# Vertical perimeter versus horizontal perimeter 

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#### Abstract

Given $k \in \mathbb{N}$, the $k$ 'th discrete Heisenberg group, denoted $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$, is the group generated by the elements $a_{1}, b_{1}, \ldots, a_{k}, b_{k}, c$, subject to the commutator relations $\left[a_{1}, b_{1}\right]=\cdots=\left[a_{k}, b_{k}\right]=c$, while all the other pairs of elements from this generating set are required to commute, i.e., for every distinct $i, j \in\{1, \ldots, k\}$, we have $\left[a_{i}, a_{j}\right]=\left[b_{i}, b_{j}\right]=\left[a_{i}, b_{j}\right]=$ $\left[a_{i}, c\right]=\left[b_{i}, c\right]=1$. (In particular, this implies that $c$ is in the center of $\mathbb{H}_{Z}^{2 k+1}$.) Denote $\mathfrak{S}_{k}=\left\{a_{1}, b_{1}, a_{1}^{-1}, b_{1}^{-1}, \ldots, a_{k}, b_{k}, a_{k}^{-1}, b_{k}^{-1}\right\}$. The horizontal boundary of $\Omega \subseteq \mathbb{H}_{Z}^{2 k+1}$, denoted $\partial_{h} \Omega$, is the set of all those pairs $(x, y) \in \Omega \times\left(\mathbb{H}_{Z}^{2 k+1} \backslash \Omega\right)$ such that $x^{-1} y \in \mathfrak{S}_{k}$. The horizontal perimeter of $\Omega$ is the cardinality $\left|\partial_{h} \Omega\right|$ of $\partial_{h} \Omega$; i.e., it is the total number of edges incident to $\Omega$ in the Cayley graph induced by $\mathfrak{S}_{k}$. For $t \in \mathbb{N}$, define $\partial_{v}^{t} \Omega$ to be the set of all those pairs $(x, y) \in \Omega \times\left(\mathbb{H}_{Z}^{2 k+1} \backslash \Omega\right)$ such that $x^{-1} y \in\left\{c^{t}, c^{-t}\right\}$. Thus, $\left|\partial_{v}^{t} \Omega\right|$ is the total number of edges incident to $\Omega$ in the (disconnected) Cayley graph induced by $\left\{c^{t}, c^{-t}\right\} \subseteq \mathbb{H}_{Z}^{2 k+1}$. The vertical perimeter of $\Omega$ is defined by $\left|\partial_{v} \Omega\right|=\sqrt{\sum_{t=1}^{\infty}\left|\partial_{v}^{t} \Omega\right|^{2} / t^{2}}$. It is shown here that if $k \geqslant 2$, then $\left|\partial_{v} \Omega\right| \lesssim \frac{1}{k}\left|\partial_{h} \Omega\right|$. The proof of this "vertical versus horizontal isoperimetric inequality" uses a new structural result that decomposes sets of finite perimeter in the Heisenberg group into pieces that admit an "intrinsic corona decomposition." This allows one to deduce an endpoint $W^{1,1} \rightarrow L_{2}\left(L_{1}\right)$ boundedness of a certain singular integral operator from a corresponding lower-dimensional $W^{1,2} \rightarrow L_{2}\left(L_{2}\right)$ boundedness. Apart from its intrinsic geometric interest, the above (sharp) isoperimetric-type inequality has several (sharp) applications, including that for every $n \in \mathbb{N}$, any embedding into an $L_{1}(\mu)$ space of a ball of radius $n$ in the word metric on $\mathbb{H}_{Z}^{5}$ that is induced by the generating set $\mathfrak{S}_{2}$ incurs bi-Lipschitz distortion that is at least a universal constant multiple of $\sqrt{\log n}$. As an application to approximation algorithms, it follows that for every $n \in \mathbb{N}$, the integrality gap of the Goemans-Linial semidefinite program for the Sparsest Cut Problem on inputs of size $n$ is at least a universal constant multiple of $\sqrt{\log n}$.


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## 1. Introduction

We shall start by formulating the main result that is obtained here in the discrete setting, since this is the setting that is needed for applications to approximation algorithms and metric embeddings. However, we shall soon thereafter proceed to describe the corresponding (and equivalent) continuous setting, which is where it is most natural to carry out the bulk of the ensuing proofs.

For every $k \in \mathbb{N}$, the $k$ 'th discrete Heisenberg group, denoted $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$, is defined via generators and relations as follows: The $2 k+1$ generators are denoted $a_{1}, b_{1}, \ldots, a_{k}, b_{k}, c$. Recalling that the commutator of $g, h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ is denoted $[g, h]=g h g^{-1} h^{-1}$, the defining relations of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ are the requirements that $c=\left[a_{1}, b_{1}\right]=\cdots=\left[a_{k}, b_{k}\right]$ and $\left[a_{i}, a_{j}\right]=\left[b_{i}, b_{j}\right]=\left[a_{i}, b_{j}\right]=\left[a_{i}, c\right]=$ $\left[b_{i}, c\right]=1$ for every distinct $i, j \in\{1, \ldots, k\}$. In particular, since $c$ commutes with all of the generators, it belongs to the center of $\mathbb{H}_{Z}^{2 k+1}$. The set $\mathfrak{S}_{k}=$ $\left\{a_{1}, b_{1}, a_{1}^{-1}, b_{1}^{-1}, \ldots, a_{k}, b_{k}, a_{k}^{-1}, b_{k}^{-1}\right\}$ is a symmetric generating subset of $\mathbb{H}_{Z}^{2 k+1}$, and it therefore induces a left-invariant word metric on $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ which we denote by $d_{W}: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \times \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow[0, \infty)$. A standard concrete realization of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ is the group of all $k+2$ by $k+2$ matrices with integer entries of the following form (and usual matrix multiplication):

$$
\left(\begin{array}{cccccc}
1 & x_{1} & x_{2} & \ldots & x_{k} & z  \tag{1}\\
0 & 1 & 0 & \ldots & 0 & y_{1} \\
\vdots & \ddots & \ddots & \ddots & \vdots & \vdots \\
\vdots & \ddots & \ddots & \ddots & 0 & y_{k-1} \\
\vdots & \ddots & \ddots & \ddots & 1 & y_{k} \\
0 & \cdots & \ldots & \ldots & 0 & 1
\end{array}\right)=I_{k+2}+\sum_{i=1}^{k} x_{i} E_{1, i+1}+\sum_{j=1}^{k} y_{j} E_{j+1, k+2}+z E_{1, k+2}
$$

Here $I_{k+2}$ is the $(k+2) \times(k+2)$ identity matrix and $\left\{E_{i, j}: i, j \in\{1, \ldots, k+2\}\right\}$ is the standard basis of $M_{k+2}(\mathbb{R})$. While it is beneficial to keep in mind the above concrete realization of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$, it will later be convenient to work with a different realization of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$, and moreover the initial part of the ensuing discussion (including all of the main results and applications) can be understood while only referring to the above abstract definition of $\mathbb{H}_{Z}^{2 k+1}$ through generators and relations.

Definition 1 (Discrete boundaries). Given $\Omega \subseteq \mathbb{H}_{\mathbb{Z}}^{2 k+1}$, the horizontal boundary of $\Omega$ is defined by

$$
\begin{equation*}
\partial_{\mathrm{h}} \Omega \stackrel{\text { def }}{=}\left\{(x, y) \in \Omega \times\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} \backslash \Omega\right): x^{-1} y \in \mathfrak{S}_{k}\right\} \tag{2}
\end{equation*}
$$

Given also $t \in \mathbb{N}$, the $t$-vertical boundary of $\Omega$ is defined by

$$
\begin{equation*}
\partial_{v}^{t} \Omega \stackrel{\text { def }}{=}\left\{(x, y) \in \Omega \times\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} \backslash \Omega\right): x^{-1} y \in\left\{c^{t}, c^{-t}\right\}\right\} . \tag{3}
\end{equation*}
$$

The horizontal perimeter of $\Omega$ is defined to be the cardinality $\left|\partial_{\mathrm{h}} \Omega\right|$ of its horizontal boundary. The vertical perimeter of $\Omega$ is defined to be the quantity

$$
\begin{equation*}
\left|\partial_{v} \Omega\right| \stackrel{\text { def }}{=}\left(\sum_{t=1}^{\infty} \frac{\left|\partial_{v}^{t} \Omega\right|^{2}}{t^{2}}\right)^{\frac{1}{2}} \tag{4}
\end{equation*}
$$

Remark 2. It may be instructive to restate some of the concepts that were introduced in Definition 1 through the following graph-theoretic descriptions. The Cayley graph that is induced by a symmetric subset $\Sigma=\Sigma^{-1}$ of a group $G$ will be denoted below by $\mathcal{X}_{\Sigma}(G)$. Recall that this is the graph whose vertex set is $G$ and whose edge set is $E_{\Sigma}(G)=\{\{g, g \sigma\}: g \in G \wedge \sigma \in \Sigma\}$. In particular, $x_{\Sigma}(G)$ is connected if and only if $\Sigma$ generates $G$. With this (standard) terminology, the horizontal boundary $\partial_{\mathrm{h}} \Omega$ is the edge boundary of $\Omega$ in the Cayley graph $X_{\mathfrak{S}_{k}}\left(\mathbb{H}_{Z}^{2 k+1}\right)$, i.e., those pairs $(x, y) \in \mathbb{H}_{Z}^{2 k+1} \times \mathbb{H}_{Z}^{2 k+1}$ such that $\{x, y\} \in E_{\mathfrak{S}_{k}}\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1}\right), x \in \Omega$ and $y \notin \Omega$. In the same vein, the $t$-vertical boundary $\partial_{\mathrm{v}}^{t} \Omega$ is the edge boundary of $\Omega$ in the Cayley graph $X_{\left\{c^{t}, c^{-t}\right\}}\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1}\right)$.

The vertical perimeter as defined in (4) is a more subtle concept, and in particular, it does not have a combinatorial description that is analogous to those of Remark 2. The definition (4) was first published in [66, §4], where the isoperimetric-type conjecture that we resolve here also appeared for the first time. These were formulated by the first named author and were circulating for several years before [66] appeared, intended as a possible route towards the algorithmic application that we indeed succeed to obtain here. The basic idea is that because $c$ is equal to each of the commutators $\left\{\left[a_{i}, b_{i}\right]\right\}_{i=1}^{k}$, if the horizontal boundary of $\Omega$ is small, i.e., it is "difficult" to leave the set $\Omega$ in the "horizontal directions" $\left\{a_{i}^{ \pm 1}, b_{i}^{ \pm 1}\right\}_{i=1}^{k}$, then this should be reflected by the "smallness" of the sequence $\left\{\left|\partial_{\mathrm{v}}^{t} \Omega\right|\right\}_{t=1}^{\infty}$ that measures how difficult it is to leave $\Omega$ in the "vertical directions" $\left\{c^{ \pm t}\right\}_{t=1}^{\infty}$. That this smallness should be measured through the quantity $\left|\partial_{v} \Omega\right|$, i.e., the $\ell_{2}$ norm of the sequence $\left\{\left|\partial_{v}^{t} \Omega\right| / t\right\}_{t=1}^{\infty}$, was arrived at through trial and error, inspired by functional inequalities that were obtained in [8], [66], as explained in [66, §4]; see also Section 1.3 below.

Theorem 3. For every integer $k \geqslant 2$ and every finite subset $\Omega \subseteq \mathbb{H}_{\mathbb{Z}}^{2 k+1}$, we have $\left|\partial_{\mathrm{v}} \Omega\right| \lesssim \frac{1}{k}\left|\partial_{\mathrm{h}} \Omega\right|$.

Asymptotic notation. In Theorem 3 as well as in what follows we use the following standard conventions for asymptotic notation. Given $\alpha, \beta \in(0, \infty)$, the notation $\alpha \lesssim \beta$ and $\beta \gtrsim \alpha$ means that $\alpha \leqslant \gamma \beta$ for some universal constant $\gamma \in(0, \infty)$. The notation $\alpha \asymp \beta$ stands for $(\alpha \lesssim \beta) \wedge(\beta \lesssim \alpha)$. If we need to allow for dependence on parameters, we indicate this by subscripts. For example, in the presence of an auxiliary parameter $\psi$, the notation $\alpha \lesssim \psi \beta$ means that $\alpha \leqslant \gamma(\psi) \beta$, where $\gamma(\psi) \in(0, \infty)$ is allowed to depend only on $\psi$, and similarly for the notation $\alpha \gtrsim_{\psi} \beta$ and $\alpha \asymp_{\psi} \beta$.

The "vertical versus horizontal isoperimetric inequality" of Theorem 3 is sharp up to the implicit universal constant, as exhibited by considering the case when $\Omega$ is a singleton. (See also Remark 23 below.) In fact, this inequality has qualitatively different types of sets $\Omega$ that saturate it. Namely, for every $m \in \mathbb{N}$, there is a finite $\Omega \subseteq \mathbb{H}_{\mathbb{Z}}^{5}$ for which $\mid\left\{j \in \mathbb{N}: \sum_{t=2^{j-1}}^{2^{j}-1}\left|\partial_{v}^{t} \Omega\right|^{2} / t^{2} \asymp\right.$ $\left.\left|\partial_{\mathrm{h}} \Omega\right|^{2} / m\right\} \mid \asymp m$; this follows from an examination of the sub-level sets of the embeddings that are discussed in the paragraph following Theorem 6 below. We shall later see that by a simple argument the case $k=2$ of Theorem 3 implies its statement for general $k \geqslant 2$. The case $k=2$ is therefore the heart of the matter, and moreover all of our applications of Theorem 3 use only this special case. For these reasons, the discussion in most of the introduction will focus on $k=2$, but our ensuing proofs will be carried out for general $k$ because they derive structural results that are of interest in any dimension and do not reduce directly to a fixed dimension.

Remark 4. Note the restriction $k \geqslant 2$ in Theorem 3. The case $k=1$ remains an intriguing mystery and a subject of our ongoing research that will be published elsewhere. This ongoing work shows that Theorem 3 fails for $k=1$, but that there exists $q \in(2, \infty)$ such that for every $\Omega \subseteq \mathbb{H}_{Z}^{3}$, we have

$$
\begin{equation*}
\left(\sum_{t=1}^{\infty} \frac{\left|\partial_{\mathrm{v}}^{t} \Omega\right|^{q}}{t^{1+\frac{q}{2}}}\right)^{\frac{1}{q}} \lesssim\left|\partial_{\mathrm{h}} \Omega\right| . \tag{5}
\end{equation*}
$$

A simple argument shows that $\sup _{s \in \mathbb{N}}\left|\partial_{\mathrm{v}}^{s} \Omega\right| / \sqrt{s} \leqslant \gamma\left|\partial_{\mathrm{h}} \Omega\right|$ for some universal constant $\gamma>0$. Hence, for every $t \in \mathbb{N}$, we have

$$
\frac{\left|\partial_{\mathrm{v}}^{t} \Omega\right|^{q}}{t^{1+\frac{q}{2}}} \leqslant\left(\frac{\left|\partial_{\mathrm{v}}^{t} \Omega\right|}{t}\right)^{2} \sup _{s \in \mathbb{N}}\left(\frac{\left|\partial_{\mathrm{v}}^{s} \Omega\right|}{\sqrt{s}}\right)^{q-2} \leqslant\left(\frac{\left|\partial_{\mathrm{v}}^{t} \Omega\right|}{t}\right)^{2}\left(\gamma\left|\partial_{\mathrm{h}} \Omega\right|\right)^{q-2} .
$$

This implies that the left-hand side of (5) is bounded from above by a universal constant multiple of $\left|\partial_{\mathrm{v}} \Omega\right|^{2 / q}\left|\partial_{\mathrm{h}} \Omega\right|^{1-2 / q}$. Therefore (5) is formally weaker than the estimate $\left|\partial_{v} \Omega\right| \lesssim\left|\partial_{\mathrm{h}} \Omega\right|$ of Theorem 3. It would be interesting to determine the infimum over those $q$ for which (5) holds true for every $\Omega \subseteq \mathbb{H}_{\mathbb{Z}}^{3}$, with ongoing work indicating that it is at least 4 . It should be stressed, however, that all the applications of Theorem 3 that are obtained here (i.e., the algorithmic application in Section 1.1 and the geometric applications in Section 1.2) use the case $k=2$ of Theorem 3, and understanding the case $k=1$ would not yield any further improvements. So, while the case $k=1$ is geometrically interesting in its own right, it is not needed for the applications that we currently have in mind.
1.1. Sparsest Cut. The Sparsest Cut Problem is a central open question in approximation algorithms that has attracted major research efforts over the past decades. A remarkable aspect of this natural and versatile optimization
problem is that it admits a simple to describe algorithm that was proposed in the mid-1990s by Goemans and Linial, yet while this specific algorithm received close scrutiny by many researchers and to date it is still the best-known approximation algorithm for the Sparsest Cut Problem, despite major efforts it remained unknown for many years how well this algorithm can perform in general. Here we settle the longstanding open question of determining (up to lower order factors) the approximation ratio of the Goemans-Linial algorithm. Modulo previously published reductions, this is a quick consequence of Theorem 3. In fact, the statement of Theorem 3 was initially conjectured as a possible route towards this algorithmic application.

Fix $n \in \mathbb{N}$. The input of the Sparsest Cut Problem consists of two $n$-by- $n$ symmetric matrices with nonnegative entries $C=\left(C_{i j}\right), D=\left(D_{i j}\right) \in$ $M_{n}([0, \infty))$, which are often called capacities and demands, respectively. The goal is to design a polynomial-time algorithm to evaluate the quantity

$$
\begin{equation*}
\mathrm{OPT}(C, D) \stackrel{\text { def }}{=} \min _{\varnothing \subsetneq A \subsetneq\{1, \ldots, n\}} \frac{\sum_{(i, j) \in A \times(\{1, \ldots, n\} \backslash A)} C_{i j}}{\sum_{(i, j) \in A \times(\{1, \ldots, n\} \backslash A)} D_{i j}} . \tag{6}
\end{equation*}
$$

In view of the extensive literature on the Sparsest Cut Problem, it would be needlessly repetitive to recount here the rich and multifaceted impact of this optimization problem on computer science and mathematics; see instead the articles [61], [72], the surveys [106], [75], [19], [90], Chapter 10 of the monograph [35], Chapter 15 of the monograph [80], Chapter 1 of the monograph [99], and the references therein. It suffices to say that by tuning the choice of matrices $C, D$ to the problem at hand, the minimization in (6) finds a partition of the "universe" $\{1, \ldots, n\}$ into two parts, namely, the sets $A$ and $\{1, \ldots, n\} \backslash A$, whose appropriately weighted interface is as small as possible, thus allowing for inductive solutions of various algorithmic tasks, a procedure known as divide and conquer. (Not all of the uses of the Sparsest Cut Problem fit into this framework. A recent algorithmic application of a different nature can be found in [78].)

It is NP-hard to compute $\operatorname{OPT}(C, D)$ [105]. By [29] there exists $\varepsilon_{0}>0$ such that it is even NP-hard to compute $\operatorname{OPT}(C, D)$ within a multiplicative factor of less than $1+\varepsilon_{0}$. If one assumes Khot's Unique Games Conjecture [56], [57], [109], then by [21], [59] there does not exist a polynomial-time algorithm that can compute $\mathrm{OPT}(C, D)$ within any universal constant factor.

Due to the above hardness results, a much more realistic goal would be to design a polynomial-time algorithm that takes as input the capacity and demand matrices $C, D \in M_{n}([0, \infty))$ and outputs a number $\operatorname{ALG}(C, D)$ that is guaranteed to satisfy $\operatorname{ALG}(C, D) \leqslant \mathrm{OPT}(C, D) \leqslant \rho(n) \operatorname{ALG}(C, D)$, with (hopefully) the quantity $\rho(n)$ growing to $\infty$ slowly as $n \rightarrow \infty$. Determining the best
possible asymptotic behaviour of $\rho(n)$ (assuming $\mathrm{P} \neq \mathrm{NP}$ ) is an open problem of major importance.

In [77], [7] an algorithm was designed, based on linear programming (through the connection to multicommodity flows) and Bourgain's embedding theorem [15], which yields $\rho(n)=O(\log n)$. An algorithm based on semidefinite programming (to be described precisely below) was proposed by Goemans and Linial in the mid-1990s. To the best of our knowledge this idea first appeared in the literature in [43, p. 158], where it was speculated that it might even yield a constant factor approximation for the Sparsest Cut Problem; see also [75], [76].

The Goemans-Linial algorithm is simple to describe. It takes as input the symmetric matrices $C=\left(C_{i j}\right), D=\left(D_{i j}\right) \in M_{n}([0, \infty))$ and proceeds to compute the following quantity:

$$
\begin{equation*}
\operatorname{SDP}(C, D) \stackrel{\text { def }}{=} \inf _{\left(v_{1}, \ldots, v_{n}\right) \in \operatorname{NEG}_{n}} \frac{\sum_{i=1}^{n} \sum_{j=1}^{n} C_{i j}\left\|v_{i}-v_{j}\right\|_{\ell_{2}^{n}}^{2}}{\sum_{i=1}^{n} \sum_{j=1}^{n} D_{i j}\left\|v_{i}-v_{j}\right\|_{2}^{2}}, \tag{7}
\end{equation*}
$$

where
$\operatorname{NEG}_{n} \stackrel{\text { def }}{=}\left\{\left(v_{1}, \ldots v_{n}\right) \in\left(\mathbb{R}^{n}\right)^{n}:\left\|v_{i}-v_{j}\right\|_{\ell_{2}^{n}}^{2}\right.$

$$
\left.\leqslant\left\|v_{i}-v_{k}\right\|_{\ell_{2}^{n}}^{2}+\left\|v_{k}-v_{j}\right\|_{\ell_{2}^{n}}^{2} \text { for all } i, j, k \in\{1, \ldots, n\}\right\} .
$$

Therefore $\mathrm{NEG}_{n}$ is the set of $n$-tuples $\left(v_{1}, \ldots v_{n}\right)$ of vectors in $\mathbb{R}^{n}$ such that $\left(\left\{v_{1}, \ldots, v_{n}\right\}, v_{n}\right)$ is a semi-metric space, where $v_{n}: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow[0, \infty)$ is defined by
$\forall x=\left(x_{1}, \ldots, x_{n}\right), y=\left(y_{1}, \ldots, y_{n}\right) \in \mathbb{R}^{n}, \quad v_{n}(x, y) \stackrel{\text { def }}{=} \sum_{j=1}^{n}\left(x_{j}-y_{j}\right)^{2}=\|x-y\|_{\ell_{2}^{n}}^{2}$.
A semi-metric space $(m, d)$ is said (see, e.g., [35]) to be of negative type if $(m, \sqrt{d})$ embeds isometrically into a Hilbert space. So, $\mathrm{NEG}_{n}$ can be described as the set of all (ordered) negative type semi-metrics of size $n$. The evaluation of the quantity $\operatorname{SDP}(C, D)$ in (7) can be cast as a semidefinite program (SDP), so it can be achieved (up to $o(1)$ precision) in polynomial time [46]. One has $\mathrm{SDP}(C, D) \leqslant \mathrm{OPT}(C, D)$ for all symmetric matrices $C, D \in M_{n}([0, \infty))$. See, e.g., $[80, \S 15.9]$ or $[90, \S 4.3]$ for an explanation of the above assertions about $\operatorname{SDP}(C, D)$, as well as additional background and motivation. The pertinent question is therefore to evaluate the asymptotic behavior as $n \rightarrow \infty$ of the following quantity, which is commonly known as the integrality gap of the Goemans-Linial semidefinite programming relaxation (7) for the Sparsest Cut Problem:

$$
\rho_{\mathrm{GL}}(n) \stackrel{\text { def }}{=} \sup _{\substack{C, D \in M_{n}([0, \infty)) \\ C, D \text { symmetric }}} \frac{\operatorname{OPT}(C, D)}{\operatorname{SDP}(C, D)} .
$$

Since $\rho_{\mathrm{GL}}(n)$ can be computed in polynomial time, we have $\rho(n) \leqslant \rho_{\mathrm{GL}}(n)$; the algorithmic output $\operatorname{ALG}(C, D)$ in this case is simply $\operatorname{SDP}(C, D)$. The above
mentioned Goemans-Linial conjecture was that $\rho_{\mathrm{GL}}(n)=O(1)$. This hope was dashed in the remarkable work [59], where the lower bound $\rho_{\mathrm{GL}}(n) \gtrsim$ $\sqrt[6]{\log \log n}$ was proven. An improved analysis of the ideas of [59] was conducted in [64], yielding the estimate $\rho_{\mathrm{GL}}(n) \gtrsim \log \log n$. An entirely different approach based on the geometry of the Heisenberg group was introduced in [69]. In combination with the important works [23], [24] it gives a different proof that $\lim _{n \rightarrow \infty} \rho_{\mathrm{GL}}(n)=\infty$. In [25], [26] the previously best-known lower bound $\rho_{\mathrm{GL}}(n) \gtrsim(\log n)^{\delta}$ was obtained for an effective (but small) positive universal constant $\delta$.

Despite these lower bounds, the Goemans-Linial algorithm yields an approximation ratio of $o(\log n)$, so it is asymptotically more accurate than the linear program of [77], [7]. Specifically, in [20] it was shown that $\rho_{\mathrm{GL}}(n) \lesssim$ $(\log n)^{\frac{3}{4}}$, and this was improved in [4] to $\rho_{\mathrm{GL}}(n) \lesssim(\log n)^{\frac{1}{2}+o(1)}$. See Section 10 below for additional background on the results quoted above. No other polynomial-time algorithm for the Sparsest Cut problem is known (or conjectured) to have an approximation ratio that is asymptotically better than that of the Goemans-Linial algorithm. However, despite major scrutiny by researchers in approximation algorithms, the asymptotic behavior of $\rho_{\mathrm{GL}}(n)$ as $n \rightarrow \infty$ remained unknown. Theorem 5 below resolves this question up to lower-order factors.

Theorem 5. For every integer $n \geqslant 2$, we have $\sqrt{\log n} \lesssim \rho_{\mathrm{GL}}(n) \lesssim$ $(\log n)^{\frac{1}{2}+o(1)}$.

Since the upper bound on $\rho_{\mathrm{GL}}(n)$ in Theorem 5 is due to [4], our contribution is the corresponding lower bound, thus determining the asymptotic behavior of $\rho_{\mathrm{GL}}(n)$ up to lower order factors.
1.2. Embeddings. For $p \in[1, \infty]$, the $L_{p}$ (bi-Lipschitz) distortion of a separable metric space $(m, d)$, denoted $c_{p}(m, d) \in[1, \infty]$, is the infimum over those $D \in[1, \infty]$ for which there exists an embedding $f: m \rightarrow L_{p}(\mathbb{R})$ such that $d(x, y) \leqslant\|f(x)-f(y)\|_{L_{p}(\mathbb{R})} \leqslant D d(x, y)$ for every $x, y \in m$.

Recall that $d_{W}: \mathbb{H}_{\mathbb{Z}}^{5} \times \mathbb{H}_{\mathbb{Z}}^{5} \rightarrow \mathbb{N} \cup\{0\}$ denotes the left-invariant word metric that is induced by the generating set $\mathfrak{S}_{2}=\left\{a_{1}, a_{1}^{-1}, b_{1}, b_{1}^{-1}, a_{2}, a_{2}^{-1}, b_{2}, b_{2}^{-1}\right\}$. For every $r \in[0, \infty)$, denote the corresponding (closed) ball of radius $r$ centered at the identity element by $\mathcal{B}_{r}=\left\{h \in \mathbb{H}_{\mathbb{Z}}^{5}: d_{W}(h, 1) \leqslant r\right\}$. It is well known (see, e.g., [11]) that $\left|\mathcal{B}_{r}\right| \asymp r^{6}$ and $d_{W}\left(1, c^{r}\right) \asymp \sqrt{r}$ for every $r \in \mathbb{N}$. By [69, Th. 2.2], the metric $d_{W}$ is bi-Lipschitz equivalent to a metric on $\mathbb{H}_{\mathbb{Z}}^{5}$ that is of negative type. We remark that [69] makes this assertion for a different metric on a larger continuous group that contains $\mathbb{H}_{\mathbb{Z}}^{5}$ as a discrete co-compact subgroup, but by a simple general result (see, e.g., [18, Th. 8.3.19]) the word
metric $d_{W}$ is bi-Lipschitz equivalent to the metric considered in [69]. A wellknown duality argument (see, e.g., [90, Lemma 4.5] or [25, §1]), which to the best of our knowledge was first derived by Rabinovich in the 1990s, establishes that for every $n \in \mathbb{N}$, the integrality gap $\rho_{\mathrm{GL}}(n)$ is equal to the supremum of the $L_{1}$ distortion $c_{1}(X, d)$ over all $n$-point metric spaces $(X, d)$ of negative type. Therefore, in order to establish Theorem 5 it suffices to prove the following theorem.

Theorem 6. For every $r \geqslant 2$, we have $c_{1}\left(\mathcal{B}_{r}, d_{W}\right) \asymp \sqrt{\log r} \asymp \sqrt{\log \left|\mathcal{B}_{r}\right|}$.
The new content of Theorem 6 is the lower bound $c_{1}\left(\mathcal{B}_{r}, d_{W}\right) \gtrsim \sqrt{\log r}$. The matching upper bound $c_{1}\left(\mathcal{B}_{r}, d_{W}\right) \lesssim \sqrt{\log r}$ has several proofs in the literature; see, e.g., the discussion immediately following Corollary 1.3 in [66] or Remark 7 below. The previously best-known estimate [26] was that there exists a universal constant $\delta>0$ such that $c_{1}\left(\mathcal{B}_{r}, d_{W}\right) \geqslant(\log r)^{\delta}$.

