2 y-f_{x}\right) \partial_{y}+\left(2 x f_{x}+2 y f_{y}\right) \partial_{t}
$$

which is both horizontal and tangent to the graph of $f$ at $\mathcal{H}^{2}$-a.e. point. Then the horizontal lift of the plane curve $s \mapsto \Phi(z, s)$ given by Theorem 4.1 is the curve

$$
\gamma_{z}(s)=(\Phi(z, s), f(\Phi(z, s))), \quad s \in[-\varrho, \varrho] .
$$

By Theorem 4.1, this curve is a geodesic with curvature $H$ and with length bounded from below by a positive constant depending on $K$. This curve exists for all $z \in$ $K_{0} \subset K$ with $\mathcal{L}^{2}\left(K \backslash K_{0}\right)=0$. If $z \in K \backslash K_{0}$, there are points $z_{n} \in K_{0}, n \in \mathbb{N}$, such that $z_{n} \rightarrow z$ as $n \rightarrow+\infty$. We can use the theorem of Ascoli-Arzelà to extract a subsequence, such that, after relabeling, the curves $\gamma_{z_{n}}$ converge uniformly to a curve $\gamma:[-\varrho, \varrho] \rightarrow \partial E$. This curve is a geodesic passing through $(z, f(z))$ at time $s=0$, with curvature $H$ and with length bounded from below by the same positive constant as above.

Analogously, we have:
Corollary 5.2 (Foliation by geodesics II). Under the assumptions of Theorem 4.2, for all $\zeta \in K$, there is a geodesic $\gamma_{\zeta}:[-\varrho, \varrho] \rightarrow \partial E$ with curvature $H$, such that $\gamma_{\zeta}(0)=(h(\zeta), \zeta)$. Moreover, the length of $\gamma_{\zeta}$ is bounded from below by a positive constant depending on $K$.

We identify the horizontal plane spanned by the vector fields $X$ and $Y$ at the point $p=0$ with the $x y$-plane

$$
\mathrm{H}_{0}=\left\{(z, 0) \in \mathbb{H}^{1} \mid z \in \mathbb{C}\right\}
$$

In general, if $p=(z, t) \in \mathbb{H}^{1}$ we let

$$
\mathrm{H}_{p}=p \mathrm{H}_{0}=\left\{\left(z+z^{\prime}, t+2 \operatorname{Im} z \bar{z}^{\prime}\right) \in \mathbb{H}^{1} \mid z^{\prime} \in \mathbb{C}\right\} .
$$

The plane $\mathrm{H}_{p}$ is the boundary of the halfspaces

$$
\begin{align*}
& \mathrm{H}_{p}^{-}=\left\{\left(z^{\prime}, t^{\prime}\right) \in \mathbb{H}^{1} \mid t^{\prime} \leq t+2 \operatorname{Im} z \bar{z}^{\prime}\right\}  \tag{5.4}\\
& \mathrm{H}_{p}^{+}=\left\{\left(z^{\prime}, t^{\prime}\right) \in \mathbb{H}^{1} \mid t^{\prime} \geq t+2 \operatorname{Im} z \bar{z}^{\prime}\right\}
\end{align*}
$$

We say that a line $\ell \subset \mathbb{H}^{1}$ is horizontal if $\ell \subset \mathrm{H}_{p}$ for one (equivalently: all) $p \in \ell$. A horizontal segment is a segment of a horizontal line. A plane $\pi \subset \mathbb{H}^{1}$ is a supporting plane for a set $E \subset \mathbb{H}^{1}$ at a point $p \in \partial E$ if $\pi \cap \bar{E} \subset \partial E$.
Definition 5.3 (Characteristic set). The characteristic set of a convex set $E \subset \mathbb{H}^{1}$ is $\Sigma(E)=\left\{p \in \partial E \mid \mathrm{H}_{p}\right.$ is a supporting plane for $E$ at $\left.p\right\}$.

The structure of $\Sigma(E)$ is described in Theorem A. 1 in the appendix.
Proof of Theorem 1.1. The curvature $H=3 P(E) / 4|E|$ of a (convex) isoperimetric set $E$ is homogeneous with degree -1 with respect to the dilations $\delta_{\lambda}(z, t)=$ $\left(\lambda z, \lambda^{2} t\right), \lambda>0$. Thus, we can assume $H=2$ without loss of generality.