Remark 7. Theorem 6 also yields a sharp result for the general problem of finding the asymptotically largest-possible $L_{1}$ distortion of a finite doubling metric space with $n$ points. A metric space $(m, d)$ is said to be $K$-doubling for some $K \in \mathbb{N}$ if every ball in $m$ (centered anywhere and of any radius) can be covered by $K$ balls of half its radius. By [63],

$$
\begin{equation*}
c_{1}(m, d) \lesssim \sqrt{(\log K) \log |m|} . \tag{8}
\end{equation*}
$$

As noted in [47], the dependence on $|m|$ in (8), but with a worse dependence on $K$, follows by combining [6] and [104]. (This dependence on $K$ was improved in [47].) The metric space $\left(\mathbb{H}_{\mathbb{Z}}^{5}, d_{W}\right)$ is $O(1)$-doubling because $\left|\mathcal{B}_{r}\right| \asymp r^{6}$ for every $r \geqslant 1$. Theorem 6 shows that (8) is sharp when $K=O(1)$, thus improving over the previously best-known construction [71] of arbitrarily large $O(1)$-doubling finite metric spaces $\left\{\left(m_{i}, d_{i}\right)\right\}_{i=1}^{\infty}$ for which $c_{1}\left(m_{i}, d_{i}\right) \gtrsim$ $\sqrt{\left(\log \left|m_{i}\right|\right) / \log \log \left|m_{i}\right|}$. Probably (8) is sharp for all $K \leqslant|m|$; conceivably this could be proven by incorporating Theorem 6 into the argument of [50], but we shall not pursue this here. Theorem 6 establishes for the first time the existence of a metric space that simultaneously has several useful geometric properties yet poor (indeed, worst-possible) embeddability properties into $L_{1}$. By virtue of being $O(1)$-doubling, $\left(\mathbb{H}_{\mathbb{Z}}^{5}, d_{W}\right)$ also has Markov type 2 due to [36] (which improves over [93], where the conclusion that it has Markov type $p$ for every $p<2$ was obtained). For more on the bi-Lipschitz invariant Markov type and its applications, see [9], [91]. The property of having Markov type 2 is shared by the construction of $[71]$, which is also $O(1)$-doubling, but $\left(\mathbb{H}_{Z}^{5}, d_{W}\right)$ has additional features that the example of [71] fails to have. For one, it is a group; for another, by [73], [74] we know that $\left(\mathbb{H}_{\mathbb{Z}}^{5}, d_{W}\right)$ has Markov convexity 4 (and no less). (See [70], [86] for background on the bi-Lipschitz invariant

Markov convexity and its consequences.) By [86, §3] the example of [71] does not have Markov convexity $p$ for any finite $p$. No examples of arbitrarily large finite metric spaces $\left\{\left(m_{i}, d_{i}\right)\right\}_{i=1}^{\infty}$ with bounded Markov convexity (and with the associated Markov convexity constants uniformly bounded) such that $c_{1}\left(m_{i}, d_{i}\right) \gtrsim \sqrt{\log \left|m_{i}\right|}$ were previously known to exist. Analogous statements are known to be impossible for Banach spaces [87], so it is natural in the context of the Ribe program (see the surveys [91], [10] for more on this research program) to ask whether there is a potential metric version of [87]; the above discussion shows that there is not.

Remark 8. It was proved in [65] that for every $p>2$, there is a doubling subset $D_{p}$ of $L_{p}(\mathbb{R})$ that does not admit a bi-Lipschitz embedding into $L_{q}(\mathbb{R})$ for $q \in(1, p)$, and furthermore there is $p_{0}>2$ such that $D_{p}$ does not even admit a bi-Lipschitz embedding into $L_{1}(\mathbb{R})$ for $p>p_{0}$. By substituting Theorem 6 into the proof of this statement in [65], we see that it actually holds true with $p_{0}=2$.

The following precise theorem about $L_{1}$ embeddings that need not be biLipschitz implies Theorem 6 by considering the special case of the modulus $\omega(t)=t / D$ for $D \geqslant 1$ and $t \in[0, \infty)$.

Theorem 9. Fix $r \geqslant 2$ and a nondecreasing function $\omega:[0, \infty) \rightarrow[0, \infty)$ such that $\omega(s) \leqslant s$ for all $s \geqslant 0$. Then there exists a mapping $\phi: \mathcal{B}_{r} \rightarrow L_{1}(\mathbb{R})$ that satisfies

$$
\begin{equation*}
\forall x, y \in \mathcal{B}_{r}, \quad \omega\left(d_{W}(x, y)\right) \lesssim\|\phi(x)-\phi(y)\|_{L_{1}(\mathbb{R})} \leqslant d_{W}(x, y) \tag{9}
\end{equation*}
$$

## if and only if

$$
\begin{equation*}
\int_{1}^{2 r} \frac{\omega(s)^{2}}{s^{3}} \mathrm{~d} s \lesssim 1 \tag{10}
\end{equation*}
$$

The fact that the integrability requirement (10) implies the existence of the desired embedding $\phi$ is due to [108, Cor. 5]. The new content of Theorem 9 is that the existence of the embedding $\phi$ implies (10). By letting $r \rightarrow \infty$ in Theorem 9, we see that there is $\phi: \mathbb{H}_{\mathbb{Z}}^{5} \rightarrow L_{1}(\mathbb{R})$ that satisfies

$$
\begin{equation*}
\forall x, y \in \mathbb{Z}^{5}, \quad \omega\left(d_{W}(x, y)\right) \lesssim\|\phi(x)-\phi(y)\|_{L_{1}(\mathbb{R})} \leqslant d_{W}(x, y) \tag{11}
\end{equation*}
$$

## if and only if

$$
\begin{equation*}
\int_{1}^{\infty} \frac{\omega(s)^{2}}{s^{3}} \mathrm{~d} s \lesssim 1 \tag{12}
\end{equation*}
$$

In [26] it was shown that if $\phi: \mathbb{H}_{\mathbb{Z}}^{5} \rightarrow L_{1}(\mathbb{R})$ satisfies (11), then there must exist arbitrarily large $t \geqslant 2$ for which $\omega(t) \lesssim t /(\log t)^{\delta}$, where $\delta>0$ is a universal constant. This follows from (12) with $\delta=\frac{1}{2}$, which is the largest possible constant for which this conclusion holds true. This positively answers
a question that was asked in [26, Rem. 1.7]. In fact, it provides an even better conclusion, because (12) implies that, say, there must exist arbitrarily large $t \geqslant 4$ for which

$$
\omega(t) \lesssim \frac{t}{\sqrt{(\log t) \log \log t}} .
$$

(The precise criterion is determined by the integrability condition (12).) Finally, by considering $\omega(t)=t^{1-\varepsilon} / D$ for $\varepsilon \in(0,1)$ and $D \geqslant 1$, we obtain the following noteworthy corollary.

Corollary 10 ( $L_{1}$ distortion of snowflakes). For every $\varepsilon \in(0,1)$, we have $c_{1}\left(\mathbb{H}_{\mathbb{Z}}^{5}, d_{W}^{1-\varepsilon}\right) \asymp \frac{1}{\sqrt{\varepsilon}}$.

The fact that we have $c_{1}\left(X, d^{1-\varepsilon}\right) \lesssim 1 / \sqrt{\varepsilon}$ for every $O(1)$-doubling metric space ( $X, d$ ) follows from an argument of [68]. (See also [95, Th. 5.2]). Corollary 10 shows that this is sharp. More generally, it follows from Theorem 9 that for every $r \geqslant 2$ and $\varepsilon \in(0,1)$, we have

$$
c_{1}\left(\mathcal{B}_{r}, d_{W}^{1-\varepsilon}\right) \asymp \min \left\{\frac{1}{\sqrt{\varepsilon}}, \sqrt{\log r}\right\} .
$$

1.3. An endpoint estimate. An equivalent formulation of the case $k=2$ of Theorem 3 is that every finitely supported function $\phi: \mathbb{H}_{\mathbb{Z}}^{5} \rightarrow \mathbb{R}$ satisfies the following Poincaré-type inequality:

$$
\begin{align*}
& \left(\sum_{t=1}^{\infty} \frac{1}{t^{2}}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{5}}\left|\phi\left(h c^{t}\right)-\phi(h)\right|\right)^{2}\right)^{\frac{1}{2}} \lesssim \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{5}} \sum_{\sigma \in \mathfrak{S}_{2}}|\phi(h \sigma)-\phi(h)| \\
& =2 \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{5}}\left(\left|\phi\left(h a_{1}\right)-\phi(h)\right|+\left|\phi\left(h b_{1}\right)-\phi(h)\right|\right.  \tag{13}\\
& \left.\quad+\left|\phi\left(h a_{2}\right)-\phi(h)\right|+\left|\phi\left(h b_{2}\right)-\phi(h)\right|\right) .
\end{align*}
$$

See Lemma 28 below for a simple proof of this equivalence. One direction is immediate, because Theorem 3 is nothing more than the special case $\phi=\mathbf{1}_{\Omega}$ of (13). Furthermore, by a straightforward convexity argument that appears in Lemma 29 below, the estimate (13) has a vector-valued version that asserts that for every finitely supported function $\phi: \mathbb{H}_{\mathbb{Z}}^{5} \rightarrow L_{1}(\mathbb{R})$, we have

$$
\begin{equation*}
\left(\sum_{t=1}^{\infty} \frac{1}{t^{2}}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{5}}\left\|\phi\left(h c^{t}\right)-\phi(h)\right\|_{L_{1}(\mathbb{R})}\right)^{2}\right)^{\frac{1}{2}} \lesssim \sum_{h \in \mathbb{H}_{\frac{\mathbb{Z}}{5}}} \sum_{\sigma \in \mathfrak{S}_{2}}\|\phi(h \sigma)-\phi(h)\|_{L_{1}(\mathbb{R})} . \tag{14}
\end{equation*}
$$

Next, as explained in Lemma 43 below (mimicking an argument that appears in Section 3.2 of [66]), the vector-valued inequality (14) formally implies its local counterpart, which asserts that there exists a universal constant $\alpha \geqslant 1$
such that for every $n \in \mathbb{N}$ and every $\phi: \mathbb{H}_{\mathbb{Z}}^{5} \rightarrow L_{1}(\mathbb{R})$, we have

$$
\begin{equation*}
\left(\sum_{t=1}^{n^{2}} \frac{1}{t^{2}}\left(\sum_{h \in \mathcal{B}_{n}}\left\|\phi\left(h c^{t}\right)-\phi(h)\right\|_{L_{1}(\mathbb{R})}\right)^{2}\right)^{\frac{1}{2}} \lesssim \sum_{h \in \mathcal{B}_{\alpha n}} \sum_{\sigma \in \mathfrak{S}_{2}}\|\phi(h \sigma)-\phi(h)\|_{L_{1}(\mathbb{R})} . \tag{15}
\end{equation*}
$$

To deduce Theorem 9 from (15), suppose that $r \geqslant 2$, that $\omega:[0, \infty) \rightarrow$ $[0, \infty)$ is a nondecreasing function satisfying $\omega(s) \leqslant s$ for all $s \geqslant 0$, and that the mapping $\phi: \mathcal{B}_{r} \rightarrow L_{1}(\mathbb{R})$ satisfies (9). For notational convenience, fix universal constants $\beta \in(0,1)$ and $\gamma \in(1, \infty)$ such that for every $t \in \mathbb{N}$, we have $\beta \sqrt{t} \leqslant d_{W}\left(c^{t}, 1\right) \leqslant \gamma \sqrt{t}$. Since $\omega(s) \leqslant s$ for all $s \geqslant 0$, the lefthand side of $(10)$ is at most $\log 2 r$. Hence, it suffices to prove Theorem 9 for $r \geqslant 1+\max \{\alpha, \gamma\}$, where $\alpha$ is the universal constant in (15). Denote $n=$ $\lfloor\min \{r /(1+\gamma),(r-1) / \alpha\}\rfloor$. If $t \in\left\{1, \ldots, n^{2}\right\}$ and $h \in \mathcal{B}_{n}$, then $d_{W}\left(h c^{t}, 1\right) \leqslant$ $n+\gamma \sqrt{t} \leqslant(1+\gamma) n \leqslant r$, and therefore we may apply (9) with $x=h c^{t}$ and $y=h$ to deduce that $\left\|\phi\left(h c^{t}\right)-\phi(h)\right\|_{L_{1}(\mathbb{R})} \gtrsim \boldsymbol{\omega}\left(d_{W}\left(c^{t}, 1\right)\right) \geqslant \boldsymbol{\omega}(\beta \sqrt{t})$. Consequently,

$$
\begin{align*}
& \sum_{t=1}^{n^{2}} \frac{1}{t^{2}}\left(\sum_{h \in \mathcal{B}_{n}}\left\|\phi\left(h c^{t}\right)-\phi(h)\right\|_{L_{1}(\mathbb{R})}\right)^{2}  \tag{16}\\
& \quad \gtrsim \sum_{t=1}^{n^{2}} \frac{\left|\mathcal{B}_{n}\right|^{2} \omega(\beta \sqrt{t})^{2}}{t^{2}} \gtrsim n^{12} \sum_{t=1}^{n^{2}} \int_{t}^{t+1} \frac{\omega(\beta \sqrt{u / 2})^{2}}{u^{2}} \mathrm{~d} u \\
& \quad=\beta^{2} n^{12} \int_{\frac{\beta}{\sqrt{2}}}^{\frac{\beta \sqrt{n^{2}+1}}{\sqrt{2}}} \frac{\omega(s)^{2}}{s^{3}} \mathrm{~d} s \geqslant \frac{\beta^{2}(r / 2)^{12}}{\max \left\{(1+\gamma)^{12}, \alpha^{12}\right\}} \int_{1}^{\frac{\beta r}{2 \max \{1+\gamma, \alpha\}}} \frac{\omega(s)^{2}}{s^{3}} \mathrm{~d} s
\end{align*}
$$

where the second inequality in (16) uses the fact that $\omega$ is nondecreasing, the penultimate step of (16) uses the change of variable $s=\beta \sqrt{u / 2}$, and for the final step of (16), recall that $\beta<1$ and the definition of $n$. Recalling that $\beta, \gamma>0$ are universal constants, it follows that

$$
\begin{align*}
\int_{1}^{2 r} \frac{\omega(s)^{2}}{s^{3}} \mathrm{~d} s & =\int_{\frac{\beta r}{2 \max \{1+\gamma, \alpha\}}}^{2 r} \frac{\omega(s)^{2}}{s^{3}} \mathrm{~d} s+\int_{1}^{\frac{\beta r}{2 \max \{1+\gamma, \alpha\}}} \frac{\omega(s)^{2}}{s^{3}} \mathrm{~d} s \\
& \lesssim 1+\frac{1}{r^{12}} \sum_{t=1}^{n^{2}} \frac{1}{t^{2}}\left(\sum_{h \in \mathcal{B}_{n}}\left\|\phi\left(h c^{t}\right)-\phi(h)\right\|_{L_{1}(\mathbb{R})}\right)^{2}, \tag{17}
\end{align*}
$$

where the final step of (17) uses (16) and that $\omega(s) \leqslant s$ for all $s \geqslant 0$. By our choice of $n$, we have $h \sigma \in \mathcal{B}_{\alpha n+1} \subseteq \mathcal{B}_{r}$ for $h \in \mathcal{B}_{\alpha n}$ and $\sigma \in S$. So, $\|\phi(h \sigma)-\phi(h)\|_{L_{1}(\mathbb{R})} \leqslant d_{W}(h \sigma, h)=1$, by (9). The right-hand side of (15) is therefore at most a universal constant multiple of $\left|\mathcal{B}_{\alpha n}\right| \cdot\left|\mathfrak{S}_{2}\right| \lesssim(\alpha n)^{6} \lesssim r^{6}$. Hence, by combining (17) with (15) we obtain that the desired estimate (10) indeed holds true.

We have thus shown that all of the new results that were stated above follow from Theorem 3. The bulk of the ensuing discussion will therefore be devoted to the proof of Theorem 3. Prior to doing so, we shall now conclude the introduction by explaining the analytic context of inequality (13) and describing some interesting (and likely quite challenging) questions that remain open.

Fix $q \in[2, \infty)$. A Banach space $\left(X,\|\cdot\|_{X}\right)$ is said to be uniformly convex if for every $\varepsilon \in(0,1)$ there exists $\delta \in(0,1)$ such that $\|x+y\|_{X} \leqslant 2(1-\delta)$ for every $x, y \in X$ that satisfy $\|x\|_{X}=\|y\|_{X}=1$ and $\|x-y\|_{X} \geqslant \varepsilon$. If one can take $\delta \gtrsim X \varepsilon^{q}$ here, then $\left(X,\|\cdot\|_{X}\right)$ is said to have a modulus of uniform convexity of power-type $q$. An important theorem of Pisier [102] asserts that in the setting of uniform convexity (of Banach spaces), power-type behavior is automatic; i.e., every uniformly convex space admits an equivalent norm whose modulus of uniform convexity is of power-type $q$ for some $q \in[2, \infty)$. By [66], for every $k \in \mathbb{N}, p>1$ and $q \geqslant 2$, if $\left(X,\|\cdot\|_{X}\right)$ has a modulus of uniform convexity of power-type $q$, then every finitely supported mapping $\phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow X$ satisfies

$$
\begin{align*}
\left(\sum_{t=1}^{\infty} \frac{1}{t^{1+\frac{q}{2}}}\right. & \left.\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left\|\phi\left(h c^{t}\right)-\phi(h)\right\|_{X}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}}  \tag{18}\\
& \quad \sum_{X, p, q, k}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\phi(h \sigma)-\phi(h)\|_{X}^{p}\right)^{\frac{1}{p}}
\end{align*}
$$

The special case of (18) when $X=\mathbb{R}$ and $p=q=2$ is due to [8], where the quadratic nature of the inequality allows for its proof using representation theory. The general case of (18) was proved in [66] using a semigroup argument that relies on vector-valued Littlewood-Paley-Stein theory. We remark that both [8] and [66] treat only the case $k=1$ of (18) but the proofs in [8], [66] carry over effortlessly to general $k \in \mathbb{N}$; alternatively, a simple argument shows that one can formally deduce (18) for any $k \in \mathbb{N}$ from its validity for $k=1$, as explained in Lemma 32 below.

If $X=\mathbb{R}$ and $q=2$, then (18) says that every finitely supported $\phi:$ $\mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow \mathbb{R}$ satisfies the following estimate for every $p>1$ :

$$
\begin{align*}
& \left(\sum_{t=1}^{\infty} \frac{1}{t^{2}}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left|\phi\left(h c^{t}\right)-\phi(h)\right|^{p}\right)^{\frac{2}{p}}\right)^{\frac{1}{2}}  \tag{19}\\
& \quad \lesssim_{p, k}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}|\phi(h \sigma)-\phi(h)|^{p}\right)^{\frac{1}{p}}
\end{align*}
$$

Thus, the equivalent formulation (13) of Theorem 3 (when $k=2$ ) is the endpoint case $p=1$ of (19). Equation (18) asserts that the following operator is
bounded from $W^{1, p}\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} ; X\right)$ to $\ell_{q}\left(\ell_{p}\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} ; X\right)\right)$ :
$\forall \phi \in W^{1, p}\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} ; X\right), \forall(h, t) \in \mathbb{H}_{\mathbb{Z}}^{2 k+1} \times \mathbb{N}, T \phi(h, t) \stackrel{\text { def }}{=} \frac{1}{t^{\frac{1}{2}+\frac{1}{q}}}\left(\phi\left(h c^{t}\right)-\phi(h)\right)$.
Here, $W^{1, p}\left(\mathbb{H}_{Z}^{2 k+1} ; X\right)$ denotes the (discrete) $X$-valued Sobolev space on $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ that is induced by the generators $\mathfrak{S}_{k}$, i.e., the space of all $\phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow X$ for which the following semi-norm is finite:

$$
\|\phi\|_{W^{1, p}\left(\mathbb{H}_{Z}^{2 k+1} ; X\right)} \stackrel{\text { def }}{=}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\phi(h \sigma)-\phi(h)\|_{X}^{p}\right)^{\frac{1}{p}} .
$$

Also, $\ell_{q}\left(\ell_{p}\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} ; X\right)\right)$ is the space of all $\mathfrak{a}: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \times \mathbb{N} \rightarrow X$ for which the following norm is finite.

$$
\|\mathfrak{a}\|_{\ell_{q}\left(\ell_{p}\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} ; X\right)\right)} \stackrel{\text { def }}{=}\left(\sum_{t=1}^{\infty}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\|\mathfrak{a}(h, t)\|_{X}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}} .
$$

Our result is that when $k \geqslant 2, T$ is actually bounded from $W^{1,1}\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} ; \mathbb{R}\right)$ to $\ell_{2}\left(\ell_{1}\left(\mathbb{H}_{Z}^{2 k+1} ; \mathbb{R}\right)\right)$. This assertion suffices for the geometric and algorithmic applications that are established here, but since our proof relies heavily on the fact that we are dealing with real-valued functions, the availability of such an endpoint estimate in the vector-valued setting remains an elusive open question that is stated explicitly below. This question is important because its positive solution would probably involve the introduction of a markedly new approach that is likely to be valuable elsewhere.

Question 11. Fix $q \geqslant 2$. Suppose that $X$ is a Banach space whose modulus of uniform convexity has power-type $q$. Does there exist $k=k(X) \in \mathbb{N}$ for which the operator $T$ that is defined in (20) is bounded from $W^{1,1}\left(\mathbb{H}_{Z}^{2 k+1} ; X\right)$ to $\ell_{q}\left(\ell_{1}\left(\mathbb{H}_{Z}^{2 k+1} ; X\right)\right)$ ? Does $k=2$ suffice here?

Remark 12. In Remark 4, we described our ongoing work on the case $k=1$ of Theorem 3, which can be rephrased as follows in terms of the operator $T$ that is defined in (20). By Lemma 28 below, the isoperimetric-type inequality (5) is equivalent to the assertion that there exists $q>2$ for which $T$ is bounded from $W^{1,1}\left(\mathbb{H}_{\mathbb{Z}}^{3} ; \mathbb{R}\right)$ to $\ell_{q}\left(\ell_{1}\left(\mathbb{H}_{\mathbb{Z}}^{3} ; \mathbb{R}\right)\right)$. Also, as we described in Remark 4, if $T$ is bounded from $W^{1,1}\left(\mathbb{H}_{Z}^{3} ; \mathbb{R}\right)$ to $\ell_{q}\left(\ell_{1}\left(\mathbb{H}_{Z}^{3} ; \mathbb{R}\right)\right)$, then necessarily $q \geqslant 4$. Thus, even though by [66] we know that $T$ is bounded from $W^{1, p}\left(\mathbb{H}_{Z}^{3} ; \mathbb{R}\right)$ to $\ell_{2}\left(\ell_{p}\left(\mathbb{H}_{Z}^{3} ; \mathbb{R}\right)\right)$ for every $p>1$, at the endpoint $p=1$ the operator $T$ is unbounded from $W^{1,1}\left(\mathbb{H}_{Z}^{3} ; \mathbb{R}\right)$ to $\ell_{2}\left(\ell_{1}\left(\mathbb{H}_{Z}^{3} ; \mathbb{R}\right)\right)$, but the infimum over those $q>2$ for which $T$ is bounded from $W^{1,1}\left(\mathbb{H}_{\mathbb{Z}}^{3} ; \mathbb{R}\right)$ to $\ell_{q}\left(\ell_{1}\left(\mathbb{H}_{\mathbb{Z}}^{3} ; \mathbb{R}\right)\right)$ is finite (and is at least 4). We have evidence that suggests that this infimum is equal to 4 , but at present we do not have a proof of this statement. Regardless, the above results already establish that in three dimensions the boundedness of $T$ exhibits a somewhat
curious jump discontinuity as $p \rightarrow 1$. This phenomenon has some geometric applications of independent interest (to dimensionality reduction) that we will present in forthcoming work.

If one could somehow construct a Banach space $X$ for which Question 11 has a negative answer, then given that we establish here that this question has a positive answer when, say, $X=\ell_{p}$ for $p \in[1,2]$ and $q=2$ (this follows from (14), since $\ell_{p}$ is isometric to a subspace of $L_{1}(\mathbb{R})$, by [53]), it would then be interesting to characterize intrinsically the class of Banach spaces for which the answer to Question 11 is positive. (At present, the answer to Question 11 is unknown when $X=\ell_{p}$ for some $p \in(2, \infty)$ and $q=p$.) In a related vein, we ask the following intriguing question.

Question 13. Let $\mathcal{F}_{\mathbb{H}}$ denote the class of all those Banach spaces $X$ into which $\mathbb{H}_{\mathbb{Z}}^{3}$ does not admit a bi-Lipschitz embedding. ${ }^{1}$ Could $\mathcal{F}_{\mathbb{H}}$ be described using intrinsic geometric properties?

By [22], [69] (see Remark 15 below), $\mathcal{F}_{\text {H }}$ contains all the Banach spaces that admit an equivalent uniformly convex norm. By [23], $\mathcal{F}_{\text {H }}$ also contains all the $L_{1}(\mu)$ spaces. Due to these results, a natural guess for an answer to Question 13 would be that $\mathcal{F}_{\mathbb{H}}$ coincides with the Banach spaces that have finite cotype; see, e.g., [81]. If true, this would be a remarkable geometric result, since (by the Maurey-Pisier theorem [82]) it would mean that the existence of a bi-Lipschitz embedding of $\mathbb{H}_{Z}^{3}$ into a Banach space $X$ implies that every finite metric space embeds into $X$ with bi-Lipschitz distortion $1+\varepsilon$ for every $\varepsilon>0$. As a concrete example of a classical Banach space for which it is unknown whether or not it belongs to $\mathcal{F}_{\text {H }}$, consider the Schatten-von Neumann trace class $S_{1}$ (the space of all operators on $\ell_{2}$, equipped with the nuclear norm; see, e.g., [111, §III.G]). We suspect that $S_{1} \in \mathcal{F}_{\mathbb{H}}$, and that moreover any embedding of $\mathcal{B}_{n}$ into $S_{1}$ incurs bi-Lipschitz distortion that is at least a universal constant multiple of $\sqrt{\log n}$. However, for this strengthening of Theorem 6 to hold true, one would need to find a way to prove it without relying on an isoperimetric-type inequality as we do here. The proofs in [23], [24], [26] that $\mathbb{H}^{3}$ does not admit a bi-Lipschitz embedding into $L_{1}(\mathbb{R})$ also rely on special properties of real-valued functions through the reduction to questions about subsets of finite perimeter in the continuous Heisenberg group. (See Section 2 below.) So, perhaps even as a step towards Question 13 in its full generality, it would be of great interest to

[^1]devise an approach that applies to mappings that take value in $S_{1}$ rather than $L_{1}(\mathbb{R})$.

By a fundamental theorem of Ostrovskii [98], $X \in \mathcal{F}_{\text {HI }}$ if and only if the bi-Lipschitz distortion of any embedding of $\mathcal{B}_{n}$ into $X$ tends to $\infty$ as $n \rightarrow \infty$. The following question asks for a quantitative refinement of this assertion, motivated by similar dichotomic phenomena that occur in metric embedding (see [84], [85], [91], [86], [3]); by [8] (in combination with Pisier's renorming theorem [102]), its answer is positive when $X$ admits an equivalent uniformly convex norm (see [73], [66] for different proofs of this fact), and by [26] its answer is also positive when $X=L_{1}(\mathbb{R})$.

Question 14. Suppose that $X \in \mathcal{F}_{\text {H }}$. Does this imply that there exists $\theta=\theta(X)>0$ such that any embedding of $\mathcal{B}_{n}$ into $X$ incurs bi-Lipschitz distortion at least $(\log n)^{\theta}$ ?

Roadmap. The introduction contained a description of the new results that are established here, but we did not yet present an overview of the ideas that go into our proof of Theorem 3. In particular, we did not yet explain how the fact that the underlying Heisenberg group is of dimension at least 5 becomes relevant to the proof of Theorem 3, despite the fact that this dimension had no role whatsoever in the precursor [66] of this result, i.e., its $\ell_{p}$-version (19) for $p>1$. The reason why we are postponing these (important) explanations is that, even though the discrete setting that was described above is needed for applications, our proof of Theorem 3 actually takes place in a continuous setting that requires the presentation of additional concepts and basic facts. We therefore postpone the overview of the steps of the proof of Theorem 3 to Section 4, which follows Section 2, where various concepts related to the continuous Heisenberg group are presented, and Section 3, where initial reductions are performed, including a reduction of Theorem 3 to its continuous counterpart. With this groundwork in place, it becomes more natural to explain the ideas of our proof in Section 4. Since many readers may be familiar with the (standard) continuous setting, those who wish to see the proof overview can read Section 2.2, which defines the vertical perimeter in the continuous setting, then skip to Section 4 on first reading, though prior to doing so we recommend familiarization with the notation in Section 2.3 since it treats the (somewhat less-standard) notion of intrinsic Lipschitz graph, which has a central role in our proof.

In Section 5 we derive, as a crucial new ingredient of our proof of Theorem 3, a special case of the (continuous counterpart of) the vertical-versushorizontal isoperimetric inequality of Theorem 3. It is important to stress that this argument of Section 5 is the only place where the assumption $k>1$ is used in our proof of Theorem 3. The remaining steps of the proof work for

Heisenberg groups of any dimension, yielding structural information that will be used in future work also when $k=1$.

Section 6 contains definitions and basic facts related to (local) intrinsic corona decompositions, which are structural properties of sets in the continuous Heisenberg group (existence of certain well-behaved multi-scale covers) that will be used to deduce our isoperimetric-type inequality in its full generality from the special case established in Section 5. In Section 7 we prove that sets that admit an intrinsic corona decomposition satisfy the desired isoperimetrictype inequality.

In Section 8 we state Theorem 57, which is a technical structural result asserting that any subset $E$ of the continuous Heisenberg group such that $E$, the complement of $E$ and the boundary of $E$ are all locally Ahlfors regular, admits a suitable local intrinsic corona decomposition. Thus, by the results of Section 7, such sets satisfy the desired isoperimetric-type inequality. To apply this fact, in Section 8 we also show how to decompose a cellular set (see Section 2.1.6) of finite perimeter in the continuous Heisenberg group into parts that are locally Ahlfors regular as above, in such a way that if we sum up the isoperimetric-type inequalities that follow from Theorem 57 for each of these parts, then we obtain the desired inequality for the initial cellular set.

As part of the basic reductions that are contained in Section 3, a simple approximation argument shows that it suffices to treat cellular sets. (See Lemma 40 below.) Therefore, Section 8 implies that it remains to prove Theorem 57, i.e., to construct a local intrinsic corona decomposition for every set that satisfies the above local Ahlfors-regularity (i.e., for the set itself, its complement, and its boundary). This construction is performed in Section 9, thus completing the proof of Theorem 3.

We conclude this article with Section 10, which contains further historical background on the Sparsest Cut Problem, as well as descriptions of directions for future research (some of which we will pursue in forthcoming works).

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## 2. The continuous setting

In what follows, it will be beneficial to allow only one exception to the convention for asymptotic notation that was described after Theorem 3. Specifically, since throughout we will fix an integer $k \in \mathbb{N}$ and many constant factors do depend on $k$, in order to not overburden the notation we shall allow the notation $\lesssim, \gtrsim, \asymp$ to coincide with $\lesssim_{k}, \gtrsim_{k}, \asymp_{k}$, respectively. Nevertheless, our main application is when $k=2$, so in this case the implicit constant factors are in fact universal constants.
2.1. Definition of continuous Heisenberg group. Fix $k \in \mathbb{N}$. In light of the matrix realization (1) that we chose in the introduction for the discrete Heisenberg group $\mathbb{H}_{Z}^{2 k+1}$, an obvious way to define the continuous Heisenberg group $\mathbb{H}^{2 k+1}$ is that $\mathbb{H}^{2 k+1}$ consists of all those matrices as in (1), but now with the entries $x_{1}, \ldots, x_{k}, y_{1}, \ldots, y_{k}, z$ allowed to be arbitrary real numbers. This would work just fine in the ensuing discussion, but for notational convenience we prefer to consider a different realization of the Heisenberg group that arises by identifying the above real matrix group with its Lie algebra via the exponential map and taking the Baker-Campbell-Hausdorff formula as the definition of the group product. We shall now describe the continuous Heisenberg group in this (standard) way and proceed to adhere exclusively to this specific realization in what follows.

Fix a Hilbertian norm $\|\cdot\|$ on $\mathbb{R}^{2 k+1}$. Also fix an orthonormal basis $X_{1}, \ldots, X_{k}, Y_{1}, \ldots, Y_{k}, Z$ of $\mathbb{R}^{2 k+1}$, so that every $h \in \mathbb{R}^{2 k+1}$ can be written as $h=\sum_{i=1}^{k} \alpha_{i} X_{i}+\sum_{i=1}^{k} \beta_{i} Y_{i}+\gamma Z$ for some $\alpha_{1}, \ldots, \alpha_{k}, \beta_{1}, \ldots, \beta_{k}, \gamma \in \mathbb{R}$. Denote $x_{i}(h)=\alpha_{i}, y_{i}(h)=\beta_{i}$ for all $i \in\{1, \ldots, k\}$ and $z(h)=\gamma$; i.e., $x_{1}, \ldots, x_{k}$, $y_{1}, \ldots, y_{k}, z: \mathbb{R}^{2 k+1} \rightarrow \mathbb{R}$ are the coordinate functions corresponding to the basis $X_{1}, \ldots, X_{k}, Y_{1}, \ldots, Y_{k}, Z$. We shall also write $x(h)=\sum_{i=1}^{k} x_{i}(h) X_{i}$, $y(h)=\sum_{i=1}^{k} y_{i}(h) Y_{i}$ and $\pi(h)=x(h)+y(h)$ (thus $\left.\pi(Z)=0\right)$. In what follows, we shall canonically identify the span of $\left\{X_{1}, \ldots, X_{k}, Y_{1}, \ldots, Y_{k}\right\}$ with $\mathbb{R}^{2 k}$, so that the mapping $\pi: \mathbb{R}^{2 k+1} \rightarrow \mathbb{R}^{2 k}$ is the orthogonal projection onto the first $2 k$ coordinates.

For every $u, v \in \mathbb{R}^{2 k+1}$, write $\boldsymbol{\omega}(u, v)=\sum_{i=1}^{k}\left(x_{i}(u) y_{i}(v)-y_{i}(u) x_{i}(v)\right)$; i.e., $\omega(u, v)$ is the standard symplectic form on $\mathbb{R}^{2 k}$ applied to the vectors $\pi(u), \pi(v)$. In particular, $\boldsymbol{\omega}\left(X_{i}, Y_{i}\right)=-\omega\left(Y_{i}, X_{i}\right)=1$, and for every two basis elements $u, v \in\left\{X_{1}, \ldots, X_{k}, Y_{1}, \ldots, Y_{k}, Z\right\}$ such that $\{u, v\} \neq\left\{X_{i}, Y_{i}\right\}$ for all $i \in\{1, \ldots, k\}$, we have $\omega(u, v)=0$. The (continuous) Heisenberg group $\mathbb{H}^{2 k+1}$ is defined to be the group whose underlying set is $\mathbb{R}^{2 k+1}$, equipped with the following product:

$$
\begin{equation*}
\forall u, v \in \mathbb{H}^{2 k+1}, \quad u v \stackrel{\text { def }}{=} u+v+\frac{\omega(u, v)}{2} Z . \tag{21}
\end{equation*}
$$

It is straightforward to check that this turns $\mathbb{H}^{2 k+1}$ into a group whose identity element is the all- 0 vector $\mathbf{0} \in \mathbb{H}^{2 k+1}$ and the inverse of $h \in \mathbb{H}^{2 k+1}$ is $-h$. The resulting group is isomorphic to the real matrix group (with usual matrix multiplication) that we described above. Indeed, consider the isomorphism that assigns to every $h \in \mathbb{R}^{2 k+1}$ the following upper triangular $k+2$ by $k+2$ matrix:

$$
\psi(h) \stackrel{\text { def }}{=}\left(\begin{array}{cccccc}
1 & x_{1}(h) & x_{2}(h) & \ldots & x_{k}(h) & w(h) \\
0 & 1 & 0 & \ldots & 0 & y_{1}(h) \\
\vdots & \ddots & \ddots & \ddots & \vdots & \vdots \\
\vdots & \ddots & \ddots & \ddots & 0 & y_{k-1}(h) \\
\vdots & \ddots & \ddots & \ddots & 1 & y_{k}(h) \\
0 & \ldots & \ldots & \ldots & 0 & 1
\end{array}\right)
$$

where

$$
w(h) \stackrel{\text { def }}{=} z(h)+\frac{1}{2} \sum_{i=1}^{k} x_{i}(h) y_{i}(y)
$$

The discrete Heisenberg group $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ can now be realized as the subgroup of $\mathbb{H}^{2 k+1}$ that is generated by $\left\{X_{1}, \ldots, X_{k}, Y_{1}, \ldots, Y_{k}\right\}$. Note that the $z$-coordinate of an element in $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ is either an integer or a half-integer. From now on, we shall work exclusively with this specific realization of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$.

In order to avoid confusing multiplication by scalars with the group law of $\mathbb{H}^{2 k+1}$, for every $h=\sum_{i=1}^{k} \alpha_{i} X_{i}+\sum_{i=1}^{k} \beta_{i} Y_{i}+\gamma Z \in \mathbb{H}^{2 k+1}$ and $t \in \mathbb{R}$, it will be convenient to use the exponential notation $h^{t}=\sum_{i=1}^{k} t \alpha_{i} X_{i}+\sum_{i=1}^{k} t \beta_{i} Y_{i}+t \gamma Z$. One directly computes from (21) that for every $i \in\{1, \ldots, k\}$ and $s, t \in \mathbb{R}$, we have $\left[X_{i}^{s}, Y_{i}^{t}\right]=Z^{s t}$, where we use the standard commutator notation $[u, v]=u v u^{-1} v^{-1}$ for every $u, v \in \mathbb{H}^{2 k+1}$. Given a subset $A \subseteq \mathbb{H}^{2 k+1}$ we shall denote its linear span by $\langle A\rangle$. Thus, $\langle h\rangle=\langle\{h\}\rangle=\left\{h^{t}: t \in \mathbb{R}\right\}$ is the one-parameter subgroup generated by $h \in \mathbb{H}^{2 k+1}$.
2.1.1. Horizontal derivatives. If $f: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ is smooth, then for every $i \in\{1, \ldots, k\}$, its horizontal derivatives in directions $X_{i}, Y_{i}$ are defined by setting for every $h \in \mathbb{H}^{2 k+1}$,

$$
\begin{aligned}
& X_{i} f(h) \stackrel{\text { def }}{=} \frac{\partial f}{\partial x_{i}}(h)-\frac{1}{2} y_{i}(h) \frac{\partial f}{\partial z}(h) \quad \text { and } \\
& Y_{i} f(h) \stackrel{\text { def }}{=} \frac{\partial f}{\partial y_{i}}(h)+\frac{1}{2} x_{i}(h) \frac{\partial f}{\partial z}(h)
\end{aligned}
$$

These are differential operators corresponding to left-invariant horizontal vector fields on $\mathbb{H}^{2 k+1}$, so they are left-invariant. That is, if we let $\rho_{g}$ be the left action $\rho_{g} f(h)=f\left(g^{-1} h\right)$ for all $g \in \mathbb{H}^{2 k+1}$, then $\left\{X_{i}\right\}_{i=1}^{k}$ and $\left\{Y_{i}\right\}_{i=1}^{k}$ commute with $\rho_{g}$.

The horizontal gradient of $f$ is then defined as follows:
$\forall h \in \mathbb{H}^{2 k+1}, \quad \nabla_{\mathbb{H}} f(h) \stackrel{\text { def }}{=}\left(X_{1} f(h), \ldots, X_{k} f(h), Y_{1} f(h), \ldots, Y_{k} f(h)\right) \in \mathbb{R}^{2 k}$.
Thus, for $p \in(0, \infty]$, the $p$-norm of the horizontal gradient of $f$ is given by

$$
\begin{aligned}
& \left\|\nabla_{\mathbb{H}} f(h)\right\|_{\ell_{p}^{2 k}}=\left(\sum_{i=1}^{k}\left(\left|X_{i} f(h)\right|^{p}+\left|Y_{i} f(h)\right|^{p}\right)\right)^{\frac{1}{p}} \\
& \quad=\left(\sum_{i=1}^{k}\left|\frac{\partial f}{\partial x_{i}}(h)-\frac{1}{2} y_{i}(h) \frac{\partial f}{\partial z}(h)\right|^{p}+\sum_{i=1}^{k}\left|\frac{\partial f}{\partial y_{i}}(h)+\frac{1}{2} x_{i}(h) \frac{\partial f}{\partial z}(h)\right|^{p}\right)^{\frac{1}{p}} .
\end{aligned}
$$

2.1.2. Lines and planes. From now on denote $\mathrm{H} \stackrel{\text { def }}{=}\left\langle X_{1}, \ldots, X_{k}, Y_{1}, \ldots, Y_{k}\right\rangle$ $=\pi\left(\mathbb{H}^{2 k+1}\right)$. This $2 k$-dimensional linear subspace H is called the subspace of horizontal vectors. If $U \subseteq \mathrm{H}$ is an affine hyperplane in H (thus $U$ is a translate of a linear subspace of H of dimension $2 k-1$ ), then we call $V=\pi^{-1}(U) \subseteq \mathbb{H}^{2 k+1}$ the vertical plane lying over $U$.

A horizontal line in $\mathbb{H}^{2 k+1}$ is a coset of the form $w\langle h\rangle$ for some $w \in \mathbb{H}^{2 k+1}$ and $h \in \mathrm{H}$. If $u \in \mathbb{H}^{2 k+1}$, then $P_{u}=u \mathrm{H}$ is called the horizontal plane centered at $u$. Equivalently, $P_{u}$ is the union of all of the horizontal lines that pass through $u$. We can write

$$
\begin{equation*}
P_{u}=\left\{h Z^{z(u)+\frac{\omega(\pi(u), \pi(h))}{2}}: h \in \mathrm{H}\right\} . \tag{22}
\end{equation*}
$$

(See Section 9.4.1.)
The term plane in $\mathbb{H}^{2 k+1}$ will always stand for either a vertical plane or a horizontal plane. By (22), any plane in $\mathbb{H}^{2 k+1}$ is an affine hyperplane in $\mathbb{R}^{2 k+1}$. Since $\omega$ is nondegenerate on $\mathbf{H}$, for any affine hyperplane $Q \subseteq \mathbb{R}^{2 k+1}$, either $Q$ is vertical or there is a unique $u \in \mathbb{H}^{2 k+1}$ such that $Q=P_{u}$. Every plane $P \subseteq$ $\mathbb{H}^{2 k+1}$ separates $\mathbb{H}^{2 k+1}$ into two half-spaces, which we denote by $P^{+}$and $P^{-}$.
2.1.3. The Carnot-Carathéodory metric. Suppose that $a, b \in \mathbb{R}$ satisfy $a<b$. If $\gamma:[a, b] \rightarrow \mathbb{H}^{2 k+1}$ is a curve such that the coordinate functions $\left\{x_{i} \circ \gamma\right\}_{i=1}^{k},\left\{y_{i} \circ \gamma\right\}_{i=1}^{k}, z \circ \gamma$ are all Lipschitz, then the tangent vector $\gamma^{\prime}(t) \in$ $\mathbb{H}^{2 k+1}$ is defined for almost all $t \in[a, b]$. We say that $\gamma$ is a horizontal curve if $\gamma^{\prime}(t) \in P_{\gamma(t)}$ for almost every $t \in[a, b]$. Equivalently, for $\gamma$ to be horizontal, we require that $\gamma^{-1} \gamma^{\prime} \in \mathrm{H}$ almost everywhere. The length of a horizontal curve $\gamma:[a, b] \rightarrow \mathbb{H}^{2 k+1}$ is defined to be $\ell(\gamma)=\int_{a}^{b}\left\|\pi\left(\gamma^{\prime}(t)\right)\right\| \mathrm{d} t=\ell(\pi \circ \gamma)$. Note that segments of horizontal lines are horizontal curves. Since $\left[X_{i}^{s}, Y_{i}^{t}\right]=Z^{s t}$ for every $i \in\{1, \ldots, k\}$ and $s, t \in \mathbb{R}$, any two points in $\mathbb{H}^{2 k+1}$ can be connected by a horizontal curve consisting of segments of horizontal lines. We can therefore define the Carnot-Carathéodory distance between any two points $v, w \in \mathbb{H}^{2 k+1}$ by

$$
d(v, w) \stackrel{\text { def }}{=} \inf \left\{\ell(\gamma): \gamma:[0,1] \rightarrow \mathbb{H}^{2 k+1} \text { is horizontal, } \gamma(0)=v, \gamma(1)=w\right\} .
$$

See, e.g., [45], [89] for more on this metric; it suffices to recall that there exists $C_{k} \in[1, \infty)$ such that
(23) $\forall h \in \mathbb{H}^{2 k+1}$,

$$
d(\mathbf{0}, h) \leqslant \sum_{i=1}^{k}\left|x_{i}(h)\right|+\sum_{i=1}^{k}\left|y_{i}(h)\right|+4 \sqrt{|z(h)|} \leqslant C_{k} d(\mathbf{0}, h)
$$

and

$$
\begin{equation*}
\forall h \in \mathbb{H}^{2 k+1}, \quad d(\mathbf{0}, h) \geqslant\|\pi(h)\| . \tag{24}
\end{equation*}
$$

In what follows, for every $r \in(0, \infty)$, we shall denote by $B_{r} \subseteq \mathbb{H}^{2 k+1}$ the open ball in the metric $d$ that is centered at $\mathbf{0}$, i.e., $B_{r}=\left\{h \in \mathbb{H}^{2 k+1}\right.$ : $d(h, \mathbf{0})<r\}$. For every $\Omega \subseteq \mathbb{H}^{2 k+1}$, the Lipschitz constant of a mapping $f: \Omega \rightarrow \mathbb{R}$ relative to the metric $d$ will be denoted by $\|f\|_{\text {Lip }(\Omega)}$. For every $s \in(0, \infty)$, the $s$-dimensional Hausdorff measure that is induced by the metric $d$ on $\mathbb{H}^{2 k+1}$ will be denoted below by $\mathscr{H}^{s}$. Note that it follows from (23) that $\mathbb{H}^{2 k+1}$ is a metric space with Hausdorff dimension $2 k+2$. One checks that $d$ is a left-invariant metric on $\mathbb{H}^{2 k+1}$ and that the Lebesgue measure is a Haar measure of $\mathbb{H}^{2 k+1}$. Thus $\mathscr{H}^{2 k+2}$ is proportional to the Lebesgue measure. Also, for every $h \in \mathbb{H}^{2 k+1}$ and $r \in[0, \infty)$, the open ball of radius $r$ in the metric $d$ is $h B_{r}=B_{r}(h)$.

If $h \in \mathbb{H}^{2 k+1}$ and $v \in \mathbf{H} \backslash\{\mathbf{0}\}$, then the horizontal line $h\langle v\rangle$ is a geodesic. Consequently,

$$
\begin{equation*}
d(h, h v)=d_{\mathrm{Euc}}(\pi(h), \pi(h v))=\|v\|, \tag{25}
\end{equation*}
$$

where $d_{\text {Euc }}(\cdot, \cdot)$ denotes the Euclidean metric that is induced by the Hilbertian norm $\|\cdot\|$. Furthermore, by the classical isoperimetric inequality for curves in $\mathbb{R}^{2}$, one can show (e.g., [89]) that

$$
\forall r \in(0, \infty), \quad d\left(\mathbf{0}, Z^{\pi r^{2}}\right)=2 \pi r .
$$

2.1.4. Isometries and scaling automorphisms. The isometry group of $\mathbb{H}^{2 k+1}$ acts transitively on $\mathbb{H}^{2 k+1}$ by left-multiplication. There are also many isometries that fix the identity element $\mathbf{0}$. Any element $v \in \mathbb{H}^{2 k+1}$ can be written uniquely as $v=h Z^{t}$ for some $h \in \mathrm{H}$ and $t \in \mathbb{R}$. If $f: \mathbb{R}^{2 k} \rightarrow \mathbb{R}^{2 k}$ is a Euclidean isometry that leaves the symplectic form $\omega$ invariant, then the map $h Z^{t} \mapsto f(h) Z^{t}$ is an isometry of $\mathbb{H}^{2 k+1}$. The group consisting of isometries of this type is isomorphic to the unitary group $U(k)$, and it acts transitively on the unit sphere in H .

Another important class of automorphisms of $\mathbb{H}^{2 k+1}$ is the following family of scalings. For every $t \in(0, \infty)$, define $\mathfrak{s}_{t}: \mathbb{H}^{2 k+1} \rightarrow \mathbb{H}^{2 k+1}$ by

$$
\mathfrak{s}_{t}\left(\sum_{i=1}^{k} \alpha_{i} X_{1}+\sum_{i=1}^{k} \beta_{i} Y_{i}+\gamma Z\right) \stackrel{\text { def }}{=} \sum_{i=1}^{k} t \alpha_{i} X_{1}+\sum_{i=1}^{k} t \beta_{i} Y_{i}+t^{2} \gamma Z .
$$

It is straightforward to check that $\mathfrak{s}_{t}$ is an automorphism of $\mathbb{H}^{2 k+1}$, that it sends horizontal curves to horizontal curves, and that if $\gamma$ is a horizontal curve, then $\ell\left(\mathfrak{s}_{t} \circ \gamma\right)=t \ell(\gamma)$. It follows that $\mathfrak{s}_{t}$ acts as a scaling of the Carnot-Carathéodory metric, i.e., $d\left(\mathfrak{s}_{t}(u), \mathfrak{s}_{t}(v)\right)=t d(u, v)$ for all $u, v \in \mathbb{H}^{2 k+1}$.

Remark 15. Analogously to Question 13, one can consider the class $\mathcal{F}_{H i}^{\mathbb{R}}$ of those Banach spaces $X$ into which the continuous Heisenberg group $\mathbb{H}^{3}$ does not admit a bi-Lipschitz embedding. By [69], [22], $\mathcal{F}_{\mathbb{R}}^{\mathbb{R}}$ contains all the Banach spaces that have the Radon-Nikodým property (RNP) hence, in particular, all the reflexive Banach spaces and all the separable dual Banach spaces. (See, e.g., [13, Ch. 5] for these assertions, as well as much more on the RNP.) By [23], $\mathcal{F}_{\mathbb{H}}^{\mathbb{R}}$ also contains all the $L_{1}(\mu)$ spaces. We ask whether $X \in \mathcal{F}_{H}^{\mathbb{R}}$ implies that also $L_{1}(\mu ; X) \in \mathcal{F}_{H}^{\mathbb{R}}$; if true, this would be a beautiful strengthening of [23]. We also do not know whether $\mathcal{F}_{\mathbb{H}}^{\mathbb{R}}$ contains all the noncommutative $L_{1}$ spaces (see, e.g., [103]), but since the Schatten-von Neumann trace class $S_{1}$ has the RNP (since $S_{1}$ is separable and it is the dual of the space of all the compact operators on $\ell_{2}$, equipped with the operator norm), by [69], [22] we do know that $S_{1} \in \mathcal{F}_{\mathbb{R}}^{\mathbb{R}}$. (Recall that it is unknown whether or not $S_{1}$ belongs to $\mathcal{F}_{\mathbb{H}}$.) Since $\mathbb{H}_{\mathbb{Z}}^{3}$ is bi-Lipschitz equivalent to a subset of $\mathbb{H}^{3}$, we trivially have $\mathcal{F}_{\text {Hi }} \subseteq \mathcal{F}_{\mathbb{R}}^{\mathbb{R}}$. This inclusion is strict because by [12] we know that $\mathbb{H}_{\mathbb{Z}}^{3}$ (indeed, any locally-finite metric space) admits a bi-Lipschitz embedding into a Banach space $Z$ with the RNP. (In fact, one can take $Z$ to be the $\ell_{2}$ direct sum of the finite dimensional spaces $\left\{\ell_{\infty}^{n}\right\}_{n=1}^{\infty}$, so $Z$ can even be reflexive.) By considering the re-scaled copies $\mathfrak{s}_{t}\left(\mathbb{H}_{\mathbb{Z}}^{3}\right) \subseteq \mathbb{H}^{3}$ as $t \rightarrow 0^{+}$, one sees that if $X$ belongs to $\mathcal{F}_{\text {H }}$, then for any nonprincipal ultrafilter $U$, the ultrapower $X^{u}$ belongs to $\mathcal{F}_{H}^{\mathbb{R}}$; see, e.g., [49] for ultrapowers of Banach spaces. By [17], $X^{u}$ has the RNP if and only if $X$ admits an equivalent uniformly convex norm. Hence, it follows from [69], [22] that $\mathcal{F}_{\text {HI }}$ contains all the Banach spaces that admit an equivalent uniformly convex norm. (Different proofs of this fact were obtained in [8], [73], [66].) As in Question 13, it is natural to ask for an intrinsic geometric characterization of the class $\mathcal{F}_{\mathbb{R}}^{\mathbb{R}}$. We chose to highlight Question 13 rather than this continuous variant due to its local nature, i.e., because by [98] the balls $\mathcal{B}_{n}$ embed with distortion $O_{X}(1)$ (as $n \rightarrow \infty$ ) into a Banach space $X$ if and only if $X \in \mathcal{F}_{\mathrm{H}}$. In other words, membership in $\mathcal{F}_{\mathrm{H}}$ is determined by the geometry of finite subsets of the Banach space in question while membership in $\mathcal{F}_{\mathbb{H}}^{\mathbb{R}}$ is an inherently infinite dimensional property.
2.1.5. Measure and perimeter. If $E \subseteq \mathbb{H}^{2 k+1}$ has smooth or piecewise smooth boundary, then $\partial E$ has Hausdorff dimension $2 k+1$. In this paper we will primarily deal with sets $E \subseteq \mathbb{H}^{2 k+1}$ such that $\mathscr{H}^{2 k+1}(\partial E)<\infty$ or surfaces $A \subseteq \mathbb{H}^{2 k+1}$ of topological dimension $2 k+1$ such that $\mathscr{H}^{2 k+1}(A)<\infty$; such surfaces include the boundaries of the cellular sets that are introduced in

Section 2.1.6 below and the intrinsic Lipschitz graphs that are introduced in Section 2.3 below. In Section 9.2, however, we will need to work with limits of cellular sets in order to apply a compactness argument. The boundaries of such limits can be more complicated, and therefore we will need in Section 9.2 the more sophisticated notion of Heisenberg perimeter. See [38] for the definition of this notion, which is also recalled in [23], [26]. We will not need to work here with this definition directly, so it suffices to state that as in [38] one associates to every measurable subset $E \subseteq \mathbb{H}^{2 k+1}$ a measure $\operatorname{Per}_{E}$ that is called the perimeter measure of $E$ and has the properties that are stated in Propositions 16, 17, 18 and 19 below, all of which are proven in [38].

Proposition 16. There exists $c_{k} \in(0, \infty)$ such that if $E \subseteq \mathbb{H}^{2 k+1}$ has piecewise smooth boundary and $U \subseteq \mathbb{H}^{2 k+1}$ is an open set, then $\operatorname{Per}_{E}(U)=$ $c_{k} \mathscr{H}^{2 k+1}(U \cap \partial E)$.

Proposition 17. The perimeter measure is scale-invariant in the sense that for any $E \subseteq \mathbb{H}^{2 k+1}$, any open set $U$ and any $t \in(0, \infty)$, we have $\operatorname{Per}_{\mathfrak{s}_{t}(E)}\left(\mathfrak{s}_{t}(U)\right)=t^{2 k+1} \operatorname{Per}_{E}(U)$.

If $\operatorname{Per}_{E}(B)<\infty$ for every ball $B \subseteq \mathbb{H}^{2 k+1}$, we shall say that $E$ has locally finite perimeter.

Proposition 18. If $E \subseteq \mathbb{H}^{2 k+1}$ has locally finite perimeter, then $\operatorname{Per}_{E}(U)$ $\lesssim \mathscr{H}^{2 k+1}(U \cap \partial E)$ for every open set $U \subseteq \mathbb{H}^{2 k+1}$.

Finally, the following compactness result appears in [42, Th. 1.28] and [38, Th. 2.10].

Proposition 19. Let $\left\{E_{i}\right\}_{i=1}^{\infty}$ be a sequence of measurable subsets of $\mathbb{H}^{2 k+1}$ such that for every ball $B \subseteq \mathbb{H}^{2 k+1}$, we have $\sup _{i \in \mathbb{N}} \operatorname{Per}_{E_{i}}(B)<\infty$. Then there exists a subsequence $\left\{i_{n}\right\}_{n=1}^{\infty} \subseteq \mathbb{N}$ and a measurable $E \subseteq \mathbb{H}^{2 k+1}$ such that $\lim _{n \rightarrow \infty} \mathbf{1}_{E_{i_{n}}}=\mathbf{1}_{E}$ in $L_{1}^{\text {loc }}\left(\mathcal{H}^{2 k+2}\right)$ and $\operatorname{Per}_{E}(U) \leqslant \sup _{n \in \mathbb{N}} \operatorname{Per}_{E_{n_{i}}}(U)$ for every open $U \subseteq \mathbb{H}^{2 k+1}$.
2.1.6. Cellular sets. Recall that the discrete Heisenberg group $\mathbb{H}_{Z}^{2 k+1}$ is the subgroup of $\mathbb{H}^{2 k+1}$ generated by $\left\{X_{1}, Y_{1}, \ldots, X_{k}, Y_{k}\right\}$. The unit cube $\left[-\frac{1}{2}, \frac{1}{2}\right]^{2 k+1}$ is a fundamental domain for the left action of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ on $\mathbb{H}^{2 k+1}$. Its translates tile $\mathbb{H}^{2 k+1}$ by parallelepipeds and give $\mathbb{H}^{2 k+1}$ the structure of a polyhedral complex $\tau^{2 k+1}$. In what follows, we call this complex the unit lattice and say that a subset of $\mathbb{H}^{2 k+1}$ is cellular if it is the union of closed $(2 k+1)$-cells of the unit lattice.

If $E \subseteq \mathbb{H}^{2 k+1}$ is measurable and satisfies $\mathscr{H}^{2 k+1}(\partial E)<\infty$, then $E$ can be approximated by scalings of cellular sets. To justify this (well-known, but not easily quotable) fact, we need to recall the following relative version of
the isoperimetric inequality for $\mathbb{H}^{2 k+1}$. Below, and in what follows, given $E \subseteq \mathbb{H}^{2 k+1}$ we use the notation $E^{c}=\mathbb{H}^{2 k+1} \backslash E$. Also, throughout the ensuing discussion it is convenient to use the following standard notation for normalized integrals. If $\mu$ is a measure on $\mathbb{H}^{2 k+1}$ and $U \subseteq \mathbb{H}^{2 k+1}$ is $\mu$-measurable with $\mu(U)>0$, then for every $f \in L_{1}(U)$, write

$$
f_{U} f \mathrm{~d} \mu \stackrel{\text { def }}{=} \frac{1}{\mu(U)} \int_{U} f \mathrm{~d} \mu .
$$

Lemma 20. There is $\Lambda_{k} \in(0, \infty)$ such that if $E \subseteq \mathbb{H}^{2 k+1}$ is measurable and $r \in(0, \infty)$, then

$$
\begin{equation*}
\left(\frac{\mathscr{H}^{2 k+2}\left(E \cap B_{r}\right) \cdot \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap B_{r}\right)}{\mathscr{H}^{2 k+2}\left(B_{r}\right)^{2}}\right)^{\frac{2 k+1}{2 k+2}} \lesssim \frac{r \mathscr{H}^{2 k+1}\left(\partial E \cap B_{\Lambda_{k} r}\right)}{\mathscr{H}^{2 k+2}\left(B_{\Lambda_{k} r}\right)} . \tag{26}
\end{equation*}
$$

Proof. The Poincaré-Sobolev inequality on $\mathbb{H}^{2 k+1}$ (see [48] or equation (2.5) in [26]) asserts that there exists $\Lambda_{k} \in(0, \infty)$ such that every smooth $f$ : $\mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ satisfies

$$
\begin{align*}
& \left(f_{B_{r} \times B_{r}}|f(x)-f(y)|^{\frac{2 k+2}{2 k+1}} \mathrm{~d} \mathscr{H}^{2 k+2}(x) \mathrm{d} \mathscr{H}^{2 k+2}(y)\right)^{\frac{2 k+1}{2 k+2}}  \tag{27}\\
& \quad \lesssim r f_{B_{\Lambda_{k} r}}\left\|\nabla_{\mathbb{H}} f(x)\right\|_{\ell_{1}^{2 k}} \mathrm{~d} \mathscr{H}^{2 k+2}(x) .
\end{align*}
$$

The desired estimate (26) follows by applying (27) when $f$ is a smooth approximation to $\mathbf{1}_{E}$.