Claim 1. For all $p \in \partial E \backslash \Sigma(E)$, there is a curve of maximal length passing through $p$ and contained in $\partial E$ whose projection onto the $x y$-plane is an arc of circle with curvature $H$. One endpoint of the curve belongs to $\Sigma^{-}$, the other one to $\Sigma^{+}$.

Let $f: D \rightarrow \mathbb{R}$ be the function in (1.5), and let $h: F \rightarrow \mathbb{R}$ be the function in (1.10). Using a rotation of $\mathbb{R}^{3}$ fixing the $t$-axis, which is an isometry for the Carnot-Carathéodory metric, and a left translation in the $y$ direction if necessary, we can assume that $p \in \partial E \backslash \Sigma(E)$ is either of the form

1) $p=(z, f(z))$ for some $z \in \operatorname{int}(D) \backslash \Sigma(f)$, or of the form
2) $p=(h(\zeta), \zeta)$ for some $\zeta \in \operatorname{int}(F) \backslash(\{y=0\} \cup \Sigma(h))$.

The existence of a piece of geodesic with curvature $H$ passing through $p$ and contained in $\partial E$ follows either from Corollary 5.1 or from Corollary 5.2. The existence of a curve with maximal length whose projection onto the $x y$-plane is an arc of circle with curvature $H$ follows from a compactness argument. If one endpoint of the maximal curve does not belong to $\Sigma(E)$, then a continuation argument provides a proper extension of the curve. To prove this continuation, notice that two geodesics contained in $\partial E$, with the same curvature and meeting at one point in $\partial E \backslash \Sigma(E)$ coincide in a neighborhood of the point. Finally, the endpoints of the maximal curve cannot both belong to $\Sigma^{-}$or $\Sigma^{+}$.

Claim 2. The sets $\Sigma^{-}$and $\Sigma^{+}$are points.
Assume by contradiction that $\Sigma^{-}$contains more than one point. According to Theorem A.1, we have two cases:

1) $\Sigma^{-}$is a closed horizontal segment of positive length;
2) $\Sigma^{-}=\Sigma_{1}^{-} \cup \Sigma_{2}^{-}$with $\Sigma_{1}^{-}$and $\Sigma_{2}^{-}$disjoint closed horizontal segments lying on the same horizontal line.

In case 1 ), after a left translation and a rotation fixing the $t$-axis, we have $\Sigma^{-}=$ $\left\{(0, y, 0) \in \mathbb{H}^{1}| | y \mid \leq y_{0}\right\}$ for some $y_{0}>0$. Moreover, letting $p_{0}=\left(0, y_{0}, 0\right)$, by the discussion of case 1 ) in the proof of Theorem A.1, we have

$$
\begin{equation*}
E \subset \mathrm{H}_{p_{0}}^{+} \cap \mathrm{H}_{-p_{0}}^{+} \tag{5.5}
\end{equation*}
$$

Claim 1 combined with an approximation-compactness argument shows that there is a geodesic $\gamma:[0, L] \rightarrow \partial E$ parameterized proportionally to arc-length, with curvature $H=2$ and such that $\gamma(0)=0$. We have $\dot{\gamma}(0)=(\alpha, \beta, 0)$ for some $\alpha, \beta \in \mathbb{R}$ with $\alpha^{2}+\beta^{2}=1 / 4$. Assume $\alpha=0$ and $\beta=1 / 2$ (if $\alpha \neq 0$, the proof is easier). Then, by (5.3), we have

$$
\begin{equation*}
\gamma(s)=\frac{1}{2}(1-\cos s, \sin s, s-\sin s) \in \partial E \quad \text { for all } s \in[0, L] \tag{5.6}
\end{equation*}
$$

However, for any $c_{0}>0$, there is $\delta>0$ with $s-\sin s<c_{0}(1-\cos s)$ for all $s \in(0, \delta)$. This contradicts (5.5).