Lemma 21. Let $E \subseteq \mathbb{H}^{2 k+1}$ be a measurable set. Suppose that $\mathscr{H}^{2 k+1}(\partial E)$ $<\infty$. Then for every $\rho \in(0, \infty)$, there is a set $F=F_{\rho}$ with the following properties:

- $\mathfrak{s}_{1 / \mathrm{\rho}}\left(F_{\mathrm{\rho}}\right)$ is a cellular set,
- $\mathscr{H}^{2 k+1}\left(\partial F_{\rho}\right) \lesssim \mathscr{H}^{2 k+1}(\partial E)$,
- $\mathscr{H}^{2 k+2}\left(E \triangle F_{\rho}\right) \lesssim \rho \mathscr{H}^{2 k+1}(\partial E)$.

Here, $A \triangle B \stackrel{\text { def }}{=}(A \backslash B) \cup(B \backslash A)$ is (as usual) the symmetric difference of $A, B \subseteq \mathbb{H}^{2 k+1}$.

Proof. It suffices to prove Lemma 21 when $\rho=1$. Indeed, for general $\rho \in(0, \infty)$, apply the case $\rho=1$ of Lemma 21 with $E$ replaced by $\mathfrak{s}_{1 / \rho}(E)$ to construct a cellular set $F_{0}$, and then take $F=\mathfrak{s}_{\rho}\left(F_{0}\right)$. Consider the following set of cells in the unit lattice $\tau^{2 k+1}$ :

$$
e \stackrel{\text { def }}{=}\left\{C \in \tau^{2 k+1}: \mathscr{H}^{2 k+2}(C \cap E)>\frac{\mathscr{H}^{2 k+2}(C)}{2}\right\} .
$$

Write $F=\bigcup_{C \in \mathcal{C}} C$. Then $F$ is cellular by design, and we shall next show that $F$ satisfies the remaining two desired properties.

Let $A \subseteq \tau^{2 k+1}$ be the set of cells that intersect $\partial F$. Then $\mathscr{H}^{2 k+1}(\partial F)$ $\asymp|A|$. For each $C \in \tau^{2 k+1}$, let $h_{C} \in C$ be the center of $C$. If $C \in A$, then $C$ is adjacent to a cell $D \in \tau^{2 k+1}$ with $|\{C, D\} \cap C|=1$. Let $r=\operatorname{diam}(C)$ $=\operatorname{diam}\left(\left[-\frac{1}{2}, \frac{1}{2}\right]^{2 k+1}\right)$, and let $B=B_{2 r}\left(h_{C}\right)=h_{C} B_{2 r}$, so that $C \cup D \subseteq B$. Then $\mathscr{H}^{2 k+2}(E \cap B) \geqslant \mathscr{H}^{2 k+2}(C) / 2 \asymp \mathscr{H}^{2 k+2}(B)$ and $\mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap B\right) \geqslant$ $\mathscr{H}^{2 k+2}(D) / 2 \asymp \mathscr{H}^{2 k+2}(B)$. Lemma 20 therefore implies $\mathscr{H}^{2 k+1}\left(\partial E \cap B_{2 \Lambda_{k} r}\left(h_{C}\right)\right)$ $\gtrsim 1$. Since the balls $\left\{B_{2 \Lambda_{k} r}\left(h_{C}\right)\right\}_{C \in \tau^{2 k+1}}$ cover $\mathbb{H}^{2 k+1}$ with finite multiplicity, it follows that

$$
\mathscr{H}^{2 k+1}(\partial E) \gtrsim \sum_{C \in A} \mathscr{H}^{2 k+1}\left(\partial E \cap B_{2 \Lambda_{k} r}\left(h_{C}\right)\right) \gtrsim|A| \asymp \mathscr{H}^{2 k+1}(\partial F) .
$$

Next, by the definition of $F$, for every $C \in \tau^{2 k+1}$, we have

$$
\min \left\{\mathscr{H}^{2 k+2}(E \cap C), \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap C\right)\right\}=\mathscr{H}^{2 k+2}((E \triangle F) \cap C),
$$

and also we trivially have

$$
\max \left\{\mathscr{H}^{2 k+2}(E \cap C), \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap C\right)\right\} \geqslant \frac{\mathscr{H}^{2 k+2}(C)}{2} .
$$

Consequently, if we denote $B=B_{r}\left(h_{C}\right)$, then

$$
\begin{aligned}
\mathscr{H}^{2 k+2}(E \cap B) \cdot \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap B\right) & \geqslant \mathscr{H}^{2 k+2}(E \cap C) \cdot \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap C\right) \\
& \geqslant \mathscr{H}^{2 k+2}((E \triangle F) \cap C) \cdot \frac{\mathscr{H}^{2 k+2}(C)}{2}
\end{aligned}
$$

and, using Lemma 20 (and recalling that $r=\operatorname{diam}\left(\left[-\frac{1}{2}, \frac{1}{2}\right]^{2 k+1}\right) \asymp 1$ ),

$$
\mathscr{H}^{2 k+1}\left(\partial E \cap B_{\Lambda_{k} r}\right) \gtrsim \mathscr{H}^{2 k+2}((E \triangle F) \cap C)^{\frac{2 k+1}{2 k+2}} \gtrsim \mathscr{H}^{2 k+2}((E \triangle F) \cap C) .
$$

These estimates imply that

$$
\begin{aligned}
\mathcal{H}^{2 k+1}(\partial E) & \gtrsim \sum_{C \in \tau^{2 k+1}} \mathscr{H}^{2 k+1}\left(\partial E \cap B_{2 \Lambda_{k} r}\left(h_{C}\right)\right) \\
& \gtrsim \sum_{C \in \tau^{2 k+1}} \mathscr{H}^{2 k+2}((E \triangle F) \cap C)=\mathscr{H}^{2 k+2}(E \triangle F) .
\end{aligned}
$$

For future use, we shall record the following lemma, which is a consequence of Lemma 21.

Lemma 22. Suppose that $E \subseteq \mathbb{H}^{2 k+1}$ is a measurable set such that $\mathcal{H}^{2 k+1}\left(B_{r} \cap \partial E\right)<\infty$ for all $r \in(0, \infty)$. Then for any open set $U \subseteq \mathbb{H}^{2 k+1}$, we have $\operatorname{Per}_{E}(U) \lesssim \mathscr{H}^{2 k+1}(U \cap \partial E)$.

Proof. First, consider the case $\mathscr{H}^{2 k+1}(\partial E)<\infty$. By Lemma 21, there is a sequence $\left\{F_{2^{-i}}\right\}_{i=1}^{\infty}$ of scalings of cellular sets such that $\lim _{i \rightarrow \infty} \mathbf{1}_{F_{2-i}}=$ $\mathbf{1}_{E}$ in $L_{1}^{\text {loc }}\left(\mathscr{H}^{2 k+2}\right)$ and $\mathscr{H}^{2 k+1}\left(\partial F_{i}\right) \lesssim \mathscr{H}^{2 k+1}(\partial E)$ for every $i \in \mathbb{N}$. Since by Proposition 16 we have $\operatorname{Per}_{F_{2-i}}\left(\mathbb{H}^{2 k+1}\right) \asymp \mathscr{H}^{2 k+1}\left(\partial F_{2^{-i}}\right)$, it follows that $\operatorname{Per}_{F_{2-i}}\left(\mathbb{H}^{2 k+1}\right) \lesssim \mathscr{H}^{2 k+1}(\partial E)$. So, by Proposition 19 we have $\operatorname{Per}_{E}\left(\mathbb{H}^{2 k+1}\right) \lesssim$
$\mathscr{H}^{2 k+1}(\partial E)<\infty$. By Proposition 18 we thus have $\operatorname{Per}_{E}(U) \lesssim \mathscr{H}^{2 k+1}(U \cap \partial E)$ for every open set $U \subseteq \mathbb{H}^{2 k+1}$.

For the general case, for every $i \in \mathbb{N}$, let $E_{i}=B_{i} \cap E$. So $\lim _{i \rightarrow \infty} \mathbf{1}_{E_{i}}=\mathbf{1}_{E}$ in $L_{1}^{\text {loc }}\left(\mathscr{H}^{2 k+2}\right)$. By the validity of Lemma 22 in the special case $\mathscr{H}^{2 k+1}(\partial E)$ $<\infty$ that we just established, Proposition 19 implies $E$ has locally finite perimeter. So, by Proposition 18 we have $\operatorname{Per}_{E}(U) \lesssim \mathscr{H}^{2 k+1}(U \cap \partial E)$ for every open set $U \subseteq \mathbb{H}^{2 k+1}$.
2.2. Vertical perimeter. The vertical perimeter of a subset of the Heisenberg group was introduced in [66]. If $E \subseteq \mathbb{H}^{2 k+1}$ and $s \in \mathbb{R}$, then for every $\rho \in \mathbb{R}$, denote

$$
\mathrm{D}_{s} E \stackrel{\text { def }}{=} E \triangle E Z^{2^{2 s}} \subseteq \mathbb{H}^{2 k+1}
$$

If $U, E \subseteq \mathbb{H}^{2 k+1}$ are measurable, then define $\overline{\mathrm{v}}_{U}(E): \mathbb{R} \rightarrow \mathbb{R}$ by setting for all $\rho \in \mathbb{R}$,

$$
\begin{align*}
\overline{\mathrm{V}}_{U}(E)(\rho) & \stackrel{\text { def }}{=} \frac{\mathscr{H}^{2 k+2}\left(U \cap \mathrm{D}_{\rho} E\right)}{2^{\rho}} \\
& =\frac{1}{2^{\rho}} \int_{U}\left|\mathbf{1}_{E}(x)-\mathbf{1}_{E}\left(x Z^{-2^{2 \rho}}\right)\right| \mathrm{d} \mathscr{H}^{2 k+2}(x) . \tag{28}
\end{align*}
$$

Note that $\overline{\mathrm{v}}_{U}(E)$ is left-invariant in the sense that $\overline{\mathrm{v}}_{g} U(g E)=\overline{\mathrm{v}}_{U}(E)$ for every $g \in \mathbb{H}^{2 k+1}$. When $U=\mathbb{H}^{2 k+1}$ it will be convenient to use the simpler notation $\overline{\mathrm{v}}_{\mathbb{H}^{2 k+1}}(E)=\overline{\mathrm{v}}(E)$. The vertical perimeter of $E$ is then defined to be the quantity

$$
\|\overline{\mathrm{v}}(E)\|_{L_{2}(\mathbb{R})}=\left(\int_{-\infty}^{\infty} \overline{\mathrm{v}}(E)(\rho)^{2} \mathrm{~d} \rho\right)^{\frac{1}{2}}
$$

Remark 23. The following inequality and example may shed some light on the behavior of the vertical perimeter. If $E \subseteq \mathbb{H}^{2 k+1}$ is measurable and $s \in \mathbb{R}$, then there is a geodesic $\gamma_{s}$ from $\mathbf{0}$ to $Z^{2^{2 s}}$ with $\ell\left(\gamma_{s}\right) \asymp 2^{s}$. We can connect each point $h$ of $E$ to a point of $E Z^{2^{2 s}}$ by a translate of $\gamma_{s}$, and if $h Z^{2^{2 s}} \in E \triangle E Z^{2^{2 s}}$, then $h \gamma_{s}$ crosses $\partial E$. It follows from these observations that

$$
\begin{equation*}
\overline{\mathrm{v}}(E)(s) \lesssim \frac{\mathcal{H}^{2 k+1}(\partial E) d\left(\mathbf{0}, Z^{2^{2 s}}\right)}{2^{s}} \lesssim \mathscr{H}^{2 k+1}(\partial E) \tag{29}
\end{equation*}
$$

Thus, $\|\overline{\mathrm{v}}(E)(s)\|_{L_{\infty}(\mathbb{R})} \lesssim \mathscr{H}^{2 k+1}(\partial E)$. For an example where this bound is sharp, consider the box

$$
\begin{aligned}
C_{r} & \stackrel{\text { def }}{=}[-r, r]^{2 k} \times\left[-r^{2}, r^{2}\right] \\
& =\left\{h: \in \mathbb{H}^{2 k+1}: \max \left\{\left|x_{1}(h)\right|,\left|y_{1}(h)\right|, \ldots,\left|x_{k}(h)\right|,\left|y_{k}(h)\right|, \sqrt{|z(h)|}\right\} \leqslant r\right\}
\end{aligned}
$$

for some $r \in(0, \infty)$. A straightforward computation shows that

$$
\forall s \in \mathbb{R}, \quad \overline{\mathrm{v}}\left(C_{r}\right)(s) \asymp \begin{cases}2^{s} r^{2 k} & \text { if } 2^{s}<r, \\ \frac{r^{2 k+2}}{2^{s}} & \text { if } 2^{s} \geqslant r .\end{cases}
$$

So, $\overline{\mathrm{v}}\left(C_{r}\right)(s)$ has a maximum of roughly $r^{2 k+1} \asymp \mathscr{H}^{2 k+1}\left(\partial C_{r}\right)$ near $s=\log _{2} r$ and exponentially decaying tails. In general, $\overline{\mathrm{v}}(E)(s)$ measures the "bumpiness" of $\partial E$ at scale $2^{s}$, so the vertical perimeter of $E$ can include contributions from many different scales.

For ease of future use, we record some elementary properties of the vertical perimeter function.

Lemma 24. Suppose that $A, B, U \subseteq \mathbb{H}^{2 k+1}$ are measurable sets of finite measure. Then
(1) $\overline{\mathrm{v}}(A)=\overline{\mathrm{v}}\left(A^{\mathrm{c}}\right)$,
(2) $|\overline{\mathrm{v}}(A)-\overline{\mathrm{v}}(B)| \leqslant \overline{\mathrm{v}}(A \triangle B)$,
(3) $\overline{\mathrm{v}}(A \cap B) \leqslant \overline{\mathrm{v}}_{B}(A)+\overline{\mathrm{v}}(B)$,
(4) $\overline{\mathrm{v}}\left(\mathfrak{s}_{t}(A)\right)(\rho)=t^{2 k+1} \overline{\mathrm{v}}(A)\left(\rho-\log _{2} t\right)$ for every $t \in(0, \infty)$ and every $\rho \in \mathbb{R}$.

Proof. The first property follows from the fact that $\mathrm{D}_{\rho} A=\mathrm{D}_{\rho} A^{\mathrm{c}}$. For the second property, note that since the Boolean operation $\triangle$ is associative and commutative, we have

$$
\begin{aligned}
2^{\rho}|\overline{\mathrm{v}}(A)(\rho)-\overline{\mathrm{v}}(B)(\rho)| & \stackrel{(28)}{=}\left|\left\|\mathbf{1}_{\mathrm{D}_{\rho} A}\right\|_{L_{1}\left(\mathcal{H}^{2 k+2}\right)}-\left\|\mathbf{1}_{\mathrm{D}_{\rho} B}\right\|_{L_{1}\left(\mathcal{H}^{2 k+2}\right)}\right| \\
& \leqslant \mathscr{H}^{2 k+2}\left(\mathrm{D}_{\rho} A \triangle \mathrm{D}_{\rho} B\right) \\
& =\mathscr{H}^{2 k+2}\left((A \triangle B) \triangle\left(A Z^{2^{2 \rho}} \triangle B Z^{2^{2 \rho}}\right)\right) \\
& =\mathscr{H}^{2 k+2}\left(\mathrm{D}_{\rho}(A \triangle B)\right)=2^{\rho} \overline{\mathrm{v}}(A \triangle B)(\rho) .
\end{aligned}
$$

For the third property, it suffices to show that $\mathrm{D}_{\rho}(A \cap B) \subseteq\left(B \cap \mathrm{D}_{\rho} A\right) \cup \mathrm{D}_{\rho}(B)$. To this end, fix $h \in \mathrm{D}_{\rho}(A \cap B)$. If $\left|\left\{h Z^{-2^{2 \rho}}, h\right\} \cap B\right|=1$, then $h \in \mathrm{D}_{\rho}(B)$, as required. Otherwise either $\left\{h Z^{-2^{2 \rho}}, h\right\} \subseteq B$ or $\left\{h Z^{-2^{2 \rho}}, h\right\} \cap B=\varnothing$. The latter case cannot occur because it implies that $h Z^{-2^{2 \rho}}, h \notin A \cap B$, and therefore $h \notin \mathrm{D}_{\rho}(A \cap B)$ in contradiction to our choice of $h$. It remains to consider the case $\left\{h Z^{-2^{2 \rho}}, h\right\} \subseteq B$, but then the assumption $h \in \mathrm{D}_{\rho}(A \cap B)$ implies that also $h \in \mathrm{D}_{\rho} A$, as required. Finally, the fourth property holds since for every $(t, \rho) \in(0, \infty) \times \mathbb{R}$, we have

$$
\begin{aligned}
\overline{\mathrm{v}}\left(\mathfrak{s}_{t}(A)\right)(\rho) & =\frac{\mathcal{H}^{2 k+2}\left(\mathfrak{s}_{t}(A) \triangle\left(\mathfrak{s}_{t}(A) Z^{2^{2 \rho}}\right)\right)}{2^{\rho}} \\
& =\frac{\mathscr{H}^{2 k+2}\left(\mathfrak{s}_{t}(A) \triangle \mathfrak{s}_{t}\left(A Z^{\frac{1}{t^{2} 2^{2}}}\right)\right)}{2^{\rho}}=\frac{\mathscr{H}^{2 k+2}\left(\mathfrak{s}_{t}\left(\mathrm{D}_{\rho-\log _{2} t}(A)\right)\right)}{2^{\rho}} \\
& =\frac{t^{2 k+2} \mathscr{H}^{2 k+2}\left(\mathrm{D}_{\rho-\log _{2} t} A\right)}{2^{\rho}}=t^{2 k+1} \overline{\mathrm{v}}(A)\left(\rho-\log _{2} t\right) .
\end{aligned}
$$

2.3. Intrinsic Lipschitz graphs. One of the key ideas in the present work is that surfaces in the Heisenberg group can be approximated by intrinsic Lipschitz graphs. These were introduced by Franchi, Serapioni, and Serra Cassano [40] as a natural class of surfaces in the Heisenberg group that are analogues of Lipschitz graphs in Euclidean space. David and Semmes proved [30] that a set in Euclidean space is uniformly rectifiable if it has a certain decomposition (a corona decomposition; see Section 6 below) into pieces that can be approximated by Lipschitz graphs. In Section 6 below we will define a similar intrinsic version of this decomposition for surfaces in $\mathbb{H}^{2 k+1}$.

We recall the definitions of intrinsic graphs and intrinsic Lipschitz graphs, which are due to [40]. Let $V \subseteq \mathbb{H}^{2 k+1}$ be any vertical plane; typically, we take $\left\{h \in \mathbb{H}^{2 k+1}: x_{k}(h)=0\right\}$. Let $W=V^{\perp}$ be the orthogonal complement of $V$, i.e., the horizontal line through the origin that is perpendicular to $V$. Suppose that $U \subseteq V$ and that $f: U \rightarrow W$ is any function. Then the intrinsic graph of $f$ is defined to be the following set:

$$
\Gamma_{f} \stackrel{\text { def }}{=}\{u f(u): u \in U\} .
$$

Below, any such set is called an intrinsic graph over $U$. It is often convenient to discuss intrinsic graphs of real-valued functions, which is achieved by identifying $W$ with $\mathbb{R}$. Specifically, if $w \in W$ is a vector of unit Euclidean length and $f: U \rightarrow \mathbb{R}$, then we slightly abuse notation by denoting

$$
\begin{equation*}
\Gamma_{f} \stackrel{\text { def }}{=}\left\{u w^{f(u)}: u \in U\right\} . \tag{30}
\end{equation*}
$$

Since the unitary group $U(k)$ acts transitively on the unit sphere of $\mathbb{R}^{2 k}$, if $\Gamma$ is any intrinsic graph, then there exists an isometric automorphism of $\mathbb{H}^{2 k+1}$ that takes $\Gamma$ to an intrinsic graph over $U \subseteq\left\{h \in \mathbb{H}^{2 k+1}: x_{k}(h)=0\right\}$ of the form $\Gamma=\left\{v X_{k}^{f(v)}: v \in U\right\}$. If, in addition, we have $U=V=\left\{h \in \mathbb{H}^{2 k+1}\right.$ : $\left.x_{k}(h)=0\right\}$, then $\Gamma$ bounds two half-spaces $\left\{v X_{k}^{t}: v \in V \wedge t \geqslant f(v)\right\}$ and $\left\{v X_{k}^{t}: v \in V \wedge t \leqslant f(v)\right\}$, which we refer to below as $\Gamma^{+}$and $\Gamma^{-}$. (The superscripts + and - will be used interchangeably and are not intended to indicate the sign of the inequality.)

If $V \subseteq \mathbb{H}^{2 k+1}$ is a vertical plane and $W=V^{\perp}$, then denote the orthogonal projection onto $W$ by $\operatorname{Proj}_{W}: \mathbb{R}^{2 k+1} \rightarrow W$. For every $\lambda \in(0,1)$ define the following double cone:

$$
\operatorname{Cone}_{\lambda}(V) \stackrel{\text { def }}{=}\left\{h \in \mathbb{H}^{2 k+1}:\left\|\operatorname{Proj}_{W}(h)\right\|>\lambda d(\mathbf{0}, h)\right\} .
$$

When $V=\left\{h \in \mathbb{H}^{2 k+1}: x_{k}(h)=0\right\}$, we shall use the simpler notation

$$
\begin{aligned}
\operatorname{Cone}_{\lambda} & =\operatorname{Cone}_{\lambda}\left(\left\{h \in \mathbb{H}^{2 k+1}: x_{k}(h)=0\right\}\right) \\
& =\left\{h \in \mathbb{H}^{2 k+1}:\left|x_{k}(h)\right|>\lambda d(\mathbf{0}, h)\right\} .
\end{aligned}
$$

If $U \subseteq V$ and $\Gamma \subseteq \mathbb{H}^{2 k+1}$ is an intrinsic graph over $U$, then we say that $\Gamma$ is an intrinsic Lipschitz graph over $U$ if there exists $\lambda \in(0,1)$ such that for every $h \in \Gamma$, we have $\left(h \operatorname{Cone}_{\lambda}(V)\right) \cap \Gamma=\varnothing$. In this case we say that $\Gamma$ is an intrinsic $\lambda$-Lipschitz graph. Correspondingly, $\Gamma^{+}$and $\Gamma^{-}$are then called intrinsic $\lambda$-Lipschitz half-spaces.

We warn that this definition is different from the definition of the Lipschitz constant of an intrinsic Lipschitz graph that is given in [41], but we believe that the above modification is more natural in the present context. One reason for this is that the above definition leads to Theorem 27 below, which is a sharp Lipschitz extension statement that is an intrinsic version of the classical nonlinear Hahn-Banach theorem [83]. (See also [110] or [13, Lemma 1.1].)

In the concrete setting $V=\left\{h \in \mathbb{H}^{2 k+1}: x_{k}(h)=0\right\}$, given $\lambda \in(0,1)$ and an intrinsic graph $\Gamma$, the requirement that $\left(h\right.$ Cone $\left._{\lambda}\right) \cap \Gamma=\varnothing$ for all $h \in \Gamma$ is the same as the requirement that for every $w_{1}, w_{2} \in \Gamma$, we have $\left|x_{k}\left(w_{1}\right)-x_{k}\left(w_{2}\right)\right| \leqslant \lambda d\left(w_{1}, w_{2}\right)$.

LEMMA 25. If $h \in \mathbb{H}^{2 k+1}$ and $\Gamma$ is an intrinsic $\lambda$-Lipschitz graph over $U \subseteq V$, then $h \Gamma$ is an intrinsic $\lambda$-Lipschitz graph over some set $U^{\prime} \subseteq V$.

Proof. Let $f: U \rightarrow \mathbb{R}$ be such that $\Gamma=\left\{v X_{k}^{f(v)}: v \in U\right\}$, and let $v_{0} \in V$, $t_{0} \in \mathbb{R}$ be the unique elements such that $h=v_{0} X_{k}^{t_{0}}$. Then

$$
\begin{aligned}
v c h \Gamma & =\left\{h v X_{k}^{f(v)}: v \in U\right\} \\
& =\left\{v_{0} X_{k}^{t_{0}} v X_{k}^{f(v)}: v \in U\right\}=\left\{v_{0} v Z^{t_{0} y_{k}(v)} X_{k}^{f(v)+t_{0}}: v \in U\right\}
\end{aligned}
$$

If $v^{\prime}=v_{0} v Z^{t_{0} y_{k}(v)}$, then $v=v_{0}^{-1} v^{\prime} Z^{-t_{0} y_{k}\left(v_{0}^{-1} v^{\prime}\right)}$. So, by setting

$$
U^{\prime} \stackrel{\text { def }}{=}\left\{v_{0} v Z^{t_{0} y_{k}(v)}: v \in U\right\}
$$

and defining $f^{\prime}: U^{\prime} \rightarrow \mathbb{R}$ by

$$
\forall v^{\prime} \in U^{\prime}, \quad f^{\prime}\left(v^{\prime}\right) \stackrel{\text { def }}{=} f\left(v_{0}^{-1} v^{\prime} Z^{-t_{0} y_{k}\left(v_{0}^{-1} v^{\prime}\right)}\right)+t_{0}
$$

we conclude that $h \Gamma$ is the intrinsic graph of $f^{\prime}$ over $U^{\prime}$. Since the metric on $\mathbb{H}^{2 k+1}$ is invariant under translation, $h \Gamma$ is also $\lambda$-Lipschitz.

Remark 26. If $\lambda \in(0,1)$ and $\Gamma=\Gamma_{f}$ is an intrinsic graph of a function $f: V \rightarrow W$ as above, then there is no relationship between the requirement that $\Gamma$ is an intrinsic Lipschitz graph and the requirement that $f$ is a Lipschitz function between the corresponding subsets of $\mathbb{H}^{2 k+1}$ (equipped with the Carnot-Carathéodory metric); see Remark 3.13 in [41].

A variant of the following extension result is Theorem 4.25 of [41], with the difference being that in [41] this is obtained with a weaker bound on the Lipschitz constant of the extended graph. Control on the Lipschitz constant of
the extension will be important for us in Section 9.6 below, so we will include a proof of Theorem 27 here, while also taking the opportunity to state it in a sharp form. The important property for us will turn out to be the weaker assertion that if the Lipschitz constant of $\Gamma$ is small, then the Lipschitz constant of the extended graph $\widetilde{\Gamma}$ is also small.

Theorem 27. Fix $\lambda \in(0,1)$ and a vertical plane $V \subseteq \mathbb{H}^{2 k+1}$. Write $W=V^{\perp}$. Suppose that $\varnothing \neq U \subseteq V$ and that $f: U \rightarrow W$ is a function such that $\Gamma=\Gamma_{f}$ is an intrinsic $\lambda$-Lipschitz graph over $U$. Then there exists an intrinsic $\lambda$-Lipschitz graph over $V$, denoted $\widetilde{\Gamma}$, such that $\Gamma \subseteq \widetilde{\Gamma}$.

Proof. We may assume without loss of generality that $V=\left\{h \in \mathbb{H}^{2 k+1}\right.$ : $\left.x_{k}(h)=0\right\}$ and $W=\left\langle X_{k}\right\rangle$. Due to Lemma 25, by applying a translation we may suppose that $\mathbf{0} \in \Gamma$.

The double cone Cone ${ }_{\lambda}$ consists of two halves, which we denote by Cone $_{\lambda}^{+}$ and Cone ${ }_{\lambda}^{-}$. Namely,

$$
\text { Cone }_{\lambda}^{+} \stackrel{\text { def }}{=}\left\{h: \in \mathbb{H}^{2 k+1}: x_{1}(h)>\lambda d(\mathbf{0}, h)\right\}
$$

and

$$
\text { Cone }_{\lambda}^{-} \stackrel{\text { def }}{=}\left\{h: \in \mathbb{H}^{2 k+1}: x_{1}(h)<-\lambda d(\mathbf{0}, h)\right\} .
$$

Thus Cone ${ }_{\lambda}^{-}=\left(\text {Cone }_{\lambda}^{+}\right)^{-1}$, and for all $g, h \in \mathbb{H}^{2 k+1}$, we have $g \in h$ Cone $_{\lambda}^{+}$if and only if $h \in g$ Cone $_{\lambda}^{-}$. Define

$$
\widetilde{\Gamma}^{+} \stackrel{\text { def }}{=} \bigcup_{h \in \Gamma} h \text { Cone }_{\lambda}^{+} .
$$

We claim that $\widetilde{\Gamma} \stackrel{\text { def }}{=} \partial \widetilde{\Gamma}^{+}$satisfies the desired conditions.
Observe first that we have

$$
\begin{equation*}
\text { Cone }_{\lambda}^{+} \cdot \text { Cone }_{\lambda}^{+} \subseteq \text { Cone }_{\lambda}^{+} \text {. } \tag{31}
\end{equation*}
$$

Indeed, if $h_{1}, h_{2} \in \operatorname{Cone}_{\lambda}^{+}$, then $x_{1}\left(h_{1}\right)>\lambda d\left(\mathbf{0}, h_{1}\right)$ and $x_{1}\left(h_{2}\right)>\lambda d\left(\mathbf{0}, h_{2}\right)$. Thus,

$$
\begin{align*}
x_{1}\left(h_{1} h_{2}\right) & =x_{1}\left(h_{1}\right)+x_{1}\left(h_{2}\right)>\lambda\left(d\left(\mathbf{0}, h_{1}\right)+d\left(\mathbf{0}, h_{2}\right)\right)  \tag{32}\\
& =\lambda\left(d\left(\mathbf{0}, h_{1}^{-1}\right)+d\left(\mathbf{0}, h_{2}\right)\right) \geqslant \lambda d\left(h_{1}^{-1}, h_{2}\right)=\lambda d\left(\mathbf{0}, h_{1} h_{2}\right) .
\end{align*}
$$

By the definition of $\mathrm{Cone}_{\lambda}^{+}$, (32) means that $h_{1} h_{2} \in \mathrm{Cone}_{\lambda}^{+}$, thus completing the verification of (31). A key consequence is that

$$
\begin{equation*}
\widetilde{\Gamma}^{+} . \text {Cone }_{\lambda}^{+} \subseteq \widetilde{\Gamma}^{+} . \tag{33}
\end{equation*}
$$

We claim that $\widetilde{\Gamma}$ is an intrinsic graph, that is, that for any $v \in V$,

$$
\Phi_{v} \stackrel{\text { def }}{=}\left\{t \in \mathbb{R}: v X_{k}^{t} \in \widetilde{\Gamma}^{+}\right\} \subseteq \mathbb{R}
$$

is a half-line. First, we show that $\Phi_{v} \neq \varnothing$. Since $\mathbf{0} \in \Gamma$, it suffices to show that $v X_{k}^{t} \in$ Cone $_{\lambda}^{+}$when $t$ is sufficiently large. For every $t \in \mathbb{R}$, we have

$$
\begin{equation*}
d\left(v X_{k}^{t}, \mathbf{0}\right) \leqslant d(\mathbf{0}, v)+d\left(v, v X_{k}^{t}\right) \stackrel{(25)}{=} d(\mathbf{0}, v)+|t| \tag{34}
\end{equation*}
$$

Since $x_{k}\left(v X_{k}^{t}\right)=t$, it follows from (34) that

$$
t>\frac{\lambda}{1-\lambda} d(\mathbf{0}, v) \Longrightarrow x_{k}\left(v X_{k}^{t}\right)>\lambda d\left(v X_{k}^{t}, \mathbf{0}\right)
$$

Thus there exists $t>0$ for which $v X_{k}^{t} \in \operatorname{Cone}_{\lambda}^{+} \subseteq \widetilde{\Gamma}^{+}$. Hence $\Phi_{v} \neq \varnothing$. Similarly,

$$
\begin{equation*}
\exists t \in(-\infty, 0), \quad v X_{k}^{t} \in \text { Cone }_{\lambda}^{-} \tag{35}
\end{equation*}
$$

Next, we claim that $\Phi_{v} \neq \mathbb{R}$. In fact, we show that for all $h \in \Gamma$,

$$
\begin{equation*}
h \text { Cone }_{\lambda}^{-} \cap \widetilde{\Gamma}^{+}=\varnothing . \tag{36}
\end{equation*}
$$

By (35), this implies the claim.
Suppose that $w \in h$ Cone $_{\lambda}^{-} \cap \widetilde{\Gamma}^{+} ;$then there is an $h^{\prime} \in \Gamma$ such that $w \in$ $h^{\prime}$ Cone $_{\lambda}^{+}$. We have $h^{-1} w \in$ Cone $_{\lambda}^{-}$, so $w^{-1} h \in$ Cone $_{\lambda}^{+}$and therefore

$$
h=w \cdot w^{-1} h \in h^{\prime} \text { Cone }_{\lambda}^{+} \cdot \text { Cone }_{\lambda}^{+} \subseteq h^{\prime} \text { Cone }_{\lambda}^{+} .
$$

Since $h, h^{\prime} \in \Gamma$, this contradicts the underlying assumption (since $f$ is an intrinsic $\lambda$-Lipschitz function) that $\left(h^{\prime}\right.$ Cone $\left._{\lambda}\right) \cap \Gamma=\varnothing$.

Since $X_{\underset{k}{s}}^{s} \in$ Cone $_{\lambda}^{+}$for $s>0$, by (33) we see that if $t \in \Phi_{v}$, then $v X_{k}^{t+s}=$ $\left(v X_{k}^{t}\right) X_{k}^{s} \in \widetilde{\Gamma}^{+}$. Cone $_{\lambda}^{+} \subseteq \widetilde{\Gamma}^{+}$, i.e., $t+s \in \Phi_{v}$. Thus $t \in \Phi_{v} \Longrightarrow[t, \infty) \subseteq \Phi_{v}$. Since $\Phi_{v}$ is an open nonempty proper subset of $\mathbb{R}$, this shows that $\Phi_{v}=$ $(\phi(v), \infty)$ for some $\phi(v) \in \mathbb{R}$. Consequently, $\widetilde{\Gamma}=\Gamma_{\phi}$, i.e., $\widetilde{\Gamma}$ is an intrinsic graph over $V$.

If $h=u X_{k}^{f(u)} \in \Gamma$, then $u X_{k}^{f(u)+s} \in h$ Cone $_{\lambda}^{+}$for all $s>0$ and $u X_{k}^{f(u)+s} \in$ $h$ Cone $_{\lambda}^{-}$for all $s<0$. Equation (36) then implies that $h \in \widetilde{\Gamma}$, so $\widetilde{\Gamma} \supseteq \Gamma$.

It remains to prove that $\widetilde{\Gamma}$ is an intrinsic $\lambda$-Lipschitz graph. To this end, the goal is to show that if $h_{1}, h_{2} \in \widetilde{\Gamma}$, then $h_{2} \notin h_{1}$ Cone $_{\lambda}$. Because $h_{1} \in \widetilde{\Gamma}$, for every $\varepsilon>0$, there is a $p \in \widetilde{\Gamma}^{+}$such that $d\left(p, h_{1}\right)<\varepsilon$. Let $g \in \Gamma$ be such that $p \in g$ Cone $_{\lambda}^{+}$. Using (31) we therefore have

$$
\begin{equation*}
p \text { Cone }_{\lambda}^{+} \subseteq g \text { Cone }_{\lambda}^{+} \cdot \text { Cone }_{\lambda}^{+} \subseteq g \text { Cone }_{\lambda}^{+} \subseteq \widetilde{\Gamma}^{+} \tag{37}
\end{equation*}
$$

Since this holds for every $\varepsilon>0$, we have $h_{1} \operatorname{Cone}_{\lambda}^{+} \subseteq \widetilde{\Gamma}^{+}$. Since $\widetilde{\Gamma}^{+}$is an open subset of $\mathbb{H}^{2 k+1}$ and $h_{2}$ is assumed to be in its boundary $\widetilde{\Gamma}$, it follows from (37) that $h_{2} \notin h_{1}$ Cone $_{\lambda}^{+}$. By the same reasoning with the roles of $h_{1}$ and $h_{2}$ reversed, we also see that

$$
h_{1} \notin h_{2} \text { Cone }_{\lambda}^{+}=h_{2}\left(\text { Cone }_{\lambda}^{-}\right)^{-1} \Longrightarrow h_{2} \notin h_{1} \text { Cone }_{\lambda}^{-}
$$

Since $h_{1}$ Cone $_{\lambda}=\left(h_{1}\right.$ Cone $\left._{\lambda}^{+}\right) \cup\left(h_{1}\right.$ Cone $\left._{\lambda}^{-}\right)$, this concludes the proof that $h_{2} \notin$ $h_{1}$ Cone $_{\lambda}$.

## 3. Initial reductions

The purpose of this section is to present simple reductions between various questions so as to set the stage for the main steps of the proof of Theorem 3. In Section 3.1 below we shall relate various inequalities for functions on the discrete Heisenberg group. Some of these reductions were already quoted in Section 1.3 to deduce Theorem 9 (and hence also all of the new results that were stated in the introduction) from Theorem 3. The arguments in Section 3.1 are for the most part due to [66], except that in [66] they were carried out only in the special case $k=1$. This was done because when [66] was written the relevance of the assumption $k \geqslant 2$ in Theorem 3 was not known, and it was therefore believed that all of the Heisenberg groups $\left\{\mathbb{H}_{\mathbb{Z}}^{2 k+1}\right\}_{k=1}^{\infty}$ have identical roles in the present context. Since it turns out that the underlying dimension does play a role, for completeness we include in Section 3.1 an explanation of the straightforward modifications of the arguments in [66] so as to obtain the desired statements for general $k \in \mathbb{N}$. Section 3.2 shows via a partition of unity argument that in order to prove Theorem 3 it suffices to prove a certain (singular) Sobolev-type inequality on the continuous Heisenberg group. Again, the argument is included for completeness, but it follows steps that were carried out in [66] for the special case $k=1$. Section 3.3 is devoted to reductions that take place entirely in the continuous setting. We first show that in order to establish the Sobolev-type inequality of Section 3.2 it suffices to prove an isoperimetric-type inequality on the continuous Heisenberg group that is analogous to its discrete counterpart that we stated in Theorem 3; this step is a quick application of the coarea formula. We conclude Section 3.3 by establishing the technical (but convenient) statement that it suffices to prove the desired continuous isoperimetric-type inequality for cellular sets.
3.1. Reductions in the discrete setting. We shall start by presenting simple arguments that allow one to relate various inequalities on the discrete Heisenberg group. For ease of reference in future work, we shall present our statements in a form that is more general than what is needed here (but the proofs of the more general case are identical to the special case that we will use).
3.1.1. From sets to functions. The following lemma explains how certain isoperimetric-type inequalities, as in Theorem 3 or in (5), are equivalent to certain functional inequalities. This step is a standard and simple application of the classical co-area formula, in combination with convexity.

Lemma 28. Fix $k \in \mathbb{N}$, a sequence $\left\{w_{t}\right\}_{t=1}^{\infty} \subseteq[0, \infty), C \in(0, \infty)$ and $q \in[1, \infty)$. The inequality

$$
\begin{equation*}
\left(\sum_{t=1}^{\infty} w_{t}\left|\partial_{\mathrm{v}}^{t} \Omega\right|^{q}\right)^{\frac{1}{q}} \leqslant C\left|\partial_{\mathrm{h}} \Omega\right| \tag{38}
\end{equation*}
$$

holds true for every finite $\Omega \subseteq \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ if and only if every finitely supported $\phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow \mathbb{R}$ satisfies

$$
\begin{equation*}
\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left|\phi\left(h Z^{t}\right)-\phi(h)\right|\right)^{q}\right)^{\frac{1}{q}} \leqslant \frac{C}{2} \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}|\phi(h \sigma)-\phi(h)| \tag{39}
\end{equation*}
$$

Proof. The estimate (38) is nothing more than the special case $\phi=\mathbf{1}_{\Omega}$ of (39). Conversely, to show that (39) follows from (38), suppose that $\phi$ : $\mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow \mathbb{R}$ is finitely supported, and for every $u \in \mathbb{R}$, denote $\Omega_{u}=\{h \in$ $\left.\mathbb{H}_{\mathbb{Z}}^{2 k+1}: \phi(h)<u\right\}$. Since $\phi$ is finitely supported, $\Omega_{u}$ is finite if $u \leqslant 0$ and $\mathbb{H}_{\mathbb{Z}}^{2 k+1} \backslash \Omega_{u}$ is finite when $u>0$. In both cases (38) holds true with $\Omega$ replaced by $\Omega_{u}$, and therefore (39) holds true with $\phi$ replaced by $\mathbf{1}_{\Omega_{u}}$ for every $u \in \mathbb{R}$, i.e.,

$$
\begin{align*}
\left(\sum _ { t = 1 } ^ { \infty } w _ { t } \left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \mid\right.\right. & \left.\left.\mathbf{1}_{\left\{\phi\left(h Z^{t}\right)<u\right\}}-\mathbf{1}_{\{\phi(h)<u\}} \mid\right)^{q}\right)^{\frac{1}{q}}  \tag{40}\\
& \leqslant \frac{C}{2} \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\left|\mathbf{1}_{\{\phi(h \sigma)<u\}}-\mathbf{1}_{\{\phi(h)<u\}}\right|
\end{align*}
$$

The right-hand side of (39) is equal to the integral of the right-hand side of (40) with respect to $u$. Hence, it suffices to use the triangle inequality in $\ell_{q}$ to conclude as follows:

$$
\begin{aligned}
\frac{C}{2} \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} & \sum_{\sigma \in \mathfrak{S}_{k}}|\phi(h \sigma)-\phi(h)| \\
& \stackrel{(40)}{\geqslant} \int_{-\infty}^{\infty}\left(w_{t}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left|\mathbf{1}_{\left\{\phi\left(h Z^{t}\right)<u\right\}}-\mathbf{1}_{\{\phi(h)<u\}}\right|\right)^{q}\right)^{\frac{1}{q}} \mathrm{~d} u \\
& \geqslant\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \int_{-\infty}^{\infty}\left|\mathbf{1}_{\left\{\phi\left(h Z^{t}\right)<u\right\}}-\mathbf{1}_{\{\phi(h)<u\}}\right| \mathrm{d} u\right)^{q}\right)^{\frac{1}{q}} \\
& =\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left|\phi\left(h Z^{t}\right)-\phi(h)\right|\right)^{q}\right)^{\frac{1}{q}}
\end{aligned}
$$

3.1.2. From real-valued to vector-valued. The following lemma is another simple use of convexity that shows that certain real-valued functional inequalities are equivalent to the analogous inequalities for functions that take value in an $L_{p}(\mu)$ space.

Lemma 29. Fix $k \in \mathbb{N}$, a sequence $\left\{w_{t}\right\}_{t=1}^{\infty} \subseteq[0, \infty)$, and $p, q \in(0, \infty)$ with $p \leqslant q$. Suppose also that $\eta \in(0, \infty)$ is such that for every finitely supported
$\phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow \mathbb{R}$ we have
$\eta\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left|\phi\left(h Z^{t}\right)-\phi(h)\right|^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}} \leqslant\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}|\phi(h \sigma)-\phi(h)|^{p}\right)^{\frac{1}{p}}$.
Then also for every measure space $(\mathcal{S}, \mu)$ and every finitely supported $\Phi$ : $\mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow L_{p}(\mu)$, we have

$$
\begin{gather*}
\eta\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left\|\Phi\left(h Z^{t}\right)-\Phi(h)\right\|_{L_{p}(\mu)}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}} \\
\quad \leqslant\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\Phi(h \sigma)-\Phi(h)\|_{L_{p}(\mu)}^{p}\right)^{\frac{1}{p}} . \tag{42}
\end{gather*}
$$

Proof. For $\mu$-almost every $s \in \mathcal{S}$, the function $\phi_{s}: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow \mathbb{R}$ given by $\phi_{s}(h)=\Phi(h)(s)$ is finitely supported, and therefore by (41) we have

$$
\begin{aligned}
& \sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\Phi(h \sigma)-\Phi(h)\|_{L_{p}(\mu)}^{p} \\
& \quad=\int_{\mathcal{S}} \sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\left|\phi_{s}(h \sigma)-\phi_{s}(h)\right|^{p} \mathrm{~d} \mu(s) \\
& \quad \stackrel{(41)}{\geqslant} \eta \int_{\mathcal{S}}\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}}\left|\phi_{s}\left(h Z^{t}\right)-\phi_{s}(h)\right|^{p}\right)^{\frac{q}{p}}\right)^{\frac{p}{q}} \mathrm{~d} \mu(s) \\
& \quad \geqslant \eta\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \int_{S}\left|\phi_{s}\left(h Z^{t}\right)-\phi_{s}(h)\right|^{p} \mathrm{~d} \mu(s)\right)^{\frac{q}{p}}\right)^{\frac{p}{q}} \\
&=\eta\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}}\left\|\Phi\left(h Z^{t}\right)-\Phi(h)\right\|_{L_{p}(\mu)}^{p}\right)^{\frac{q}{p}}\right)^{\frac{p}{q}},
\end{aligned}
$$

where in the penultimate step we used the triangle inequality in $\ell_{q / p}$; recall that $q \geqslant p$.
3.1.3. From global to local. The proof of the following lemma follows the steps of a proof that appears in Section 3.2 of [66]. We include the details because [66] makes an analogous statement only when $X$ is uniformly convex, $k=1$ and $p>1$, while here we need to use the case $X=L_{1}(\mathbb{R}), k \geqslant 2$ and $p=1$. The argument below is a straightforward adaptation of the argument in [66].

Lemma 30. Fix $k \in \mathbb{N}$, a sequence $\left\{w_{t}\right\}_{t=1}^{\infty} \subseteq[0, \infty)$, $p, q \in[1, \infty)$ and $K \in(0, \infty)$. Let $\left(X,\|\cdot\|_{X}\right)$ be a normed space such that every finitely supported $\phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow X$ satisfies

$$
\begin{align*}
& \left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}}\left\|\phi\left(h Z^{t}\right)-\phi(h)\right\|_{X}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}}  \tag{43}\\
& \quad \leqslant K\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\phi(h \sigma)-\phi(h)\|_{X}^{p}\right)^{\frac{1}{p}}
\end{align*}
$$

Then there is a universal constant $c \in \mathbb{N}(c=21$ works here $)$ such that for every $f: \mathbb{H}_{Z}^{2 k+1} \rightarrow X$ and every $n \in \mathbb{N}$, we have

$$
\begin{equation*}
\left(\sum_{t=1}^{n^{2}} w_{t}\left(\sum_{h \in \mathcal{B}_{n}}\left\|f\left(h Z^{t}\right)-f(h)\right\|_{X}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}} \lesssim_{k} K\left(\sum_{h \in \mathcal{B}_{c n}} \sum_{\sigma \in \mathfrak{S}_{k}}\|f(h \sigma)-f(h)\|_{X}^{p}\right)^{\frac{1}{p}} \tag{44}
\end{equation*}
$$

Before proving Lemma 30, for ease of later reference we record the following general estimates.

Lemma 31. Let $\Gamma$ be a finitely generated group, and fix a finite symmetric generating set $\Sigma \subseteq \Gamma$. Let $\rho_{\Sigma}: \Gamma \times \Gamma \rightarrow[0, \infty)$ denote the left-invariant word metric that is induced by $\Sigma$ on $\Gamma$ and for every $n \in \mathbb{N}$, denote by $B_{\Sigma}(n)=$ $\left\{\gamma \in \Gamma: \rho_{\Sigma}\left(g, 1_{\Gamma}\right) \leqslant n\right\}$ the $\rho_{\Sigma}$-ball of radius $n$ centered at the identity element $1_{\Gamma} \in \Gamma$. Then, for every $p \in[1, \infty)$, every $n \in \mathbb{N}$, every metric space $\left(m, d_{m}\right)$, and every finitely supported function $\phi: \Gamma \rightarrow m$, we have

$$
\begin{align*}
& \left(\sum_{x \in \Gamma} \frac{1}{\left|B_{\Sigma}(n)\right|} \sum_{y \in B_{\Sigma}(n)} d m(\phi(x y), \phi(x))^{p}\right)^{\frac{1}{p}} \\
& \quad \leqslant n\left(\sum_{x \in \Gamma} \max _{\sigma \in \Sigma} d_{m}(\phi(x \sigma), \phi(x))^{p}\right)^{\frac{1}{p}} . \tag{45}
\end{align*}
$$

Also, if $\Gamma$ has polynomial growth, then, fixing $r \in \mathbb{N}$ and $\alpha, \beta>0$ such that $\frac{1}{\alpha} m^{r} \leqslant\left|B_{\Sigma}(m)\right| \leqslant \beta m^{r}$ for every $m \in \mathbb{N}$ (such $\alpha, \beta, r$ exist by Gromov's theorem [44]), we have

$$
\begin{align*}
& \left(\frac{1}{\left|B_{\Sigma}(n)\right|^{2}} \sum_{x, y \in B_{\Sigma}(n)} d m(\phi(x), \phi(y))^{p}\right)^{\frac{1}{p}}  \tag{46}\\
& \quad \leqslant 2 n\left(\frac{6^{r} \alpha^{2} \beta^{2}}{\left|B_{\Sigma}(3 n)\right|} \max _{\sigma \in \Sigma} \sum_{x \in B_{\Sigma}(3 n)} d m(\phi(x \sigma), \phi(x))^{p}\right)^{\frac{1}{p}} .
\end{align*}
$$

Proof. The proof of (45) is essentially contained in the proof of Lemma 3.4 in [66], even though [66, Lemma 3.2] is stated for the special case $\Gamma=\mathbb{H}_{\mathbb{Z}}^{3}$. We
briefly recall the simple argument. For every $y \in B_{\Sigma}(n)$, fix $\sigma_{1}(y), \ldots, \sigma_{n}(y) \in$ $\Sigma \cup\left\{1_{\Gamma}\right\}$ such that $y=\sigma_{1}(y) \cdots \sigma_{n}(y)$. Then, denoting $\gamma_{j}(y)=\sigma_{1}(y) \cdots \sigma_{j-1}(y)$ for every $j \in\{1, \ldots, n\}$, by the triangle inequality we have
$\forall(x, y) \in \Gamma \times B_{\Sigma}(n), \quad d_{m}(\phi(x y), \phi(x)) \leqslant \sum_{j=1}^{n} d_{m}\left(\phi\left(x \gamma_{j}(y) \sigma_{j}(y)\right), \phi\left(x \gamma_{j}(y)\right)\right)$.
By the triangle inequality in $\ell_{p}(\Gamma)$, this implies that

$$
\begin{aligned}
\left(\sum_{x \in \Gamma} \sum_{y \in B_{\Sigma}(n)}\right. & \left.d m(\phi(x y), \phi(x))^{p}\right)^{\frac{1}{p}} \\
& \leqslant \sum_{j=1}^{n}\left(\sum_{y \in B_{\Sigma}(n)} \sum_{x \in \Gamma} d m\left(\phi\left(x \gamma_{j}(y) \sigma_{j}(y)\right), \phi\left(x \gamma_{j}(y)\right)\right)^{p}\right)^{\frac{1}{p}} \\
& =\sum_{j=1}^{n}\left(\sum_{y \in B_{\Sigma}(n)} \sum_{z \in \Gamma} d m\left(\phi\left(z \sigma_{j}(y)\right), \phi(z)\right)^{p}\right)^{\frac{1}{p}} \\
& \leqslant n\left(\left|B_{\Sigma}(n)\right| \sum_{z \in \Gamma} \max _{\sigma \in \Sigma} d m(\phi(z \sigma), \phi(z))^{p}\right)^{\frac{1}{p}}
\end{aligned}
$$

thus concluding the proof of (45).
The proof of (46) is essentially contained in the proof of Lemma 3.2 of [66], which is itself a generalization of an inequality of Kleiner [62, Th. 2.2] that treats real-valued mappings and $p=2$. Indeed, even though [66, Lemma 3.2] is stated for the special case $\Gamma=\mathbb{H}_{\mathbb{Z}}^{3}$, up to the final sentence of the proof of [66, Lemma 3.2] the argument is valid in any finitely generated group $\Gamma$, and the penultimate step of the final displayed equation in the proof of $[66$, Lemma 3.2] gives that

$$
\begin{aligned}
& \frac{1}{\left|B_{\Sigma}(n)\right|^{2}} \sum_{x, y \in B_{\Sigma}(n)} d m(\phi(x), \phi(y))^{p} \\
& \leqslant(2 n)^{p} \cdot \frac{\left|B_{\Sigma}(2 n)\right| \cdot\left|B_{\Sigma}(3 n)\right|}{\left|B_{\Sigma}(n)\right|^{2}} \cdot \max _{\sigma \in \Sigma} \frac{1}{\left|B_{\Sigma}(3 n)\right|} \sum_{x \in B_{\Sigma}(3 n)} d m(\phi(x), \phi(x \sigma))^{p} .
\end{aligned}
$$

Proof of Lemma 30. Fix $n \in \mathbb{N}$. By translating $f$ we may assume without loss of generality that $\sum_{x \in \mathcal{B}_{7 n}} f(x)=0$. Let $\xi_{n}: \mathbb{H}_{Z}^{2 k+1} \rightarrow[0,1]$ be a cutoff function that is $\frac{1}{n}$-Lipschitz (with respect to the word metric $d_{W}$ ), is equal to 1 on $\mathcal{B}_{5 n}$, and vanishes outside $\mathcal{B}_{6 n}$. Consider the finitely supported mapping $\phi_{n}=\xi_{n} f: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow X$. We claim that the following two estimates hold true:

$$
\begin{align*}
& \left(\sum_{t=1}^{n^{2}} w_{t}\left(\sum_{h \in \mathcal{B}_{n}}\left\|f\left(h Z^{t}\right)-f(h)\right\|_{X}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}} \\
& \quad \leqslant\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}}\left\|\phi_{n}\left(h Z^{t}\right)-\phi_{n}(h)\right\|_{X}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}} \tag{47}
\end{align*}
$$

and

$$
\begin{equation*}
\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\left\|\phi_{n}(h \sigma)-\phi_{n}(h)\right\|_{X}^{p}\right)^{\frac{1}{p}} \lesssim_{k}\left(\sum_{h \in \mathcal{B}_{21 n}} \sum_{\sigma \in \mathfrak{S}_{k}}\|f(h \sigma)-f(h)\|_{X}^{p}\right)^{\frac{1}{p}} \tag{48}
\end{equation*}
$$

The estimates (47) and (48) imply the desired inequality (44) by an application of (43) to $\phi_{n}$.

Suppose that $h \in \mathcal{B}_{n}$ and $t \in\left\{1, \ldots, n^{2}\right\}$. Since $d_{W}\left(\mathbf{0}, Z^{t}\right) \leqslant 4 n$ (see, e.g., [14]), we have $h, h Z^{t} \in \mathcal{B}_{5 n}$, and therefore $\phi(h)=f(h)$ and $\phi\left(h Z^{t}\right)=f\left(h Z^{t}\right)$. This immediately implies (47). Next, since $\xi_{n}$ is $\frac{1}{n}$-Lipschitz and takes values in $[0,1]$, for every $\sigma \in \mathfrak{S}_{k}$ and $h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ we have

$$
\begin{aligned}
& \left\|\phi_{n}(h \sigma)-\phi_{n}(h)\right\|_{X}=\left\|\xi_{n}(h \sigma)(f(h \sigma)-f(h))+\left(\xi_{n}(h \sigma)-\xi_{n}(h)\right) f(h)\right\|_{X} \\
& \quad \leqslant \xi_{n}(h \sigma)\|f(h \sigma)-f(h)\|_{X}+\left|\xi_{n}(h \sigma)-\xi_{n}(h)\right| \cdot\|f(h)\|_{X} \\
& \quad \leqslant\|f(h \sigma)-f(h)\|_{X}+\frac{1}{n}\|f(h)\|_{X}
\end{aligned}
$$

Therefore, by the triangle inequality in $\ell_{p}\left(\mathcal{B}_{7 n} \times \mathfrak{S}_{k}\right)$ we have

$$
\begin{align*}
& \left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\left\|\phi_{n}(h \sigma)-\phi_{n}(h)\right\|_{X}^{p}\right)^{\frac{1}{p}} \\
& \quad=\left(\sum_{h \in \mathcal{B}_{7 n}} \sum_{\sigma \in \mathfrak{S}_{k}}\left\|\phi_{n}(h \sigma)-\phi_{n}(h)\right\|_{X}^{p}\right)^{\frac{1}{p}}  \tag{49}\\
& \quad \leqslant\left(\sum_{h \in \mathcal{B}_{7 n}} \sum_{\sigma \in \mathfrak{S}_{k}}\|f(h \sigma)-f(h)\|_{X}^{p}\right)^{\frac{1}{p}}+\frac{(2 k)^{\frac{1}{p}}}{n}\left(\sum_{h \in \mathcal{B}_{7 n}}\|f(h)\|_{X}^{p}\right)^{\frac{1}{p}}
\end{align*}
$$

where the first step of (49) holds true because $\phi_{n}$ is supported on $\mathcal{B}_{6 n}$. To deduce (48) from (49), and hence also to conclude the proof of Lemma 30, it suffices to note that

$$
\begin{aligned}
\frac{1}{n}\left(\sum_{h \in \mathcal{B}_{7 n}}\|f(h)\|_{X}^{p}\right)^{\frac{1}{p}} & \leqslant \frac{1}{n}\left(\frac{1}{\left|\mathcal{B}_{7 n}\right|} \sum_{g, h \in \mathcal{B}_{7 n}}\|f(g)-f(h)\|_{X}^{p}\right)^{\frac{1}{p}} \\
& \lesssim k\left(\sum_{h \in \mathcal{B}_{21 n}} \sum_{\sigma \in \mathfrak{S}_{k}}\|f(h \sigma)-f(h)\|_{X}^{p}\right)^{\frac{1}{p}}
\end{aligned}
$$

where the first step follows from Jensen's inequality since $\sum_{x \in \mathcal{B}_{7 n}} f(x)=0$, and the final step is an application of (46), using the fact [11] that $\left|\mathcal{B}_{m}\right| \asymp_{k} m^{2 k+2}$ for every $m \in \mathbb{N}$.
3.1.4. From low dimension to high dimension. Lemma 32 below is of lesser importance for the present purposes. We shall use it only to show that the case $k=2$ of Theorem 3 implies the general case $k \in\{2,3,4, \ldots$,$\} with the$ stated dependence on $k$ (otherwise the ensuing proof yields a worse asymptotic dependence on $k$ as $k \rightarrow \infty$ ). Below, a mapping $\phi: X \rightarrow Y$ between two abstract sets $X$ and $Y$ is said to be finitely supported if there exists a point $y \in Y$ such that $\phi(x)=y$ for all but finitely many $x \in X$, i.e., $\left|\phi^{-1}(Y \backslash\{y\})\right|<$ $\infty$. This is a slight (and harmless) departure from the more common use of this term when $Y$ is a Banach space, in which case one requires that $y=0$.