Case 2) cannot occur. In fact, denoting by $\widehat{\Sigma}$ the nonempty open segment between $\Sigma_{1}^{-}$and $\Sigma_{2}^{-}$, we have $\widehat{\Sigma} \subset \partial E$. By Claim 1 , for any $p \in \widehat{\Sigma}$ there is a geodesic contained in $\partial E$ with curvature 2 which passes through $p \in \widehat{\Sigma}$ and has one endpoint in $\Sigma^{-}$. This is impossible.

Claim 3. Up to a left translation, we have $E=E_{\text {isop }}$ as in (1.4).
After a left translation, we can assume that $\Sigma^{-}=\{(0,-\pi / 2)\}$ and $\Sigma^{+}=$ $\{(z, t)\}$ with $t>-\pi / 2$. This forces $z=0$ and $t=\pi / 2$, otherwise there would be at most one geodesic with curvature $H=2$ connecting $\Sigma^{-}$to $\Sigma^{+}$.

We claim that $0 \in \operatorname{int}(D)$. If, by contradiction, $0 \in \partial D$, then $\{(0, t) \in$ $\left.\mathbb{H}^{1}| | t \mid \leq \pi / 2\right\}$ is contained in $\partial E$. But there is no geodesic with curvature $H=2$ which starts from $\Sigma^{-}$and passes through a point $(0, t)$ with $|t|<\pi / 2$. This contradicts Claim 1.

For all $z \in \partial D$, there is a geodesic contained in $\partial E$, with endpoints $(0,-\pi / 2)$, $(0, \pi / 2)$, and which passes through $(z, f(z))$. The projection of this geodesic onto the $x y$-plane is a (full) circle with radius $1 / 2$ passing through the points 0 and $z$. Hence $\partial E$ contains the surface of rotation around the $t$-axis generated by the curve (5.3), and the claim follows.

## A. Appendix

We describe the structure of the characteristic set $\Sigma(E)$ of a compact convex set $E \subset \mathbb{H}^{1}$ with nonempty interior. The notation is fixed before Definition 5.3.

Theorem A. 1 (Characteristic set). If $E \subset \mathbb{H}^{1}$ is a compact convex set with nonempty interior, then $\Sigma(E)$ admits the decomposition $\Sigma(E)=\Sigma^{-} \cup \Sigma^{+}$, where either

1) $\Sigma^{-}$(resp. $\left.\Sigma^{+}\right)$is a compact horizontal segment (possibly one point); or,
2) $\Sigma^{-}=\Sigma_{1}^{-} \cup \Sigma_{2}^{-}\left(r e s p . \Sigma^{+}=\Sigma_{1}^{+} \cup \Sigma_{2}^{+}\right)$, where $\Sigma_{1}^{-}$and $\Sigma_{2}^{-}$(resp. $\Sigma_{1}^{+}$and $\Sigma_{2}^{+}$) are disjoint compact horizontal segments lying on the same horizontal line.

Moreover, if $0 \in \operatorname{int}(E)$, then $\mathrm{H}_{0}$ separates $\Sigma^{-}$from $\Sigma^{+}$.
Proof. The proof is divided into a number of steps.
Claim 1. If $E \subset \mathbb{H}^{1}$ is a bounded strictly convex set, then $\Sigma(E)$ contains at most two points.

Let $p_{0}, p_{1} \in \Sigma(E)$ be such that $p_{0} \neq p_{1}$. After a left translation, we can assume that $p_{0}=0$ and $p_{1}=\left(z_{1}, t_{1}\right)$ with $t_{1}>0 . E$ is contained between the horizontal planes $\mathrm{H}_{p_{0}}$ and $\mathrm{H}_{p_{1}}$. By strict convexity, the segment

$$
\left\{p_{s}=\left(s z_{1}, s t_{1}\right) \in \mathbb{H}^{1} \mid 0<s<1\right\}
$$

is contained in the interior of $E$. Moreover, we have

$$
\partial E=\left\{p_{0}, p_{1}\right\} \cup \bigcup_{0<s<1} \partial E \cap \mathrm{H}_{p_{s}}
$$