Lemma 32. Fix $k, m \in \mathbb{N}$ with $m \leqslant k$, a sequence $\left\{w_{t}\right\}_{t=1}^{\infty} \subseteq[0, \infty)$, $q \in[1, \infty)$ and $\beta \in(0, \infty)$. Let $X$ be a set and $\mathrm{K}, \mathrm{L}: X \times X \rightarrow[0, \infty)$ be two symmetric functions such that $\mathrm{K}(x, x)=\mathrm{L}(x, x)=0$ for all $x \in X$. Suppose that for every finitely supported mapping $\phi: \mathbb{H}_{\mathbb{Z}}^{2 m+1} \rightarrow X$ we have

$$
\begin{equation*}
\beta\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 m+1}} \mathrm{~K}\left(\phi\left(h Z^{t}\right), \phi(h)\right)\right)^{q}\right)^{\frac{1}{q}} \leqslant \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 m+1}} \sum_{\sigma \in \mathfrak{S}_{m}} \mathrm{~L}(\phi(h \sigma), \phi(h)) . \tag{50}
\end{equation*}
$$

Then for every finitely supported mapping $\Phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow X$ we have

$$
\begin{equation*}
\frac{\beta k}{m}\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \mathrm{~K}\left(\Phi\left(h Z^{t}\right), \Phi(h)\right)\right)^{q}\right)^{\frac{1}{q}} \leqslant \sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}} \mathrm{~L}(\Phi(h \sigma), \Phi(h)) \tag{51}
\end{equation*}
$$

Proof. For every $A \subseteq\{1, \ldots, k\}$ denote $\mathfrak{S}_{A}=\left\{X_{i}^{ \pm 1}, Y_{i}^{ \pm 1}\right\}_{i \in A} \subseteq \mathfrak{S}_{k}$. Let $\mathbb{G}_{A} \triangleleft \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ be the (normal) subgroup of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ that is generated by $\mathfrak{S}_{A}$ and let $\left\{g_{j}(A)\right\}_{j=1}^{\infty} \subseteq \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ be representatives of the distinct left-cosets of $\mathbb{G}_{A}$ in $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$. If $|A|=m$, then $\mathbb{G}_{A}$ is isomorphic to $\mathbb{H}_{\mathbb{Z}}^{2 m+1}$. Therefore, for every $A \subseteq\{1, \ldots, k\}$ with $|A|=m$ and every $j \in \mathbb{N}$, an application of (50) to the mapping $\left(h \in \mathbb{G}_{A}\right) \mapsto \Phi\left(g_{j}(A) h\right)$ shows that

$$
\begin{gather*}
\beta\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{G}_{A}} \mathrm{~K}\left(\Phi\left(g_{j}(A) h Z^{t}\right), \Phi\left(g_{j}(A) h\right)\right)\right)^{q}\right)^{\frac{1}{q}}  \tag{52}\\
\leqslant \sum_{h \in \mathbb{G}_{A}} \sum_{\sigma \in \mathfrak{G}_{A}} \mathrm{~L}\left(\Phi\left(g_{j}(A) h \sigma\right), \Phi\left(g_{j}(A) h\right)\right) .
\end{gather*}
$$

Observe that

$$
\begin{equation*}
\sum_{j=1}^{\infty} \sum_{h \in \mathfrak{G}_{A}} \sum_{\sigma \in \mathfrak{S}_{A}} \mathrm{~L}\left(\Phi\left(g_{j}(A) h \sigma\right), \Phi\left(g_{j}(A) h\right)\right)=\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{A}} \mathrm{~L}(\Phi(h \sigma), \Phi(h)), \tag{53}
\end{equation*}
$$

and, by the triangle inequality in $\ell_{q}$,

$$
\begin{align*}
\sum_{j=1}^{\infty} & \left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{G}_{A}} \mathrm{~K}\left(\Phi\left(g_{j}(A) h Z^{t}\right), \Phi\left(g_{j}(A) h\right)\right)\right)^{q}\right)^{\frac{1}{q}} \\
& \geqslant\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{j=1}^{\infty} \sum_{h \in \mathbb{G}_{A}} \mathrm{~K}\left(\Phi\left(g_{j}(A) h Z^{t}\right), \Phi\left(g_{j}(A) h\right)\right)\right)^{q}\right)^{\frac{1}{q}}  \tag{54}\\
& =\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \mathrm{~K}\left(\Phi\left(h Z^{t}\right), \Phi(h)\right)\right)^{q}\right)^{\frac{1}{q}}
\end{align*}
$$

By summing (52) over $j \in \mathbb{N}$ and using (53) and (54), it follows that

$$
\begin{equation*}
\beta\left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \mathrm{~K}\left(\Phi\left(h Z^{t}\right), \Phi(h)\right)\right)^{q}\right)^{\frac{1}{q}} \leqslant \sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{A}} \mathrm{~L}(\Phi(h \sigma), \Phi(h)) . \tag{55}
\end{equation*}
$$

By summing (55) over all those $A \subseteq\{1, \ldots, k\}$ with $|A|=m$, we obtain the estimate

$$
\begin{align*}
\beta\binom{k}{m}\left(\sum_{t=1}^{\infty} w_{t}( \right. & \left.\left.\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \mathrm{~K}\left(\Phi\left(h Z^{t}\right), \Phi(h)\right)\right)^{q}\right)^{\frac{1}{q}}  \tag{56}\\
& \leqslant\binom{ k-1}{m-1} \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}} \mathrm{~L}(\Phi(h \sigma), \Phi(h)),
\end{align*}
$$

where we used the fact that for each $\sigma \in \mathfrak{S}_{k}$, the number of $A \subseteq\{1, \ldots, k\}$ with $|A|=m$ that contain $\sigma$ is equal to $\binom{k-1}{m-1}$. The desired estimate (51) is the same as (56).

Suppose that we have already proved Theorem 3 when $k=2$, i.e., that $\left|\partial_{\mathrm{v}} \Omega\right| \lesssim\left|\partial_{\mathrm{h}} \Omega\right|$ for every finite $\Omega \subseteq \mathbb{H}_{\mathbb{Z}}^{5}$. By Lemma 28 this implies that for every finitely supported $\phi: \mathbb{H}_{\mathbb{Z}}^{5} \rightarrow \mathbb{R}$,

$$
\left(\sum_{t=1}^{\infty} \frac{1}{t^{2}}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{5}}\left|\phi\left(h Z^{t}\right)-\phi(h)\right|\right)^{2}\right)^{\frac{1}{2}} \lesssim \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{5}} \sum_{\sigma \in \mathfrak{G}_{2}}|\phi(h \sigma)-\phi(h)| .
$$

Due to Lemma 32, it follows that for every integer $k \geqslant 2$ and every finitely supported $\Phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow \mathbb{R}$,

$$
\left(\sum_{t=1}^{\infty} \frac{1}{t^{2}}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}}\left|\Phi\left(h Z^{t}\right)-\Phi(h)\right|\right)^{2}\right)^{\frac{1}{2}} \lesssim \frac{1}{k} \sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}|\Phi(h \sigma)-\Phi(h)| .
$$

By (the trivial direction of) Lemma 28, it follows that $\left|\partial_{\mathrm{V}} \Omega\right| \lesssim \frac{1}{k}\left|\partial_{\mathrm{h}} \Omega\right|$ for every finite $\Omega \subseteq \mathbb{H}_{\mathbb{Z}}^{2 k+1}$; i.e., the conclusion of Theorem 3 holds true with the stated asymptotic dependence on $k$ as $k \rightarrow \infty$. This shows that in order to prove Theorem 3 we can from now on ignore the dependence on $k$.
3.2. From continuous to discrete. Our goal here is to show that Theorem 3 follows from a certain (singular) Sobolev-type inequality on the continuous Heisenberg group that is stated in Corollary 34 below. Once this assertion is established, we will henceforth focus our attention entirely on considerations that take place in the continuous setting. The proof of the following lemma is a simple adaptation of the steps of a proof that appears in Section 3.3 of [66].

Lemma 33. Fix $k \in \mathbb{N}$ and $\mathcal{W}:[1, \infty) \rightarrow[0, \infty)$ that satisfies $\|\mathcal{W}\|_{L_{1}(1, \infty)}$ $=\int_{1}^{\infty} \mathcal{W}(s) \mathrm{d} s<\infty$. Define a sequence $\left\{w_{t}\right\}_{t=1}^{\infty}$ by setting

$$
\forall t \in \mathbb{N}, \quad w_{t} \stackrel{\text { def }}{=} \int_{s}^{s+1} \mathcal{W}(s) \mathrm{d} s
$$

Fix also $C \in(0, \infty)$ and $p, q \in(0, \infty)$ with $p \geqslant 1$. Let $\left(X,\|\cdot\|_{X}\right)$ be a normed space such that every smooth and compactly supported function $\Phi: \mathbb{H}_{\mathbb{R}}^{2 k+1} \rightarrow X$ satisfies

$$
\begin{align*}
\left(\int _ { 1 } ^ { \infty } \mathcal { W } ( t ) \left(\int_{\mathbb{H}^{2 k+1}}\right.\right. & \left.\left.\left\|\Phi\left(h Z^{t}\right)-\Phi(h)\right\|_{X}^{p} \mathrm{~d} \mathscr{H}^{2 k+2}(h)\right)^{\frac{q}{p}} \mathrm{~d} t\right)^{\frac{1}{q}}  \tag{57}\\
& \leqslant C\left(\int_{\mathbb{H}^{2 k+1}}\left\|\nabla_{\mathbb{H}} \Phi(h)\right\|_{\ell_{p}^{2 k}(X)}^{p} \mathrm{~d} \mathscr{H}^{2 k+2}(h)\right)^{\frac{1}{p}} .
\end{align*}
$$

Then every finitely supported function $\phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow X$ satisfies

$$
\begin{align*}
\left(\sum _ { t = 1 } ^ { \infty } w _ { t } \left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\right.\right. & \left.\left.\left\|\phi\left(h Z^{t}\right)-\phi(h)\right\|_{X}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}}  \tag{58}\\
& \lesssim\left(C+\|\mathcal{W}\|_{L_{1}(1, \infty)}^{\frac{1}{q}}\right)\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\phi(h \sigma)-\phi(h)\|_{X}^{p}\right)^{\frac{1}{p}}
\end{align*}
$$

Proof of Lemma 33. Suppose that $\phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow X$ is finitely supported. We shall show that there exists a compactly supported smooth function $\Phi$ : $\mathbb{H}^{2 k+1} \rightarrow X$ that satisfies the following two estimates, which immediately imply the desired inequality (58) via an application of (57):

$$
\begin{equation*}
\left(\int_{\mathbb{H}^{2 k+1}}\left\|\nabla_{\mathbb{H}} \Phi(h)\right\|_{\ell_{p}^{2 k}(X)}^{p} \mathrm{~d} \mathscr{\not}^{2 k+2}(h)\right)^{\frac{1}{p}} \lesssim\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\phi(h \sigma)-\phi(h)\|_{X}^{p}\right)^{\frac{1}{p}} \tag{59}
\end{equation*}
$$

and

$$
\begin{align*}
& \left(\sum_{t=1}^{\infty} w_{t}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}}\left\|\phi\left(h Z^{t}\right)-\phi(h)\right\|_{X}^{p}\right)^{\frac{q}{p}}\right)^{\frac{1}{q}} \\
& \quad \lesssim\left(\int_{1}^{\infty} \mathcal{W}(t)\left(\int_{\mathbb{H}^{2 k+1}}\left\|\Phi\left(h Z^{t}\right)-\Phi(h)\right\|_{X}^{p} \mathrm{~d} \mathscr{H}^{2 k+2}(h)\right)^{\frac{q}{p}} \mathrm{~d} t\right)^{\frac{1}{q}}  \tag{60}\\
& \quad \quad+\|\mathcal{W}\|_{L_{1}(1, \infty)}^{\frac{1}{q}}\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\phi(h \sigma)-\phi(h)\|_{X}^{p}\right)^{\frac{1}{p}} .
\end{align*}
$$

Since $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ is a co-compact lattice of $\mathbb{H}^{2 k+1}$, from now on we can fix a compactly supported smooth function $\chi: \mathbb{H}^{2 k+1} \rightarrow[0,1]$ such that

$$
\sum_{g \in \mathbb{H}_{Z}^{2 k+1}} \chi\left(g^{-1} h\right)=1
$$

for all $h \in \mathbb{H}^{2 k+1}$. Define

$$
\forall h \in \mathbb{H}^{2 k+1}, \quad \Phi(h) \stackrel{\text { def }}{=} \sum_{g \in \mathbb{H}_{Z}^{2 k+1}} \chi\left(g^{-1} h\right) f(g) .
$$

If we denote the support of $\phi$ by $S=\operatorname{supp}(\phi) \subseteq \mathbb{H}_{\mathbb{Z}}^{2 k+1} \subseteq \mathbb{H}^{2 k+1}$ and the support of $\chi$ by $T=\operatorname{supp}(\chi) \subseteq \mathbb{H}^{2 k+1}$, then the support of $\Phi$ is contained in $S T$. Hence $\Phi$ is compactly supported.

Fix $m \in \mathbb{N}$ such that $T \cap \mathbb{H}_{Z}^{2 k+1} \subseteq \mathcal{B}_{m}$. Note that $m$ depends only on $k$ (through the choice of the bump function $\chi$, which is fixed once and for all). Arguing identically to the proof of equation (56) of [66] (which uses the triangle inequality in $\ell_{p}^{2 k}(X)$, i.e., the assumption $p \geqslant 1$ ), we see that

$$
\begin{equation*}
\forall g \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}, \quad \sup _{h \in g T}\left\|\nabla_{\mathbb{H}} \Phi(h)\right\|_{\ell_{p}^{2 k}}^{p} \lesssim \sum_{z \in \mathcal{B}_{2 m}}\|\phi(g z)-\phi(g)\|_{X}^{p} . \tag{61}
\end{equation*}
$$

Hence, since the assumption on $\chi$ implies $\bigcup_{g \in \mathbb{H}_{Z}^{2 k+1}} g T=\mathbb{H}^{2 k+1}$, we deduce that

$$
\begin{aligned}
\left(\int_{\mathbb{H}^{2 k+1}}\right. & \left.\left\|\nabla_{\mathbb{H}} \Phi(h)\right\|_{\ell_{p}^{2 k}(X)}^{p} \mathrm{~d} \mathscr{H}^{2 k+2}(h)\right)^{\frac{1}{p}} \\
& \leqslant\left(\sum_{g \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \int_{g T}\left\|\nabla_{\mathbb{H}} \Phi(h)\right\|_{\ell_{p}^{2 k}(X)}^{p} \mathrm{~d} \mathscr{H}^{2 k+2}(h)\right)^{\frac{1}{p}} \\
& \stackrel{(61)}{\lesssim}\left(\sum_{g \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{z \in \mathcal{B}_{2 m}}\|\phi(g z)-\phi(g)\|_{X}^{p}\right)^{\frac{1}{p}} \\
& \stackrel{(45)}{\lesssim}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\phi(h \sigma)-\phi(h)\|_{X}^{p}\right)^{\frac{1}{p}}
\end{aligned}
$$

This completes the justification of (59).

To prove (60), by arguing identically to the proof of equation (64) in [66], while using (45) in place of the use of [66, Lemma 3.4], we see that for every $t \in \mathbb{N}$ and $s \in[t, t+1]$, we have

$$
\begin{align*}
\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left\|\phi\left(h Z^{t}\right)-\phi(h)\right\|_{X}^{p}\right)^{\frac{1}{p}} \lesssim & \left(\int_{\mathbb{H}^{2 k+1}}\left\|\Phi\left(h Z^{s}\right)-\Phi(h)\right\|_{X}^{p} \mathrm{~d} \mathscr{H}^{2 k+2}(h)\right)^{\frac{1}{p}}  \tag{62}\\
& +\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{k}}\|\phi(h \sigma)-\phi(h)\|_{X}^{p}\right)^{\frac{1}{p}}
\end{align*}
$$

The desired estimate (60) now follows by raising (62) to the power $q$, multiplying the resulting estimate by $\mathcal{W}(s)$, integrating over $s \in[t, t+1]$ and summing over $t \in \mathbb{N}$.

The following corollary is the special case $p=1, q=2, X=\mathbb{R}$ and $\mathcal{W}(s)=1 / s^{2}$ of Lemma 33 (combined with Lemma 32 so as to obtain the stated dependence on $k$ ).

Corollary 34. In order to establish Theorem 3 it suffices to prove that for any $k \in\{2,3, \ldots$,$\} , every finitely supported smooth function f: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ satisfies

$$
\begin{aligned}
&\left(\int_{0}^{\infty}\left(\int_{\mathbb{H}^{2 k+1}}\left|f\left(h Z^{t}\right)-f(h)\right| \mathrm{d} \mathscr{H}^{2 k+2}(h)\right)^{2} \frac{\mathrm{~d} t}{t^{2}}\right)^{\frac{1}{2}} \\
& \lesssim \int_{\mathbb{H}^{2 k+1}}\left\|\nabla_{\mathbb{H}} f(h)\right\|_{\ell_{1}^{2 k}} \mathrm{~d} \mathscr{H}^{2 k+2}(h) .
\end{aligned}
$$

3.3. Reductions in the continuous setting. By Corollary 34, all of the new results that we stated in the introduction follow from Theorem 35 below, which answers Question 4.1 in [66] positively when the dimension of the underlying Heisenberg group is at least 5. Note, however, that [66, Question 4.1] ignored the role of this dimension, so the forthcoming work that we discussed in Remark 4 shows that [66, Question 4.1] actually has a negative answer when $k=1$.

Theorem 35. Fix $k \in\{2,3, \ldots$,$\} . If f: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ is smooth and compactly supported, then

$$
\begin{gather*}
\left(\int_{0}^{\infty}\left(\int_{\mathbb{H}^{2 k+1}}\left|f\left(h Z^{t}\right)-f(h)\right| \mathrm{d} \mathscr{H}^{2 k+2}(h)\right)^{2} \frac{\mathrm{~d} t}{t^{2}}\right)^{\frac{1}{2}}  \tag{63}\\
\lesssim \int_{\mathbb{H}^{2 k+1}}\left\|\nabla_{\mathbb{H}} f(h)\right\|_{\ell_{1}^{2 k}} \mathrm{~d} \mathscr{H}^{2 k+2}(h) .
\end{gather*}
$$

Theorem 35 follows from an isoperimetric-type inequality that is stated in Theorem 36 below. It relates the vertical perimeter (recall Section 2.2) of
a set to the measure of its boundary. Suppose that $E \subseteq \mathbb{H}^{2 k+1}$ is measurable and $\mathscr{H}^{2 k+1}(\partial E)<\infty$. It is not clear a priori that $\|\overline{\mathrm{v}}(E)\|_{L_{2}(\mathbb{R})}$ is necessarily finite, but the following theorem shows that it is at most a constant multiple of $\mathscr{H}^{2 k+1}(\partial E)$ when the dimension of the underlying Heisenberg group is at least 5 .

Theorem 36. If $k \in\{2,3, \ldots$,$\} and E \subseteq \mathbb{H}^{2 k+1}$ is measurable, then $\|\overline{\mathrm{v}}(E)\|_{L_{2}(\mathbb{R})} \lesssim \mathscr{H}^{2 k+1}(\partial E)$.

We shall now explain how Theorem 36 implies Theorem 35. To do so, we will need the following well-known coarea formula; see, e.g., [37, Th. 2.3.5] or [1, Th. 3.3].

Theorem 37. Let $f: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ be a smooth, compactly supported function, and for $u \in \mathbb{R}$, let $E_{u}=\left\{h \in \mathbb{H}^{2 k+1}: f(h)<u\right\}$. Then the following inequality holds true:

$$
\begin{equation*}
\int_{-\infty}^{\infty} \mathscr{H}^{2 k+1}\left(\partial E_{u}\right) \mathrm{d} u \lesssim \int_{\mathbb{H}^{2 k+1}}\left\|\nabla_{\mathbb{H}} f(h)\right\|_{\ell_{1}^{2 k}} \mathrm{~d} \mathscr{H}^{2 k+2}(h) . \tag{64}
\end{equation*}
$$

When $k \geqslant 2$ and $q=2$, the following lemma shows that Theorem 36 implies Theorem 35.

Lemma 38. Fix $k \in \mathbb{N}, q \in[1, \infty)$ and $C \in(0, \infty)$. Suppose that $\|\overline{\mathrm{v}}(E)\|_{L_{q}(\mathbb{R})} \leqslant C \not^{2 k+1}(\partial E)$ for every measurable $E \subseteq \mathbb{H}^{2 k+1}$. Then every smooth and compactly supported $f: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ satisfies

$$
\begin{aligned}
&\left(\int_{0}^{\infty}\left(\int_{\mathbb{H}^{2 k+1}}\left|f\left(h Z^{t}\right)-f(h)\right| \mathrm{d} \mathscr{H}^{2 k+2}(h)\right)^{q} \frac{\mathrm{~d} t}{t^{1+\frac{q}{2}}}\right)^{\frac{1}{q}} \\
& \lesssim C \int_{\mathbb{H}^{2 k+1}}\left\|\nabla_{\mathbb{H}} f(h)\right\|_{\ell_{1}^{2 k}} \mathrm{~d} \mathscr{H}^{2 k+2}(h) .
\end{aligned}
$$

Proof. If $f: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ is measurable, then

$$
\int_{\mathbb{H}^{2 k+1}}\left|f\left(h Z^{t}\right)-f(h)\right| \mathrm{d} \mathscr{H}^{2 k+2}(h)=\int_{-\infty}^{\infty} \mathscr{H}^{2 k+2}\left(E_{u} \triangle E_{u} Z^{t}\right) \mathrm{d} u .
$$

Hence, by the triangle inequality in $L_{q}(\mathbb{R})$ and the substitution $t=2^{2 s}$, we have

$$
\begin{align*}
& \left(\int_{0}^{\infty}\left(\int_{\mathbb{H}^{2 k+1}}\left|f\left(h Z^{t}\right)-f(h)\right| \mathrm{d} \mathscr{H}^{2 k+2}(h)\right)^{q} \frac{\mathrm{~d} t}{t^{1+\frac{q}{2}}}\right)^{\frac{1}{q}}  \tag{65}\\
\leqslant & \int_{-\infty}^{\infty}\left(\int_{0}^{\infty} \mathscr{H}^{2 k+2}\left(E_{u} \triangle E_{u} Z^{t}\right)^{q} \frac{\mathrm{~d} t}{t^{1+\frac{q}{2}}}\right)^{\frac{1}{q}} \\
\asymp & \int_{-\infty}^{\infty}\left(\int_{-\infty}^{\infty} \frac{\mathcal{H}^{2 k+2}\left(E_{u} \triangle E_{u} Z^{2^{2 s}}\right)^{q}}{2^{q s}} \mathrm{~d} s\right)^{\frac{1}{q}} \mathrm{~d} u \stackrel{(28)}{=} \int_{-\infty}^{\infty}\left\|\overline{\mathrm{v}}\left(E_{u}\right)\right\|_{L_{q}(\mathbb{R})} \mathrm{d} u .
\end{align*}
$$

Now, by applying the assumption of Lemma 38 to each $\left\{E_{u}\right\}_{u \in \mathbb{R}}$, we conclude the proof of Lemma 38 by combining (65) with the following estimate:

$$
\begin{aligned}
\int_{-\infty}^{\infty}\left\|\overline{\mathrm{v}}\left(E_{u}\right)\right\|_{L_{q}(\mathbb{R})} \mathrm{d} u & \lesssim C \int_{-\infty}^{\infty} \mathcal{H}^{2 k+1}\left(\partial E_{u}\right) \mathrm{d} u \\
& \stackrel{(64)}{\lesssim} C \int_{\mathbb{H}^{2 k+1}}\left\|\nabla_{\mathbb{H}} f(h)\right\|_{\ell_{1}^{2 k}} \mathrm{~d} \mathscr{H}^{2 k+2}(h) .
\end{aligned}
$$

Theorem 39 below is nothing more than Theorem 36 in the special case of cellular sets; recall Section 2.1.6. We prefer to state Theorem 39 separately because we shall next prove that it implies Theorem 36 in full generality. Once this is done, we can devote the rest of the discussion to the proof of Theorem 39.

Theorem 39. If $k \in\{2,3, \ldots\}$ and $F \subseteq \mathbb{H}^{2 k+1}$ is a cellular set, then $\|\overline{\mathrm{v}}(F)\|_{L_{2}(\mathbb{R})} \lesssim \mathscr{H}^{2 k+1}(\partial F)$.

The special case $k \in\{2,3, \ldots\}$ and $q=2$ of the following lemma shows that Theorem 36 in full generality is a consequence of (its special case) Theorem 39.

Lemma 40. Fix $k \in \mathbb{N}, q \in[1, \infty)$, and $C \in(0, \infty)$. Suppose that $\|\overline{\mathrm{v}}(F)\|_{L_{q}(\mathbb{R})} \leqslant C \mathscr{H}^{2 k+1}(\partial F)$ for any cellular set $F \subseteq \mathbb{H}^{2 k+1}$. Then $\|\overline{\mathrm{v}}(E)\|_{L_{q}(\mathbb{R})}$ $\lesssim(C+1) \mathcal{H}^{2 k+1}(\partial E)$ for any measurable $E \subseteq \mathbb{H}^{2 k+1}$.

Proof. Suppose that $E \subseteq \mathbb{H}^{2 k+1}$ is a measurable set satisfying $\mathscr{H}^{2 k+1}(\partial E)$ $<\infty$. Fix $i \in \mathbb{Z}$, and denote $\rho=2^{i}$. By Lemma 21, there is a set $F_{\rho} \subseteq \mathbb{H}^{2 k+1}$ such that its re-scaling $\mathfrak{s}_{1 / \rho}\left(F_{\rho}\right)$ is cellular, $\mathscr{H}^{2 k+1}\left(\partial F_{\rho}\right) \lesssim \mathscr{H}^{2 k+1}(\partial E)$, and $\mathscr{H}^{2 k+2}\left(E \triangle F_{\rho}\right) \lesssim \rho \mathscr{H}^{2 k+1}(\partial E)=2^{i} \mathscr{H}^{2 k+1}(\partial E)$. Using Lemma 24 and the assumption of Lemma 40 applied to $\mathfrak{s}_{1 / \rho}\left(F_{\rho}\right)$, we see that

$$
\begin{align*}
\left\|\overline{\mathrm{v}}\left(F_{\rho}\right)\right\|_{L_{q}(\mathbb{R})} & =\rho^{2 k+1}\left\|\overline{\mathrm{v}}\left(\mathfrak{s}_{1 / \rho}\left(F_{\rho}\right)\right)\right\|_{L_{q}(\mathbb{R})} \leqslant C \rho^{2 k+1} \mathcal{H}^{2 k+1}\left(\partial \mathfrak{s}_{1 / \rho}\left(F_{\rho}\right)\right)  \tag{66}\\
& =C \mathscr{H}^{2 k+1}\left(\partial F_{\rho}\right) \lesssim C \mathscr{H}^{2 k+1}(\partial E) .
\end{align*}
$$

If $t \in \mathbb{R}$ satisfies $t \geqslant i$, then by Lemma 24.(2), we have

$$
\begin{equation*}
\left|\overline{\mathrm{v}}(E)(t)-\overline{\mathrm{v}}\left(F_{\rho}\right)(t)\right| \leqslant \overline{\mathrm{v}}\left(E \triangle F_{\rho}\right)(t) \stackrel{(28)}{\leqslant} \frac{2 \mathcal{H}^{2 k+2}\left(E \triangle F_{\rho}\right)}{2^{t}} \lesssim 2^{i-t} \mathcal{H}^{2 k+1}(\partial E) \tag{67}
\end{equation*}
$$

Hence,

$$
\begin{gathered}
\|\overline{\mathrm{v}}(E)\|_{L_{q}(i, \infty)} \leqslant\left\|\overline{\mathrm{v}}\left(F_{\rho}\right)\right\|_{L_{q}(i, \infty)}+\left\|\overline{\mathrm{v}}(E)-\overline{\mathrm{v}}\left(F_{\rho}\right)\right\|_{L_{q}(i, \infty)} \\
\stackrel{(66) \wedge(67)}{ } \quad(C+1) \mathscr{H}^{2 k+1}(\partial E) .
\end{gathered}
$$

Since this is true for any $i \in \mathbb{Z}$, we have

$$
\|\bar{v}(E)\|_{L_{q}(\mathbb{R})}=\lim _{i \rightarrow-\infty}\|\overline{\mathrm{v}}(E)\|_{L_{q}(i, \infty)} \lesssim(C+1) \mathscr{H}^{2 k+1}(\partial E)
$$

## 4. Proof overview

By the reductions in Section 3, we know that in order to prove Theorem 3 it suffices to prove Theorem 39. Here we shall sketch in broad strokes the steps of the proof of Theorem 39. The proof uses "intrinsic corona decompositions" in the Heisenberg group to reduce the desired bound to a bound on pieces of intrinsic Lipschitz half-spaces; recall Section 2.3. Each of these intrinsic Lipschitz half-spaces is bounded by the intrinsic graph of a function defined on a vertical plane. We slice the vertical plane into cosets of $\mathbb{H}^{3}$ and bound the vertical perimeter of the intrinsic half-space using an $L_{2}$ inequality for functions on these cosets. In essence, our argument "lifts" an $L_{2}$ inequality for functions on $\mathbb{H}^{3}$ to a formally stronger endpoint $L_{1}$ inequality on $\mathbb{H}^{2 k+1}$. The assumption that $k>1$ arises from this lifting process. When $k=1$ the vertical plane is 2 -dimensional, so we cannot relate the question to a functional inequality on $\mathbb{H}^{3}$, and therefore this strategy does not work. However, the bound on the vertical perimeter of intrinsic Lipschitz half-spaces (Section 5) is the only step of our proof that fails for the 3-dimensional Heisenberg group; when $k=1$, we can still use intrinsic corona decompositions to reduce to the case of intrinsic Lipschitz half-spaces, but the desired inequality fails for these half-spaces.

Intrinsic Lipschitz half-spaces. To prove Theorem 39, we first establish a version of Theorem 36 for intrinsic Lipschitz half-spaces. Specifically, in Section 5 we prove that, provided $k \in\{2,3, \ldots$,$\} , for every \lambda \in(0,1)$, if $E \subseteq \mathbb{H}^{2 k+1}$ is an intrinsic $\lambda$-Lipschitz half-space (recall Section 2.3), then for every point $p \in \mathbb{H}^{2 k+1}$ and every radius $r>0$, we have

$$
\begin{equation*}
\left\|\overline{\mathrm{v}}_{B_{r}(p)}(E)\right\|_{L_{2}(\mathbb{R})} \lesssim \frac{r^{2 k+1}}{1-\lambda} \tag{68}
\end{equation*}
$$

When, say, $\lambda \in\left(0, \frac{1}{2}\right)$, the estimate (68) is in essence the special case of Theorem 36 for (pieces of) intrinsic Lipschitz half-spaces. This is so because, due to the isoperimetric inequality for the Heisenberg group [100], the right-hand side of (68) is at most a constant (depending on $k$ ) multiple of $\mathscr{H}^{2 k+1}\left(\partial\left(B_{r} \cap E\right)\right)$ whenever $\mathscr{H}^{2 k+2}\left(B_{r} \cap E\right) \gtrsim r^{2 k+2}$, i.e., provided that the intrinsic Lipschitz half-space $E$ occupies a constant fraction of the volume of the ball $B_{r}$. The estimate (68) will be used below only in such a nondegenerate situation.

Our proof of (68) relies crucially on a local $L_{2}$-variant of (63) for $\mathbb{H}^{3}$ that was proven in [8] using representation theory; see Theorem 42 below for a precise formulation. We leverage this lower-dimensional functional inequality to deduce (68) as follows: Suppose that $\partial E$ is an intrinsic Lipschitz graph over a $2 k$-dimensional vertical plane $V \subseteq \mathbb{H}^{2 k+1}$; that is, $\partial E=\Gamma_{f}$ for some function $f: V \rightarrow \mathbb{R}$ as in (30). The function $f$ does not have to be Lipschitz with respect to the Carnot-Carathéodory metric, but a computation that appears
in Section 5 shows that the restrictions of $f$ to cosets of copies of $\mathbb{H}^{3}$ in $V$ are in fact Lipschitz with respect to the Carnot-Carathéodory metric. This allows us to use the quadratic inequality of $[8]$ for each of these "fibers" of $f$, and by combining the resulting inequalities (using an application of Cauchy-Schwarz) one arrives at (68). It is important to stress that this proof does not work for functions on $\mathbb{H}^{3}$ because it relies on slicing $V$ into copies of $\mathbb{H}^{3}$. When $V \subseteq \mathbb{H}^{3}$, $V$ is two-dimensional and can be sliced into vertical lines, but there is no analogue of the quadratic inequality of $[8]$ that we use here for vertical lines in $\mathbb{H}^{3}$.

Intrinsic corona decompositions. In order to apply (68) to more general sets, we use intrinsic corona decompositions in the Heisenberg group. The full definitions will appear in Section 6, and a qualitative description appears below. Corona decompositions are an established tool in analysis for reducing the study of certain singular integrals on $\mathbb{R}^{n}$ to the case of Lipschitz graphs, starting with seminal works of David [31], [32] and Jones [51], [52] on the Cauchy integral and culminating with the David-Semmes theory of quantitative rectifiability [30], [33]. Our adaptation of this technique is mostly technical, but it will also involve a conceptually new ingredient, namely, the use of quantitative monotonicity to govern the decomposition.

A corona decomposition covers $\partial E$ by two types of sets, called stoppingtime regions and bad cubes. Stopping-time regions correspond to parts of $\partial E$ that are close to intrinsic Lipschitz graphs, and bad cubes correspond to parts of $\partial E$, like sharp corners, that are not. The multiplicity of this cover depends on the shape of $\partial E$ at different scales. For example, $\partial E$ might look smooth on a large neighborhood of a point $x$, jagged at a medium scale, then smooth again at a small scale. If so, then $x$ is contained in a large stopping-time region, a medium-sized bad cube, and a second small stopping-time region. A cover like this is a corona decomposition if it satisfies a Carleson packing condition (see Section 6) that bounds its average multiplicity on any ball.

Section 7 establishes that Theorem 36 holds when $\partial E$ admits a corona decomposition. We proceed as follows. The vertical perimeter of $\partial E$ comes from three sources: the bad cubes, the graphs that approximate the stoppingtime regions, and the error incurred by approximating a stopping-time region by an intrinsic Lipschitz graph. By the Carleson packing condition, there are few bad cubes, and they contribute vertical perimeter on the order of $\mathcal{H}^{2 k+1}(\partial E)$. Using (68), we argue that the intrinsic Lipschitz graphs also contribute vertical perimeter on the order of $\mathcal{H}^{2 k+1}(\partial E)$. Finally, the difference between a stopping-time region and an intrinsic Lipschitz graph is bounded by the size of the stopping-time region. The stopping-time regions also satisfy a Carleson packing condition, so these errors also contribute vertical perimeter on the order of $\mathscr{H}^{2 k+1}(\partial E)$. Summing these contributions, we obtain the desired bound.

Not every set with finite perimeter admits a corona decomposition, even in $\mathbb{R}^{n}$, but we show that the boundaries of cellular sets can be built out of pieces that admit such decompositions. (A similar technique for surfaces in $\mathbb{R}^{n}$ was used in [112].) We start by showing that in order to establish Theorem 39 it suffices to prove that for every $r>0$, we have $\left\|\bar{v}_{B_{r}}(E)\right\|_{L_{2}(\mathbb{R})} \lesssim r^{2 k+1}$ under the additional assumption that the sets $E, E^{\mathrm{C}}, \partial E$ are $r$-locally Ahlfors-regular. That is, $\mathscr{H}^{2 k+2}\left(u B_{\rho} \cap E\right) \asymp \rho^{2 k+2} \asymp \mathcal{H}^{2 k+2}\left(v B_{\rho} \backslash E\right)$ and $\mathcal{H}^{2 k+1}\left(w B_{\rho} \cap \partial E\right) \asymp$ $\rho^{2 k+1}$ for all $\rho \in(0, r)$ and $(u, v, w) \in E \times E^{c} \times \partial E$. To see this reduction, recall that in Theorem 39 we are given a cellular set $F \subseteq \mathbb{H}^{2 k+1}$. Any such set is Ahlfors-regular on sufficiently small balls. We argue that $F$ can be decomposed into sets that satisfy the desired local Ahlfors-regularity. The full construction of this decomposition is carried out in Section 8, but we remark briefly that it amounts to the following natural "greedy" iterative procedure. If one of the sets $F, F^{\mathrm{c}}, \partial F$ were not locally Ahlfors-regular, then there would be some smallest ball $B$ such that the density of $F, F^{\text {c }}$ or $\partial F$ is either too low or too high on $B$. (It is actually convenient to work here with "cellular balls" here rather than Carnot-Carathéodory balls so as to maintain the inductive hypothesis that $F$ is cellular, but we ignore this technical point for the purpose of the overview; see Section 8.) By replacing $F$ by either $F \cup B$ or $F \backslash B$, we cut off a piece of $\partial F$ and decrease $\mathscr{H}^{2 k+1}(\partial F)$. Since $B$ was the smallest ball where Ahlfors-regularity fails, $F, F^{c}, \partial F$ are Ahlfors-regular on balls smaller than $B$. Repeating this process eventually reduces $F$ to the empty set. We arrive at the conclusion of Theorem 39 for the initial set $F$ by proving the (local version of the) theorem for each piece of this decomposition, then summing the resulting inequalities.

Finally, in Section 9, we arrive at the main technical step of the decomposition procedure: showing that if $E, E^{c}$, and $\partial E$ are Ahlfors regular, then $\partial E$ admits a corona decomposition (Theorem 57). This is the longest section of the proof, but it follows the well-established methods of [30], [33]. We start by constructing a sequence of nested partitions of $\partial E$ into pieces called cubes; this is a standard construction due to Christ [27] and David [32] and only uses the Ahlfors regularity of $\partial E$. These partitions are analogues of the standard tilings of $\mathbb{R}^{n}$ into dyadic cubes. Next, we classify the cubes into good cubes, which are close to a piece of a hyperplane, and bad cubes, which are not. In order to produce a corona decomposition, there cannot be too many bad cubes; i.e., they must satisfy a Carleson packing condition. In [30], [33], this condition follows from quantitative rectifiability; the surface in question is assumed to satisfy a condition that bounds the sum of its (appropriately normalized) local deviations from hyperplanes. These local deviations are higher-dimensional versions of Jones' $\beta$-numbers [51], [52], and the quantitative rectifiability assumption leads to the desired packing condition. In the present setting, the packing condition follows instead from quantitative nonmonotonicity. The concept of
quantitative nonmonotonicity of a set $E \subseteq \mathbb{H}^{2 k+1}$ (see Section 9.1) was defined in [25], [26], where the kinematic formula for the Heisenberg group was used to show that the total nonmonotonicity of all of the cubes is at most a constant multiple of $\mathscr{H}^{2 k+1}(\partial E)$. This means that there cannot be many cubes that have large nonmonotonicity. By a result of [25], [26], if a surface has small nonmonotonicity, then it is close to a half-space. Consequently, most cubes are close to a half-space and are therefore good. (The result in [25], [26] is stronger than what we need for this proof; it provides power-type bounds on how closely a nearly-monotone surface approximates a hyperplane. For our purposes, it is enough to have some bound (not necessarily power-type) on the shape of nearly-monotone sets, and we will deduce the bound that we need by applying a quick compactness argument to a result from [24] that states that if a set is precisely monotone (i.e., every horizontal line intersects the boundary of the set in at most one point), then it is a half-space.)

Next, we partition the good cubes into stopping-time regions by using an iterative construction that corrects overpartitioning that may have occurred when the Christ cubes were constructed. If $Q$ is a largest good cube that has not been treated yet and if $P$ is its approximating half-space, we find all of the descendants of $Q$ with approximating half-spaces that are sufficiently close to $P$. If we glue these half-spaces together using a partition of unity, the result is an intrinsic Lipschitz half-space that approximates all of these descendants. By repeating this procedure for each untreated cube, we obtain a collection of stopping-time regions. These regions satisfy a Carleson packing condition because if a point $x \in \partial E$ is contained in many different stopping-time regions, then either $x$ is contained in many different bad cubes, or $x$ is contained in good cubes whose approximating hyperplanes point in many different directions. In either case, these cubes generate nonmonotonicity, so there can only be a few points with large multiplicity. In fact, this cover satisfies a Carleson packing condition, so it is an intrinsic corona decomposition.

## 5. Isoperimetry of intrinsic Lipschitz graphs

In this section, we bound the (local) vertical perimeter of intrinsic Lipschitz graphs in $\mathbb{H}^{2 k+1}$, provided that $k \geqslant 2$. This is the only place in the proof of Theorem 39 in which the assumption that $k>1$ is used; see Remark 44 below. Specifically, we will establish the following proposition.

Proposition 41. For every integer $k \geqslant 2$ and every $\lambda \in(0,1)$, if $E \subseteq$ $\mathbb{H}^{2 k+1}$ is a $\lambda$-Lipschitz half-space, then for every point $p \in \mathbb{H}^{2 k+1}$ and every radius $r>0$, we have

$$
\begin{equation*}
\left\|\overline{\mathrm{v}}_{B_{r}(p)}(E)\right\|_{L_{2}(\mathbb{R})} \lesssim \frac{r^{2 k+1}}{1-\lambda} \tag{69}
\end{equation*}
$$

Before proving Proposition 41, we shall state Theorem 42 below for ease of later reference. It follows from [8, Th. 7.5], where it is proven using representation theory. A generalization of this result (namely, an $L_{p}$ version for every $p \in(1, \infty)$ as well as more general target spaces) follows from [66, Th. 2.1], where it is proven using different methods than those of $[8]$.

Theorem 42. For every Lipschitz function $f: \mathbb{H}^{3} \rightarrow \mathbb{R}$ and every $r \in$ $(0, \infty)$, we have

$$
\begin{equation*}
\int_{0}^{r^{2}} f_{B_{r}} \frac{\left|f(h)-f\left(h Z^{-t}\right)\right|^{2}}{t^{2}} \mathrm{~d} \mathscr{H}^{4}(h) \mathrm{d} t \lesssim\|f\|_{\operatorname{Lip}\left(\mathbb{H}^{3}\right)}^{2} \tag{70}
\end{equation*}
$$

Remark 43. The statement of [8, Th. 7.5] actually treats smooth functions, asserting that there exists a universal constant $C \in[1, \infty)$ such that for every smooth $f: \mathbb{H}^{3} \rightarrow \mathbb{R}$, we have for every $r \in[1, \infty)$,

$$
\begin{align*}
\int_{1}^{r^{2}} \int_{B_{r}} & \frac{\left|f(h)-f\left(h Z^{-t}\right)\right|^{2}}{t^{2}} \mathrm{~d} \mathscr{H}^{4}(h) \mathrm{d} t  \tag{71}\\
& \lesssim \int_{B_{C r}}\left\|\nabla_{\mathbb{H}} f(h)\right\|_{\ell_{2}^{2}}^{2} \mathrm{~d} \mathscr{H}^{4}(h) .
\end{align*}
$$

When $f$ is Lipschitz but not necessarily smooth, it follows that

$$
\begin{equation*}
\forall r \in[1, \infty), \quad \int_{1}^{r^{2}} f_{B_{r}} \frac{\left|f(h)-f\left(h Z^{-t}\right)\right|^{2}}{t^{2}} \mathrm{~d} \mathscr{H}^{4}(h) \mathrm{d} t \lesssim\|f\|_{\operatorname{Lip}\left(\mathbb{H}^{3}\right)}^{2} \tag{72}
\end{equation*}
$$

Indeed, convolve $f$ on the left with a smooth nonnegative function of arbitrarily small support and integral 1. The result is a smooth function with the same or smaller Lipschitz constant that approximates $f$ in $L_{\infty}\left(B_{r}\right)$. Applying (71) to the resulting function and using the fact that the integrals appearing in (71) are over compact sets, we obtain (72). By scale-invariance, (72) implies (70). Indeed, if $r>0$ and $\varepsilon \in(0, r)$, apply (72) with $r$ replaced by $r / \varepsilon$ and $f$ replaced by $h \mapsto f\left(\mathfrak{s}_{\varepsilon}(h)\right)$, change the variable to $h^{\prime}=\mathfrak{s}_{\varepsilon}(h)$, and let $\varepsilon \rightarrow 0$ in the resulting inequality.

Proof of Proposition 41. Denote $V=\left\{h \in \mathbb{H}^{2 k+1}: x_{k}(h)=\mathbf{0}\right\}$. By applying an isometric automorphism and using Lemma 25, we may assume that $p=\mathbf{0}$ and that $E$ is bounded by an intrinsic Lipschitz graph $\Gamma$ over $V$, say, $\Gamma=\left\{v X_{k}^{f(v)}: v \in V\right\}$ for some continuous function $f: V \rightarrow \mathbb{R}$. The Lipschitz condition implies that $\left|x_{k}\left(w_{1}\right)-x_{k}\left(w_{2}\right)\right| \leqslant \lambda d\left(w_{1}, w_{2}\right)$ for all $w_{1}, w_{2} \in \Gamma$.

We shall identify $\mathbb{H}^{3}$ with the subgroup $\left\langle X_{1}, Y_{1}, Z\right\rangle \triangleleft \mathbb{H}^{2 k+1}$; importantly, $\mathbb{H}^{3}$ commutes with $X_{k}$. We shall also denote the linear subspace

$$
\left\langle X_{2}, \ldots, X_{k-1}, Y_{2}, \ldots, Y_{k}\right\rangle \subseteq \mathrm{H}
$$

by $\mathbb{G}$. Then $V=\mathbb{G}_{\mathbb{H}^{3}}$ in the sense that for each $v \in V$, there are unique elements $g \in \mathbb{G}$ and $h \in \mathbb{H}^{3}$ such that $v=g h$. For every $g \in \mathbb{G}$, define $f_{g}: \mathbb{H}^{3} \rightarrow \mathbb{R}$ by setting $f_{g}(h)=f(g h)$ for every $h \in \mathbb{H}^{3}$.

Fixing $g \in \mathbb{G}$, for every $h_{1}, h_{2} \in \mathbb{H}^{3}$ if we denote $w_{1}=g h_{1} X_{k}^{f_{g}\left(h_{1}\right)}$ and $w_{2}=g h_{2} X_{k}^{f_{g}\left(h_{2}\right)}$, then by the definition of $\Gamma$ we have $w_{1}, w_{2} \in \Gamma$. Moreover, $x_{k}\left(g h_{1}\right)=x_{k}\left(g h_{2}\right)=0$ since $g, h_{1}, h_{2} \in V$, so $x_{k}\left(w_{1}\right)=f_{g}\left(h_{1}\right)$ and $x_{k}\left(w_{2}\right)=$ $f_{g}\left(h_{2}\right)$. Hence, by our assumption on $f$ we have

$$
\begin{align*}
\left|f_{g}\left(h_{1}\right)-f_{g}\left(h_{2}\right)\right| & =\left|x_{k}\left(w_{1}\right)-x_{k}\left(w_{2}\right)\right| \leqslant \lambda d\left(w_{1}, w_{2}\right) \\
& =\lambda d\left(h_{1} X_{k}^{f_{g}\left(h_{1}\right)}, h_{2} X_{k}^{f_{g}\left(h_{2}\right)}\right) \\
& =\lambda d\left(X_{k}^{f_{g}\left(h_{1}\right)} h_{1}, X_{k}^{f_{g}\left(h_{1}\right)} h_{2} X_{k}^{f_{g}\left(h_{2}\right)-f_{g}\left(h_{1}\right)}\right) \tag{73}
\end{align*}
$$

$$
\begin{equation*}
\stackrel{(23)}{\leqslant} \lambda\left(C_{k} d\left(h_{1}, h_{2}\right)+\left|f_{g}\left(h_{1}\right)-f_{g}\left(h_{2}\right)\right|\right), \tag{23}
\end{equation*}
$$

where $C_{k}$ is the constant in (23). This simplifies to give that

$$
\left|f_{g}\left(h_{1}\right)-f_{g}\left(h_{2}\right)\right| \lesssim \frac{\lambda}{1-\lambda} d\left(h_{1}, h_{2}\right) .
$$

Note that this is not quite the condition that $f_{g}$ is $\frac{\lambda}{1-\lambda}$-Lipschitz, because the metric $d$ is the metric on $\mathbb{H}^{2 k+1}$ rather than the metric $d_{\mathbb{H}^{3}}$ on $\mathbb{H}^{3}$. By (23), however, $d_{\mathbb{H}^{3}}\left(h_{1}, h_{2}\right) \asymp d\left(h_{1}, h_{2}\right)$ for all $h_{1}, h_{2} \in \mathbb{H}^{3}$, so $\left\|f_{g}\right\|_{\text {Lip }\left(\mathbb{H}^{3}\right)} \lesssim \lambda /(1-\lambda)$. By Theorem 42, we therefore have

$$
\begin{equation*}
\forall r \in(0, \infty), \quad \int_{0}^{r^{2}} f_{B_{r} \cap \mathbb{H}^{3}}\left|f_{g}(h)-f_{g}\left(h Z^{-t}\right)\right|^{2} \mathrm{~d} \mathscr{H}^{4}(h) \frac{\mathrm{d} t}{t^{2}} \lesssim \frac{\lambda^{2}}{(1-\lambda)^{2}} . \tag{74}
\end{equation*}
$$

Every $u \in \mathbb{H}^{2 k+1}$ can be written uniquely as $u=g h X_{k}^{t}$ with $g \in \mathbb{G}$, $h \in \mathbb{H}^{3}$, and $t \in \mathbb{R}$. Indeed, necessarily $g=\sum_{i=2}^{k-1} x_{i}(u) X_{i}+\sum_{j=2}^{k} y_{j}(u) Y_{j}, h=$ $x_{1}(u) X_{1}+y_{1}(u) Y_{1}+\left(z(u)+x_{k}(u) y_{k}(u) / 2\right) Z$ and $t=x_{k}(u)$. Since $\mathbb{G}, \mathbb{H}^{3},\left\langle X_{k}\right\rangle$ are orthogonal subspaces of $\mathbb{R}^{2 k+1}$, for every $s \in \mathbb{R}$ we have

$$
\begin{align*}
& \overline{\mathrm{v}}_{B_{r}}(E)(s)= \frac{\mathcal{H}^{2 k+2}\left(B_{r} \cap\left(E \triangle E Z^{2^{2 s}}\right)\right)}{=} \\
& \frac{1}{2^{s}} \int_{\mathbb{G}} \int_{\mathbb{H}^{3}} \int_{-\infty}^{\infty}\left|\mathbf{1}_{E}\left(g h X_{k}^{t}\right)-\mathbf{1}_{E}\left(g h X_{k}^{t} Z^{-2^{2 s}}\right)\right| \\
& \cdot \mathbf{1}_{B_{r}}\left(g h X_{k}^{t}\right) \mathrm{d} t \mathrm{~d} \mathscr{\not}^{4}(h) \mathrm{d} \mathscr{\not}^{2 k-3}(g)  \tag{75}\\
&= \frac{1}{2^{s}} \int_{\mathbb{G}} \int_{\mathbb{H}^{3}} \int_{-\infty}^{\infty}\left|\mathbf{1}_{\left\{t \geqslant f_{g}(h)\right\}}-\mathbf{1}_{\left\{t \geqslant f_{g}\left(h Z^{-2^{2 s}}\right)\right\}}\right| \\
& \cdot \mathbf{1}_{B_{r}}\left(g h X_{k}^{t}\right) \mathrm{d} t \mathrm{~d} \mathscr{H}^{4}(h) \mathrm{d} \mathscr{H}^{2 k-3}(g) .
\end{align*}
$$

For every $(g, h, t) \in \mathbb{G} \times \mathbb{H}^{3} \times \mathbb{R}$, we have $\pi\left(g h X_{k}^{t}\right)=g+x_{1}(h) X_{1}+$ $y_{1}(h) Y_{1}+t X_{k}$, where $\pi: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}^{2 k}$ is the canonical projection. By (24),
this implies that if $d\left(\mathbf{0}, g h X_{k}^{t}\right)<r$, then $\|g\|<r$ and $|t|<r$, so

$$
d(\mathbf{0}, h)=d(g, g h) \leqslant d(g, \mathbf{0})+d\left(\mathbf{0}, g h X_{k}^{t}\right)+d\left(g h X_{k}^{t}, g h\right) \leqslant\|g\|+r+|t|<3 r
$$

Hence $\mathbf{1}_{B_{r}}\left(g h X_{k}^{t}\right) \leqslant \mathbf{1}_{B_{r}}(g) \mathbf{1}_{B_{3 r}}(h)$. A substitution of this point-wise inequality into (75) gives

$$
\begin{align*}
\overline{\mathrm{v}}_{B_{r}}(E)(s) \leqslant & \frac{1}{2^{s}} \int_{\mathbb{G}} \int_{\mathbb{H}^{3}} \int_{-\infty}^{\infty}\left|\mathbf{1}_{\left\{t \geqslant f_{g}(h)\right\}}-\mathbf{1}_{\left\{t \geqslant f_{g}\left(h Z^{-2^{2 s}}\right)\right\}}\right| \\
& \cdot \mathbf{1}_{B_{r}}(g) \mathbf{1}_{B_{3 r}}(h) \mathrm{d} t \mathrm{~d} \mathscr{H}^{4}(h) \mathrm{d} \mathscr{H}^{2 k-3}(g) \\
= & \frac{\alpha r^{2 k+1}}{2^{s}} f_{B_{r} \cap \mathbb{G}} f_{B_{3 r} \cap \mathbb{H}^{3}}\left|f_{g}(h)-f_{g}\left(h Z^{-2^{2 s}}\right)\right|  \tag{76}\\
& \cdot \mathrm{d} \mathscr{H}^{4}(h) \mathrm{d} \mathscr{H}^{2 k-3}(g),
\end{align*}
$$

where $\alpha=\alpha_{k}=81 \mathscr{H}^{2 k-3}\left(B_{1} \cap \mathbb{G}\right) \mathscr{H}^{4}\left(B_{1} \cap \mathbb{H}^{3}\right)$. Now,

$$
\begin{align*}
& \int_{-\infty}^{\log _{2} r} \overline{\mathrm{v}}_{B_{r}}(E)(s)^{2} d s \\
&(77) \lesssim r^{4 k+2} \int_{0}^{r^{2}}\left(f_{B_{r} \cap \mathbb{G}} f_{B_{3 r} \cap \mathbb{H}^{3}}\left|f_{g}(h)-f_{g}\left(h Z^{-t}\right)\right| \mathrm{d} \mathscr{H}^{4}(h) \mathrm{d} \mathscr{H}^{2 k-3}(g)\right)^{2} \frac{\mathrm{~d} t}{t^{2}} \\
&(78) \leqslant r^{4 k+2} f_{B_{r} \cap \mathbb{G}} \int_{0}^{r^{2}} f_{B_{3 r} \cap \mathbb{H}^{3}}\left|f_{g}(h)-f_{g}\left(h Z^{-t}\right)\right|^{2} \mathrm{~d} \mathscr{H}^{4}(h) \frac{\mathrm{d} t}{t^{2}} \mathrm{~d} \mathscr{H}^{2 k-3}(g)  \tag{78}\\
&(79) \lesssim \frac{\lambda^{2} r^{4 k+2}}{(1-\lambda)^{2}}
\end{align*}
$$

where in (77) we used (76) and the change of variable $t=2^{2 s}$, in (78) we used Jensen's inequality, and in (79) we used (74). Also, the trivial bound $\overline{\mathrm{v}}_{B_{r}}(E)(s) \leqslant \mathcal{H}^{2 k+2}\left(B_{r}\right) / 2^{s} \lesssim r^{2 k+2} / 2^{s}$ gives

$$
\begin{equation*}
\int_{\log _{2} r}^{\infty} \overline{\mathrm{v}}_{B_{r}}(E)(s)^{2} d s \lesssim_{k} r^{4 k+4} \int_{\log _{2} r}^{\infty} \frac{\mathrm{d} s}{2^{2 s}} \asymp r^{4 k+2} \tag{80}
\end{equation*}
$$

The desired estimate (69) now follows by combining (79) and (80).
Remark 44. Note that Proposition 41 used the assumption that $k>1$ crucially. The above proof would not work for $k=1$ because an analogue of Theorem 42 does not exist for 1-dimensional vertical slices of $\mathbb{H}^{3}$, while for $k \geqslant 2$ we succeeded by slicing $\mathbb{H}^{2 k+1}$ into copies of $\mathbb{H}^{3}$.

## 6. Corona decompositions

In this section we shall formulate a definition of a corona decomposition of a subset of $\mathbb{H}^{2 k+1}$. This definition mimics the analogous notion of a corona decomposition of a subset of $\mathbb{R}^{k}$, as defined by David and Semmes [30]; see also [33, Ch. 3]). Here we need a local variant of this definition, but the key
conceptual difference with the Euclidean case is the utilization of intrinsic Lipschitz graphs. Other differences occur in the ensuing proofs: Some of them are of a more technical nature, arising from geometric peculiarities of the Heisenberg group that differ from their Euclidean counterparts, but conceptually new issues arise as well, such as the use of (quantitative) monotonicity methods [24], [26] for the purpose of constructing the desired decomposition.

Definition 45 (local Ahlfors regularity). Suppose that $C, r, s \in(0, \infty)$. We shall say that a Borel subset $A \subseteq \mathbb{H}^{2 k+1}$ is $(C, r)$-Ahlfors $s$-regular if for every $x \in A$ and and $\rho \in(0, r]$, we have

$$
\frac{r^{s}}{C} \leqslant \mathcal{H}^{s}\left(B_{\rho}(x) \cap A\right) \lesssim C r^{s} .
$$

This terminology is shorthand for what one would normally refer to as $A$ being $r$-locally Ahlfors $s$-regular with constant $C$, but since we need to state this often in what follows, we prefer to use the above nonstandard shorter nomenclature. In fact, it will be convenient to introduce the following convention for even shorter terminology. Whenever we are given $E \subseteq \mathbb{H}^{2 k+1}$ and we say that $E$ is $(C, r)$-regular it will be understood by default that $E$ is assumed to be $(C, r)$ Ahlfors $s$-regular with $s=2 k+2$. When we say that $\partial E$ is $(C, r)$-regular we mean that $\partial E$ is assumed to be $(C, r)$-Ahlfors $s$-regular with $s=2 k+1$. When we say that $\left(E, E^{\mathrm{c}}, \partial E\right)$ is $(C, r)$-regular we mean that $E, E^{\mathrm{c}}$ are both $(C, r)$-Ahlfors $(2 k+2)$-regular and also $\partial E$ is $(C, r)$-Ahlfors $(2 k+1)$-regular.

Definition 46 (local cubical patchwork). Fix $K, s, r \in(0, \infty)$ and $n \in \mathbb{Z}$ for which $2^{n} \leqslant r<2^{n+1}$. Suppose that $A$ is a Borel subset of $\mathbb{H}^{2 k+1}$. A sequence $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$ of Borel partitions of $A$ is said to be an $s$-dimensional ( $K, r$ )-cubical patchwork for $A$ if it has the following properties:
(1) For every integer $i \leqslant n$ and every $Q \in \Delta_{i}$, we have

$$
\begin{equation*}
\frac{2^{i}}{K}<\operatorname{diam}(Q)<K 2^{i} \quad \text { and } \quad \frac{2^{i s}}{K}<\mathscr{H}^{s}(Q)<K 2^{i s} \tag{81}
\end{equation*}
$$

(2) For every two integers $i, j \leqslant n$ with $i \leqslant j$, the partition $\Delta_{i}$ is a refinement of the partition $\Delta_{j}$; i.e., for every $Q \in \Delta_{i}$ and $Q^{\prime} \in \Delta_{j}$, either $Q \cap Q^{\prime}=\varnothing$ or $Q \subseteq Q^{\prime}$.
(3) For every integer $i \leqslant n, Q \in \Delta_{i}$ and $t>0$, we have

$$
\begin{equation*}
\mathscr{H}^{s}\left(\partial_{\leqslant t 2^{i}} Q\right) \leqslant K t^{\frac{1}{K}} 2^{i s}, \tag{82}
\end{equation*}
$$

where for every $\rho>0$, we denote

$$
\begin{equation*}
\partial_{\leqslant \rho} Q \stackrel{\text { def }}{=}\{x \in Q: d(x, A \backslash Q) \leqslant \rho\} \cup\{x \in A \backslash Q: d(x, Q) \leqslant \rho\} . \tag{83}
\end{equation*}
$$

In the above setting, write

$$
\begin{equation*}
\Delta \stackrel{\text { def }}{=} \bigsqcup_{i=-\infty}^{n} \Delta_{i} \text {. } \tag{84}
\end{equation*}
$$

The elements of $\Delta$ (that is, any atom of one of the hierarchical partitions $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$ ) are often called Christ cubes. The disjoint union in (84) indicates that we are slightly abusing notation by considering the elements of $\Delta_{i}$ and $\Delta_{j}$ to be distinct for $i \neq j$. In other words, if $Q \in \Delta_{i}$ for multiple integers $i \leqslant n$ then $Q$ appears in $\Delta$ with that multiplicity, which by (81) can be at most $\left(2 \log _{2} K\right) / s+1$. We impose a natural partial order on $\Delta$ by saying that $Q, Q^{\prime} \in \Delta$ satisfy $Q \preceq Q^{\prime}$ if and only if $Q \subseteq Q^{\prime}$ and $Q \in \Delta_{i}, Q^{\prime} \in \Delta_{j}$ for some $i \leqslant j$. Below we shall slightly abuse notation by letting $Q \subseteq Q^{\prime}$ have the same meaning as $Q \preceq Q^{\prime}$ for every $Q, Q^{\prime} \in \Delta$. (Note that this can be an actual abuse of notation only if $Q=Q^{\prime}$.) For every integer $i \leqslant n$ and $Q \in \Delta_{i}$, we shall define $\sigma(Q)=2^{i}$. (We think of $\sigma(Q)$ as a surrogate for the "side-length" of the Christ cube $Q$.) Thus if $\Delta$ contains multiple copies of $Q$, then each such copy has a different " $\sigma$-value," corresponding to those indices $i \leqslant n$ for which $Q \in \Delta_{i}$. Under these conventions, the condition (81) becomes the requirement that for every $Q \in \Delta$, we have

$$
\begin{equation*}
\frac{\sigma(Q)}{K}<\operatorname{diam}(Q)<K \sigma(Q) \quad \text { and } \quad \frac{\sigma(Q)^{s}}{K}<\mathscr{H}^{s}(Q)<K \sigma(Q)^{s}, \tag{85}
\end{equation*}
$$

and the condition (82) becomes the requirement that for every $Q \in \Delta$ and $t>0$, we have

$$
\begin{equation*}
\mathscr{H}^{s}\left(\partial_{\leqslant t \sigma(Q)} Q\right) \leqslant K t^{\frac{1}{K}} \sigma(Q)^{s} . \tag{86}
\end{equation*}
$$

While the above definition makes sense for every $s>0$, in the present article we will only consider cubical patchworks of boundaries of open sets of $\mathbb{H}^{2 k+1}$, in which case it will be tacitly assumed that $s=2 k+1$. Thus given $E \subseteq \mathbb{H}^{2 k+1}$, when we say below that $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$ is a $(K, r)$-cubical patchwork for $\partial E \subseteq \mathbb{H}^{2 k+1}$ this will always be understood to mean that $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$ is a $(2 k+1)$-dimensional $(K, r)$-cubical patchwork for $\partial E$ in the above sense.