By strict convexity, $\partial E \cap \mathrm{H}_{p_{0}}=\left\{p_{0}\right\}$ and $\partial E \cap \mathrm{H}_{p_{1}}=\left\{p_{1}\right\}$. Using the property $p \in$ $\mathrm{H}_{p^{\prime}}$ if and only if $p^{\prime} \in \mathrm{H}_{p}$, we deduce that if $p \in \partial E \cap \mathrm{H}_{p_{s}}$ with $0<s<1$, then $p_{s} \in \mathrm{H}_{p}$ and $\mathrm{H}_{p}$ is not a supporting plane of $E$. It follows that $\Sigma(E)=\left\{p_{0}, p_{1}\right\}$.

Claim 2. If $E \subset \mathbb{H}^{1}$ is a bounded convex set with $C^{\infty}$ boundary, then $\Sigma(E) \neq \emptyset$.
For $p \in \partial E$, let us denote by $V(p) \in \mathrm{T}_{p} \partial E$ the orthogonal projection of the vector field $T$ onto $\mathrm{T}_{p} \partial E$ with respect to the left invariant inner product which makes $X, Y, T$ orthonormal. Since $\partial E$ is diffeomorphic with the sphere $S^{2}$, there exists $p \in \partial E$ such that $V(p)=0$ and therefore $p \in \Sigma(E)$.
Claim 3. If $E \subset \mathbb{H}^{1}$ is a bounded strictly convex set with $C^{\infty}$ boundary, then $\Sigma(E)$ contains exactly two points.

By Claim 2 there exists $p_{0} \in \Sigma(E)$. Without loss of generality we can assume that $p_{0}=0$ and $E \subset \mathrm{H}_{0}^{+}$. Let $V$ be the vector field defined in the proof of Claim 2. It is $V(0)=0$ and using the strict convexity of $E$ it is not difficult to check that index $(V, 0)=1$ (we omit the details). If by contradiction $V$ has no other zero on $\partial E$, then the Poincaré-Hopf index theorem would give index $(V, 0)=\chi(\partial E)=$ $\chi\left(S^{2}\right)=2$, where $\chi$ denotes the Euler-Poincaré characteristic. This is not possible.
Claim 4. If $E \subset \mathbb{H}^{1}$ is a bounded convex set with nonempty interior, then $\Sigma(E)$ contains at least two points.

This follows from Claim 3 by an approximation argument. Details are omitted.
Claim 5. We have $\Sigma(E)=\Sigma^{-} \cup \Sigma^{+}$as in the statement of Theorem A.1.
Without loss of generality, we assume that $0 \in \operatorname{int}(E)$. We define the convex sets

$$
E^{-}=E \cap \mathrm{H}_{0}^{-} \quad \text { and } \quad E^{+}=E \cap \mathrm{H}_{0}^{+} .
$$

If $p \in \partial E^{-} \cap \mathrm{H}_{0}=\partial E^{+} \cap \mathrm{H}_{0}$ with $p \neq 0$, then $0 \in \mathrm{H}_{p} \neq \mathrm{H}_{0}$ and consequently $\mathrm{H}_{p} \cap \operatorname{int}\left(E^{ \pm}\right) \neq \emptyset$. It follows that $\Sigma\left(E^{ \pm}\right) \cap \mathrm{H}_{0}=\{0\}$. By Claim 4 there exist points $p^{-}, p^{+}$such that $p^{ \pm} \in \Sigma\left(E^{ \pm}\right) \cap \operatorname{int}\left(\mathrm{H}_{0}^{ \pm}\right)$. The relative interior of the line segment connecting 0 with $p^{-}$, respectively $p^{+}$, is contained in the interior of $E^{-}$, respectively $E^{+}$. A similar reasoning as in the proof of Claim 1 gives

$$
\begin{equation*}
\Sigma\left(E^{-}\right)=\{0\} \cup\left(\Sigma\left(E^{-}\right) \cap \mathrm{H}_{p^{-}}\right) \quad \text { and } \quad \Sigma\left(E^{+}\right)=\{0\} \cup\left(\Sigma\left(E^{+}\right) \cap \mathrm{H}_{p^{+}}\right) \tag{A.1}
\end{equation*}
$$

We consider $E^{-}$and we let $p^{-}=p_{1}$. Assume there exists $p_{2} \in \Sigma\left(E^{-}\right) \cap \mathrm{H}_{p_{1}}$ with $p_{2} \neq p_{1}$. Since the equality on the left hand side of (A.1) must hold with $p_{2}$ instead of $p_{1}=p^{-}$, it follows that $\Sigma\left(E^{-}\right) \backslash\{0\}$ is contained in the horizontal line $\ell$ passing through $p_{1}$ and $p_{2}$. By a maximality argument, we can in fact assume that $\Sigma\left(E^{-}\right) \backslash\{0\}$ is contained in the compact sub-interval of $\ell$ with endpoints $p_{1}$ and $p_{2}$. Moreover, possibly interchanging $p_{1}$ and $p_{2}$, there are only two possible cases: either

1) $E^{-} \subset \mathrm{H}_{p_{1}}^{+} \cap \mathrm{H}_{p_{2}}^{+}$; or,
2) $E^{-} \subset \mathrm{H}_{p_{1}}^{+} \cap \mathrm{H}_{p_{2}}^{-}$.

In the first case, all points $p \in \ell$ between $p_{1}$ and $p_{2}$ are in $\Sigma\left(E^{-}\right)$, because $\mathrm{H}_{p}$ does not intersect the interior of $\mathrm{H}_{p_{1}}^{+} \cap \mathrm{H}_{p_{2}}^{+}$. We conclude that there exists a compact, horizontal segment $\Sigma^{-}$such that $\Sigma\left(E^{-}\right) \backslash\{0\}=\Sigma^{-}$.

In the second case, there exists a maximal compact sub-interval $\Sigma_{1}^{-}$(resp. $\Sigma_{2}^{-}$) of $\ell$ such that $p_{1} \in \Sigma_{1}^{-}$and $p \in \Sigma\left(E^{-}\right)$for all $p \in \Sigma_{1}^{-}$(resp. $p_{2} \in \Sigma_{2}^{-}$and $p \in \Sigma\left(E^{-}\right)$for all $\left.p \in \Sigma_{2}^{-}\right)$. On the other hand, there exists a point $p \in \ell$ between $p_{1}$ and $p_{2}$ such that $\mathrm{H}_{p} \cap \operatorname{int}\left(E^{-}\right) \neq \emptyset$. Hence $\Sigma_{1}^{-} \cap \Sigma_{2}^{-}=\emptyset$ and we let $\Sigma^{-}=\Sigma_{1}^{-} \cup \Sigma_{2}^{-}$.

Analogously, we get $\Sigma\left(E^{+}\right) \backslash\{0\}=\Sigma^{+}$with $\Sigma^{+}$as in the statement of the theorem. Note that neither $\Sigma^{-}$nor $\Sigma^{+}$intersects $\mathrm{H}_{0}$. Clearly, it is $\Sigma(E)=\Sigma^{-} \cup$ $\Sigma^{+}$.

Let $E$ be a compact convex set in $\mathbb{H}^{1}$ with nonempty interior of the form (1.5). We have the convex function $f: D \rightarrow \mathbb{R}$, where $D$ is a compact and convex subset of $\mathbb{R}^{2}$ with nonempty interior. We define the characteristic set $\Sigma(f)$ of $f$ as the set of points $z \in D$ such that the horizontal plane $\mathrm{H}_{p}$ at the point $p=(z, f(z)) \in$ $\mathbb{R}^{2} \times \mathbb{R}$ is a supporting plane for $E$.

An equation for the horizontal plane $\mathrm{H}_{p}, p=(z, f(z))$, is $t^{\prime}=f(z)-2 z^{\perp} \cdot z^{\prime}$. Then, we have for $z \in \operatorname{int}(D)$

$$
\begin{equation*}
z \in \Sigma(f) \Longleftrightarrow-2 z^{\perp} \in \partial f(z) \tag{A.2}
\end{equation*}
$$

where $\partial f(z)$ denotes the subdifferential of $f$ at $z$.
Proposition A. 2 (Lower bounds I). Let $\Omega \subset \operatorname{int}(D)$ be an open set. Then, for any compact set $K \subset \Omega \backslash \Sigma(f)$ there is a constant $\delta>0$ such that

$$
\begin{equation*}
\left|\nabla f(z)+2 z^{\perp}\right| \geq \delta \tag{A.3}
\end{equation*}
$$

for $\mathcal{L}^{2}$-a.e. $z \in K$.
Proof. We show that (A.3) holds at differentiability points $z \in K$ of $f$. By contradiction, assume that there exists a sequence $z_{k} \in K, k \in \mathbb{N}$, of differentiability points of $f$ such that $\left|\nabla f\left(z_{k}\right)+2 z_{k}^{\perp}\right| \rightarrow 0$ as $k \rightarrow+\infty$. Possibly taking a subsequence, we have $z_{k} \rightarrow z$ for some $z \in K$. It follows that $-2 z^{\perp} \in \partial f(z)$ and this implies $z \in \Sigma(f)$, contradicting the assumption $K \cap \Sigma(f)=\emptyset$.

Now consider a convex set $E$ of the form (1.10). We have the convex function $h: F \rightarrow \mathbb{R}$, with $F \subset \mathbb{R}^{2}$ convex set. We define the characteristic set $\Sigma(h)$ of $h$ as the set of points $\zeta \in F$ such that the horizontal plane $\mathrm{H}_{p}$ at the point $p=(h(\zeta), \zeta) \in \mathbb{R} \times \mathbb{R}^{2}$ is a supporting plane for $E$.

An equation for the horizontal plane $\mathrm{H}_{p}, p=(h(\zeta), \zeta)$ with $y \neq 0$, is

$$
x^{\prime}=h(\zeta)+\frac{h(\zeta)}{y}\left(y^{\prime}-y\right)+\frac{1}{2 y}\left(t^{\prime}-t\right)
$$

Then, we have for $\zeta=(y, t) \in \operatorname{int}(F)$ with $y \neq 0$

$$
\begin{equation*}
\zeta \in \Sigma(h) \Longleftrightarrow\left(\frac{h(\zeta)}{y}, \frac{1}{2 y}\right) \in \partial h(\zeta) \tag{A.4}
\end{equation*}
$$

Proposition A. 3 (Lower bounds II). Let $\Omega \subset \operatorname{int}(F)$ be an open set. Then, for any compact set $K \subset \Omega \backslash \Sigma(h)$ there is a constant $\delta>0$ such that

$$
\begin{equation*}
\left|h_{y}-2 h h_{t}\right|+\left|1-2 y h_{t}\right| \geq \delta \tag{A.5}
\end{equation*}
$$

$\mathcal{L}^{2}$-a.e. in $K$.
Proof. If $r>0$ is sufficiently small, we have a.e. in $K \cap\{|y| \leq r\}$

$$
\left|1-2 y h_{y}\right| \geq 1 / 2 .
$$

In fact, the derivatives of $h$ are locally bounded. We claim that at differentiability points of $h$ in $K \cap\{|y| \geq r\}$ we have

$$
\left|h_{y}-h / y\right|+\left|h_{t}-1 /(2 y)\right| \geq \delta_{1}
$$

for some $\delta_{1}>0$. If this is not the case, arguing as in the proof of Proposition A.2, we find $\zeta \in K \cap\{|y| \geq r\}$ such that $(h(\zeta) / y, 1 / 2 y) \in \partial h(\zeta)$, which is not possible by (A.4). Hence there exists $\delta_{2}>0$ such that the condition $\left|h_{t}-1 /(2 y)\right| \leq \delta_{2}$ implies $\left|h_{y}-h / y\right| \geq \delta_{1} / 2$ and $\left|2 h\left(h_{t}-1 /(2 y)\right)\right| \leq \delta_{1} / 4$, whence

$$
\left|h_{y}-2 h h_{t}\right| \geq\left|h_{y}-h / y\right|-\left|2 h\left(h_{t}-1 /(2 y)\right)\right| \geq \delta_{1} / 4 .
$$

The claim follows with $\delta=\min \left\{1 / 2, \delta_{1} / 4,2 r \delta_{2}\right\}$.

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