Christ [27] and David [32] proved that for every $C, r, s>0$, there exists $K=K(C, s)$ such that if a Borel subset $A \subseteq \mathbb{H}^{2 k+1}$ is $(C, r)$-Ahlfors $s$-regular, then $A$ also admits an $s$-dimensional ( $K, r$ )-cubical patchwork. Note that the statements of [27], [32] assume global Ahlfors regularity rather than local Ahlfors regularity; i.e., the results are stated in the case $r=\operatorname{diam}(A)$. Nevertheless, the proofs of [27], [32] work identically for general $r$ : Discard the cubes that these proofs construct for larger scales (they are irrelevant for the desired local conclusion), so that the Ahlfors regularity is needed only in the range of scales in which it is hypothesized to hold true.

For ease of later reference (specifically, in Section 9 below), we record the following very simple fact, which follows by contrasting (85) with (86); see also [33, Lemma I.3.5].

Lemma 47 (existence of "approximate centers" for Christ cubes). For every $K, r, s>0$, there exists $c=c(K)$ such that if a Borel subset $A \subseteq \mathbb{H}^{2 k+1}$ admits an s-dimensional ( $K, r$ )-cubical patchwork $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$, then for every $Q \in \Delta$, there is a point $x_{Q} \in Q$ such that

$$
d\left(x_{Q}, A \backslash Q\right)>c \sigma(Q)
$$

Proof. Choose any $c<1 / K^{2 K}$. Fix $Q \in \Delta$, and suppose for the purpose of obtaining a contradiction that $d(x, A \backslash Q) \leqslant c \sigma(Q)$ for every $x \in Q$. Recalling (83), this means that $\partial_{\leqslant c \sigma(Q)} Q \supseteq Q$. Therefore,

$$
\frac{\sigma(Q)^{s}}{K} \stackrel{(85)}{\leqslant} \mathscr{H}^{s}(Q) \leqslant \mathscr{H}^{s}\left(\partial_{\leqslant c \sigma(Q)} Q\right) \stackrel{(86)}{\leqslant} K c^{\frac{1}{K}} \sigma(Q)^{s}
$$

By cancelling $\sigma(Q)^{s}$ from both sides, this contradicts our choice of $c$.
The following very important definition is equivalent to Definition 3.9 in [33], though note that the constant $C$ that appears there differs from the constant $C$ that we use below (by a factor that depends only on the parameters of the given cubical patchwork).

Definition 48 (Carleson packing condition). Fix $r, s, K, C \in(0, \infty)$ and suppose that $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$ is an $s$-dimensional $(K, r)$-cubical patchwork of a Borel subset $A \subseteq \mathbb{H}^{2 k+1}$. If $\mathscr{D} \subseteq \Delta$ is a collection of Christ cubes, we say that $\mathcal{D}$ satisfies the $s$-dimensional Carleson packing condition with constant $C$ if for every $Q \in \Delta$, we have

$$
\sum_{\substack{R \in \mathscr{D} \\ R \subseteq Q}} \sigma(R)^{s} \leqslant C \sigma(Q)^{s}
$$

When $s=2 k+1$ it will be convenient to use below the shorter terminology " $D$ is $C$-Carleson" to indicate that $\mathscr{D}$ satisfies the $(2 k+1)$-dimensional Carleson packing condition with constant $C$.

A Carleson packing condition expresses the fact that $\mathscr{D}$ is a "small" set of cubes. It implies, for instance, that $\mathscr{H}^{s}$-almost every point $x \in A$ is contained in only finitely many elements of $\mathscr{D}$.

Definition 49 (coherent collection of cubes). Fix $r, s, K \in(0, \infty)$, and suppose that $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$ is an $s$-dimensional $(K, r)$-cubical patchwork of a Borel subset $A \subseteq \mathbb{H}^{2 k+1}$. If $\mathcal{S} \subseteq \Delta$ is a sub-collection of Christ cubes, then we say that $\mathcal{S}$ is coherent if the following three properties hold true:
(1) $\mathcal{S}$ has a (necessarily unique) maximal element with respect to inclusion; i.e, there exists a unique cube $\mathrm{Q}(\mathcal{S}) \in \mathcal{S}$ such that $\mathrm{Q}(\mathcal{S}) \supseteq Q$ for every $Q \in \mathcal{S}$.
(2) If $Q \in \mathcal{S}$ and $Q^{\prime} \in \Delta$ satisfies $Q \subseteq Q^{\prime} \subseteq Q(\mathcal{S})$, then also $Q^{\prime} \in \mathcal{S}$.
(3) If $Q \in \mathcal{S}$, then either all the children of $Q$ in $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$ belong to $\mathcal{S}$ or none of them do; i.e., if $\sigma(Q)=2^{i}$, then either $\left\{Q^{\prime} \in \Delta_{i+1}: Q^{\prime} \subseteq Q\right\} \subseteq \mathcal{S}$ or $\left\{Q^{\prime} \in \Delta_{i+1}: Q^{\prime} \subseteq Q\right\} \cap \mathcal{S}=\varnothing$.

Using Definitions 48 and 49, and following [33, Def. 3.13], we formulate the notion of a coronization, which is a partition of a cubical patchwork $\Delta$ into "good" and "bad" sets.

Definition 50 (local coronization). Fix $K, C, r>0$ and $E \subseteq \mathbb{H}^{2 k+1}$. A $(K, C, r)$-coronization of $\partial E$ is a triple $(\mathscr{B}, \mathcal{L}, \mathcal{F})$ with the following properties. There exists a $(K, r)$-cubical patchwork $\left\{\Delta_{i}\right\}_{i=-\infty}^{n}$ for $\partial E$ such that $\mathscr{B} \subseteq \Delta$ (the set of bad cubes) and $\mathcal{L} \subseteq \Delta$ (the set of good cubes) partition $\Delta$ into two disjoint sets, i.e., $\mathcal{B} \cup \mathcal{L}_{\mathcal{L}}=\Delta$ and $\mathcal{B} \cap \mathcal{L}_{\mathcal{L}}=\varnothing$, and $\mathscr{F} \subseteq 2^{\mathscr{L}}$ is a collection of subsets of $\mathcal{L}$, which below are called stopping-time regions. These sets are required to have the following properties:
(1) $\mathscr{B}$ is $C$-Carleson.
(2) The elements of $\mathcal{F}$ are pairwise disjoint and their union is $\mathcal{L}_{\mathscr{L}}$.
(3) Each $\mathcal{S} \in \mathscr{F}$ is coherent.
(4) The set of maximal cubes $\{\mathrm{Q}(\mathcal{S}): \mathcal{S} \in \mathscr{F}\}$ is $C$-Carleson.

Definition 51. If $U, V, W \subseteq \mathbb{H}^{2 k+1}$, we define the $U$-local distance between $V$ and $W$ by

$$
d_{U}(V, W) \stackrel{\text { def }}{=} \inf \left\{r:(V \triangle W) \cap U \subseteq \operatorname{nbhd}_{r}(\partial V) \cap \operatorname{nbhd}_{r}(\partial W)\right\}
$$

where we denote $\operatorname{nbhd}_{r}(A)=\left\{h \in \mathbb{H}^{2 k+1}: d(h, A)<r\right\}$ for every $A \subseteq \mathbb{H}^{2 k+1}$.
The local distance is not a metric on subsets of $\mathbb{H}^{2 k+1}$; it only satisfies the weaker version of the triangle inequality given below.

LEMMA 52. If $U, V, W, X \subseteq \mathbb{H}^{2 k+1}$, then the following assertions hold true:
(1) $d_{U}(V, W)=d_{U}(W, V)$.
(2) If $\mathbb{H}^{2 k+1} \supseteq U^{\prime} \supseteq U$, then $d_{U}(V, W) \leqslant d_{U^{\prime}}(V, W)$.
(3) $U \cap \partial V \subseteq \overline{\operatorname{nbhd}}_{s}(\partial W)$, where $s=d_{U}(V, W)$.
(4) If $r=\max \left\{d_{U}(V, W), d_{U}(W, X)\right\}$ and $U^{\prime}=\operatorname{nbhd}_{r}(U)$, then

$$
d_{U}(V, X) \leqslant d_{U^{\prime}}(V, W)+d_{U^{\prime}}(W, X)
$$

Proof. The first two properties follow from the definition. For the third property, note that if $x \in U \cap \partial V$ but $x \notin \partial W$, then there is a neighborhood of $x$ that contains points of $V$ and of $V^{c}$ but is disjoint from $\partial W$. Consequently
$x \in U \cap(\overline{V \triangle W})$, so $d(x, \partial W) \leqslant s$, as desired. For the fourth property, suppose that $x \in(V \triangle X) \cap U$. Then $x \in V \triangle W$ or $x \in W \triangle X$. Without loss of generality, suppose that $x \in V \triangle W$. Then $d(x, \partial V) \leqslant d_{U}(V, W)$ and $d(x, \partial W) \leqslant d_{U}(V, W)$. Let $w \in \partial W$ be such that $d(x, w)=d(x, \partial W)<r$. Then $w \in U^{\prime}$, so by property (3) above, $d(w, \partial X) \leqslant d_{U^{\prime}}(W, X)$ and

$$
\begin{aligned}
d(x, \partial X) & \leqslant d(x, w)+d(w, \partial X) \\
& \leqslant d_{U}(V, W)+d_{U^{\prime}}(W, X) \leqslant d_{U^{\prime}}(V, W)+d_{U^{\prime}}(W, X)
\end{aligned}
$$

A corona decomposition of a set is a coronization where every stoppingtime region is close to an intrinsic Lipschitz graph. In what follows, if $\Delta$ is a cubical patchwork for $A \subseteq \mathbb{H}^{2 k+1}$, then for every $Q \in \Delta$ and $\rho>0$, denote $N_{\rho}(Q)=\operatorname{nbhd}_{\rho \sigma(Q)}(Q)$ and $\rho Q=A \cap N_{\rho}(Q)$.

Definition 53 (local corona decomposition). Fix $K, r, C, \lambda, \theta>0$. Given $E \subseteq \mathbb{H}^{2 k+1}$, we say that the pair $(E, \partial E)$ admits a $(K, C, \lambda, \theta, r)$-corona decomposition if there is a $(K, C, r)$-coronization $(\mathscr{B}, \mathscr{G}, \mathcal{F})$ of $\partial E$ such that for each $\mathcal{S} \in \mathcal{F}$, there is an intrinsic $\lambda$-Lipschitz graph $\Gamma(\mathcal{S})$ that bounds a Lipschitz half-space $\Gamma^{+}(\mathcal{S})$, such that for all $Q \in \mathcal{S}$, we have

$$
\begin{equation*}
d_{N_{4}(Q)}\left(\Gamma^{+}(\mathcal{S}), E\right) \leqslant \theta \sigma(Q) \tag{87}
\end{equation*}
$$

We say that the pair $(E, \partial E)$ admits a ( $K, r$ )-corona decomposition if for every $\lambda, \theta>0$, there exists $C=C(\lambda, \theta)$ such that the pair $(E, \partial E)$ admits a ( $K, C, \lambda, \theta, r$ )-corona decomposition.

The number 4 in (87) is arbitrary; we can replace it with an arbitrarily large constant at the cost of increasing the Carleson packing constants of the coronization.

The following covering lemma will be helpful both to construct corona decompositions and to bound the vertical perimeter of a set with a corona decomposition. Fix $K, r>0$ and $E \subseteq \mathbb{H}^{2 k+1}$. Suppose that $\Delta$ is a $(K, r)$-cubical patchwork of $\partial E$. If $\mathcal{S} \subseteq \Delta$ is a coherent set, we define $d_{\mathcal{S}}: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ by

$$
\forall y \in \mathbb{H}^{2 k+1}, \quad d_{\mathcal{S}}(y)=\inf _{Q \in \mathcal{S}}\{d(y, Q)+\sigma(Q)\}
$$

Denote $\Gamma_{0}=d_{\mathcal{S}}^{-1}(0) \subseteq \mathbb{H}^{2 k+1}$. Since $d_{\mathcal{S}}$ is an infimum of 1-Lipschitz functions, it is also 1-Lipschitz.

Lemma 54. Let $\mathcal{S} \subseteq \Delta$ be a coherent set, and suppose that $0<\psi<\frac{1}{10}$. There is a countable set of points $\left\{c_{i}\right\}_{i \in I} \subseteq 10 \mathrm{Q}(\mathcal{S})$ such that the balls $\left\{B_{i}=\right.$ $\left.B_{\psi d_{\mathcal{S}}\left(c_{i}\right)}\left(c_{i}\right)\right\}_{i \in I}$ are disjoint and the balls $\left\{3 B_{i}=B_{3 \psi d_{\mathcal{S}}\left(c_{i}\right)}\left(c_{i}\right)\right\}_{i \in I}$ satisfy

$$
10 \mathrm{Q}(\mathcal{S}) \subseteq \bigcup_{i \in I} 3 B_{i} \cup \Gamma_{0}
$$

Furthermore, if $0<\rho<1 / \psi$, then the balls $\left\{\rho B_{i}=B_{\rho \psi d_{s}\left(c_{i}\right)}\left(c_{i}\right)\right\}_{i \in I}$ have bounded multiplicity with a bound depending on $\rho$ and $\psi$.

Proof. Write $G=10 \mathrm{Q}(\mathcal{S}) \backslash \Gamma_{0}$. Let $\left\{c_{i}\right\}_{i \in I} \subseteq G$ be a maximal set of points such that the balls $\left\{B_{i}=B_{\psi d_{S}\left(c_{i}\right)}\left(c_{i}\right)\right\}_{i \in I}$ are disjoint; this set is either finite or countably infinite. For any $v \in G$, there is an $i \in I$ such that $B_{\psi d_{S}(v)}(v)$ intersects $B_{i}$, so $d\left(v, c_{i}\right) \leqslant \psi\left(d_{\mathcal{S}}(v)+d_{\mathcal{S}}\left(c_{i}\right)\right)$. Since the mapping $d_{\mathcal{S}}$ is 1-Lipschitz, it follows that the following estimates hold true:

$$
\begin{aligned}
d\left(v, c_{i}\right) & \leqslant \psi\left(2 d_{\delta}\left(c_{i}\right)+d\left(v, c_{i}\right)\right), \\
\frac{9}{10} d\left(v, c_{i}\right) & \leqslant 2 \psi d_{\mathcal{S}}\left(c_{i}\right), \\
d\left(v, c_{i}\right) & \leqslant 3 \psi d_{\mathcal{S}}\left(c_{i}\right) .
\end{aligned}
$$

Hence $v \in 3 B_{i}$, and thus $G \subseteq \bigcup_{i \in I} 3 B_{i}$. Next, suppose that $\rho \psi<1$. If $i \in I$ and $v \in \rho B_{i}$, then

$$
\frac{1}{1+\rho \psi} d_{\mathcal{S}}(v) \leqslant d_{\mathcal{S}}\left(c_{i}\right) \leqslant \frac{1}{1-\rho \psi} d_{\mathcal{S}}(v) .
$$

It follows that there is $\beta=\beta(\rho, \psi)>0$ such that $B_{i}$ is a ball of radius at least $\beta d_{\delta}(v)$ with center in $B_{d_{\mathcal{S}}(v) / \beta}(v)$. Since the balls $\left\{B_{i}\right\}_{i \in I}$ are disjoint, the Ahlfors-regularity of $\mathbb{H}^{2 k+1}$ implies that there can be only boundedly many such balls, and thus only finitely many $i \in I$ such that $v \in \rho B_{i}$.

## 7. Isoperimetry of sets with corona decompositions

Here we show that Proposition 41 holds not only for intrinsic Lipschitz graphs, but also for subsets of $\mathbb{H}^{2 k+1}$ that have a corona decomposition. Specifically, we prove the following proposition.

Proposition 55. Fix $K, C, \Lambda, r \in(0, \infty)$, $\theta \in\left(0, \frac{1}{240}\right), \lambda \in(0,1)$ and $p \in[1, \infty)$. Suppose that for every intrinsic $\lambda$-Lipschitz half-space $\Gamma^{+} \subseteq \mathbb{H}^{2 k+1}$, every $R \in(0, \infty)$ and every $q \in \mathbb{H}^{2 k+1}$, we have

$$
\begin{equation*}
\left\|\overline{\mathrm{v}}_{B_{R}(q)}\left(\Gamma^{+}\right)\right\|_{L_{p}(\mathbb{R})} \leqslant \Lambda R^{2 k+1} . \tag{88}
\end{equation*}
$$

Then for every $E \subseteq \mathbb{H}^{2 k+1}$ such that $(E, \partial E)$ admits a ( $\left.K, C, \lambda, \theta, r\right)$-corona decomposition and any $x \in \mathbb{H}^{2 k+1}$, we have

$$
\left\|\overline{\mathrm{v}}_{B_{r}(x)}(E)\right\|_{L_{p}(\mathbb{R})} \lesssim_{k, K} C(\Lambda+1) r^{2 k+1} .
$$

Note that if $k \in\{2,3, \ldots\}$, then by Proposition 41 the assumption of Proposition 55 holds true with $p=2$ and $\Lambda \lesssim 1 /(1-\lambda)$. We stated Proposition 55 for general $p \in[1, \infty)$ rather than only for $p=2$ in anticipation of forthcoming work that treats the case $k=1$. In the next sections, we will
show that surfaces of finite perimeter in $\mathbb{H}^{2 k+1}$ can be decomposed into sets with corona decompositions and use Proposition 55 to prove Theorem 36.

The proof of Proposition 55 uses the following lemma.
Lemma 56. Suppose that $\psi \in\left(0, \frac{1}{10}\right), K, r \in(0, \infty), \lambda \in(0,1), \theta \in$ $\left(0, \frac{\psi}{12}\right)$. Let $E \subseteq \mathbb{H}^{2 k+1}$ be a set admitting a (K,r)-cubical patchwork $\Delta$. Suppose that $\mathcal{S} \subseteq \Delta$ is a coherent set with maximal cube $M$, and let $\Gamma^{+}$be a Lipschitz half-space bounded by an intrinsic $\lambda$-Lipschitz graph $\Gamma$. Let $\left\{c_{i}\right\}_{i \in I}$ and $\left\{B_{i}\right\}_{i \in I}$ be as in Lemma 54. If $\mathcal{S}$ and $\Gamma$ satisfy (87), then

$$
\left(\Gamma^{+} \triangle E\right) \cap N_{4}(M) \subseteq \bigcup_{i \in I} 4 B_{i} \cup \Gamma
$$

Proof. Let $F=\Gamma^{+} \triangle E$, and suppose that $y \in F \cap N_{4}(M)$. If $d_{\mathcal{S}}(y)=0$, then for every $\varepsilon \in(0, \sigma(M)]$, there is a $Q \in \mathcal{S}$ such that $d(y, Q)+\sigma(Q)<\varepsilon$. Let $Q^{\prime}$ be an ancestor of $Q$ such that $\frac{\varepsilon}{2} \leqslant \sigma\left(Q^{\prime}\right) \leqslant \varepsilon$. Since $d(y, Q)<\varepsilon$, we have $y \in N_{4}\left(Q^{\prime}\right)$ and $d\left(y, Q^{\prime}\right)+\sigma\left(Q^{\prime}\right)<2 \varepsilon$. By ( 87$), d(x, \Gamma) \leqslant \theta \sigma(Q) \leqslant \theta \varepsilon$. Since this holds for every $\varepsilon$, we have $y \in \Gamma$.

We can therefore suppose from now on that $y \in F \cap N_{4}(M)$ and that $0<d_{\mathcal{S}}(y) \leqslant 5 \sigma(M)$. If $d_{\mathcal{S}}(y) \geqslant \frac{\sigma(M)}{4}$, let $A=M$. Then $y \in N_{4}(A)$ and $\frac{d_{\delta}(y)}{5} \leqslant \sigma(A) \leqslant 4 d_{\delta}(y)$. Otherwise, let $Q \in \mathcal{S}$ be a cube such that $d(y, Q)+$ $\sigma(Q) \leqslant 2 d_{\mathcal{S}}(y)$, and let $A$ be the ancestor of $Q$ such that $\frac{d_{\delta}(y)}{5} \leqslant 2 d_{\delta}(y) \leqslant$ $\sigma(A)<4 d_{\mathcal{S}}(y)$. Then $d(y, A) \leqslant 2 d_{\delta}(y) \leqslant \sigma(A)$, so in either case, we have $y \in N_{4}(A)$ and $\frac{d_{\delta}(y)}{5} \leqslant \sigma(A) \leqslant 4 d_{\mathcal{\delta}}(y)$. By (87), this implies that

$$
d(y, \partial E) \leqslant \theta \sigma(A) \leqslant \frac{\psi d_{\delta}(y)}{3} .
$$

Choose $v \in \partial E$ with $d(v, y) \leqslant \frac{\psi}{2} d_{\mathcal{S}}(y)$. Then $d(v, y) \leqslant \sigma(M)$, so $v \in 5 M$. Since $d_{\mathcal{S}}$ is 1 -Lipschitz, we have $d_{\mathcal{S}}(v) \geqslant(1-\psi) d_{\mathcal{S}}(y)>0$. By Lemma 54, there is $i \in I$ such that $v \in 3 B_{i}$, and

$$
d_{\mathcal{S}}(y) \leqslant \frac{1}{1-\psi} d_{\mathcal{S}}(v) \leqslant \frac{1+3 \psi}{1-\psi} d_{\mathcal{S}}\left(c_{i}\right) \leqslant 2 d_{\mathcal{S}}\left(c_{i}\right) .
$$

Thus

$$
d\left(c_{i}, y\right) \leqslant d\left(c_{i}, v\right)+d(v, y) \leqslant 3 \psi d_{\mathcal{S}}\left(c_{i}\right)+\frac{\psi}{2} d_{\mathcal{S}}(y) \leqslant 4 \psi d_{\delta}\left(c_{i}\right) .
$$

Hence $y \in 4 B_{i}$, as desired.
Proof of Proposition 55. By re-scaling and translating $E$ we may assume that $x=\mathbf{0}$ and $r=1$. Let $\Delta=\left\{\Delta_{i}\right\}_{i=-\infty}^{0}$ be a $(K, r)$-cubical patchwork for $\partial E$, and let $(\mathscr{B}, \mathcal{G}, \mathscr{F})$ be a coronization of $\partial E$ that satisfies Definition 53. In particular, $\mathscr{B}$ and $\{\mathrm{Q}(\mathcal{S}): \mathcal{S} \in \mathscr{F}\}$ are both $C$-Carleson.

Recall that for every $s \in \mathbb{R}$, we let $\mathrm{D}_{s} E=E \triangle E Z^{2^{2 s}}$. We start by constructing a cover of $\mathrm{D}_{s} E$. Suppose that $s<0$, and let $i$ be the integer such
that $s \in(i-1, i]$. If $x \in \mathrm{D}_{s} E$, then any path from $x$ to $x Z^{-2^{2 s}}$ crosses $\partial E$. There is a path from $x$ to $x Z^{-2^{2 s}}$ of length less than $4 \cdot 2^{s}$, and this path intersects $\partial E$ in a point $y$ such that $d(x, y)<4 \cdot 2^{s}$. If $Q \in \Delta_{i}$ is the cube such that $y \in Q$, then $x \in N_{4}(Q)$. It follows that

$$
\begin{equation*}
\mathrm{D}_{s} E \subseteq \bigcup_{Q \in \Delta_{i}} N_{4}(Q) . \tag{89}
\end{equation*}
$$

Due to (89), we can decompose $\overline{\mathrm{v}}(E)$ into terms corresponding to cubes. For $Q \in \Delta$, we let the $Q$-component of $\overline{\mathrm{v}}(E)$ be

$$
\overline{\mathrm{v}}(E ; Q)(s) \stackrel{\text { def }}{=} \overline{\mathrm{V}}_{N_{4}(Q)}(E)(s) \mathbf{1}_{\left\{\frac{\sigma(Q)}{2}<2^{s} \leqslant \sigma(Q)\right\}} .
$$

Then for each integer $i \leqslant 0$ and $s \in(i-1, i]$, we have

$$
\sum_{Q \in \Delta} \overline{\mathrm{v}}(E ; Q)(s)=\sum_{Q \in \Delta_{i}} \overline{\mathrm{v}}(E ; Q)(s)=2^{-s} \sum_{Q \in \Delta_{i}} \mathscr{H}^{2 k+2}\left(N_{4}(Q) \cap \mathrm{D}_{s} E\right) .
$$

As $Q$ ranges over $\Delta_{i}$, the sets $N_{4}(Q)$ cover $\mathrm{D}_{s} E$ with multiplicity that is bounded by a quantity that depends only on $k$ and $K$, so it follows that

$$
\forall s \in(-\infty, 0], \quad \overline{\mathrm{v}}(E)(s) \asymp \sum_{Q \in \Delta} \overline{\mathrm{v}}(E ; Q)(s) .
$$

For any subset $\mathcal{J} \subseteq \Delta$ and any $F \subseteq \mathbb{H}^{2 k+1}$, denote $W_{\mathfrak{g}}(F) \stackrel{\text { def }}{=} \sum_{Q \in \mathfrak{J}} \overline{\mathrm{v}}(F ; Q)$. If $M \in \Delta$, let $\Delta_{M}=\{Q \in \Delta: Q \subseteq M\}$. We claim that

$$
\begin{equation*}
\forall M \in \Delta, \quad\left\|W_{\Delta_{M}}(E)\right\|_{L_{p}(\mathbb{R})} \lesssim_{k, K} C(\Lambda+1) \sigma(M)^{2 k+1} \tag{90}
\end{equation*}
$$

To prove (90), fix $M \in \Delta$. We first separate $W_{M}(E)$ into contributions of bad cubes and stopping-time regions. Let $\mathscr{B}_{M}=\Delta_{M} \cap \mathscr{B}$, and let $\mathscr{F}_{M}=$ $\left\{\mathcal{S} \cap \Delta_{M}: \mathcal{S} \in \mathscr{F}\right\}$. Note that $\mathscr{F}_{M}$ consists of coherent collections of cubes, so that the maximal element $\mathrm{Q}(\mathscr{T}) \in \mathscr{T}$ is well defined for every $\mathcal{T} \in \mathscr{F}_{M}$. Moreover, the elements of $\mathscr{F}_{M}$ still satisfy (87) and $\left\{\mathrm{Q}(\mathcal{S}): \mathcal{S} \in \mathscr{F}_{M}\right\}$ is still $C$-Carleson. By design we have

$$
\begin{equation*}
\forall s \in(-\infty, 0], \quad W_{\Delta_{M}}(E)(s)=W_{\mathscr{B}_{M}}(E)(s)+\sum_{\mathscr{T} \in \mathscr{F}_{M}} W_{\mathcal{T}}(E)(s) . \tag{91}
\end{equation*}
$$

We will prove (90) by bounding each of the terms in (91) separately.
First, we bound $\left\|W_{\mathscr{B}_{M}}(E)\right\|_{L_{p}(\mathbb{R})}$. If $Q \in \Delta_{i}$ and $s \in(i-1, i]$, we have the trivial bound

$$
\overline{\mathrm{v}}(E ; Q)(s) \leqslant \frac{\mathscr{H}^{2 k+2}\left(N_{4}(Q)\right)}{2^{s}} \asymp_{k, K} \frac{2^{i(2 k+2)}}{2^{s}} \asymp_{k, K} \sigma(Q)^{2 k+1},
$$

and consequently,

$$
\begin{equation*}
\|\overline{\mathrm{v}}(E ; Q)\|_{L_{p}(\mathbb{R})} \lesssim\left(\int_{i-1}^{i} \overline{\mathrm{v}}(E ; Q)(s)^{p} \mathrm{~d} s\right)^{\frac{1}{p}} \lesssim_{k, K} \sigma(Q)^{2 k+1} \tag{92}
\end{equation*}
$$

Hence,

$$
\begin{equation*}
\left\|W_{\mathscr{B}_{M}}(E)\right\|_{L_{p}(\mathbb{R})} \leqslant \sum_{Q \in \mathscr{B}_{M}}\|\overline{\mathrm{v}}(E ; Q)\|_{L_{p}(\mathbb{R})} \stackrel{(92)}{{ }_{\gtrsim}} \stackrel{(1, K}{ } \sum_{Q \in \mathscr{B}_{M}} \sigma(Q)^{2 k+1} \leqslant C \sigma(M)^{2 k+1}, \tag{93}
\end{equation*}
$$

where the final inequality in (93) is the $C$-Carleson condition for $\mathcal{B}$.
Our next goal is to show that

$$
\begin{equation*}
\forall \mathcal{T} \in \mathscr{F}_{M}, \quad\left\|W_{\mathcal{T}}(E)\right\|_{L_{p}(\mathbb{R})} \lesssim_{k, K}(\Lambda+1) \sigma(\mathrm{Q}(\mathcal{T}))^{2 k+1} \tag{94}
\end{equation*}
$$

Let $\Gamma^{+}$be a $\lambda$-Lipschitz half-space such that $d_{N_{4}(Q(\mathcal{T}))}\left(\Gamma^{+}, E\right) \leqslant \theta \sigma(Q(\mathcal{T}))$. Denote $A=\Gamma^{+} \triangle E$. Then $\mathrm{D}_{s} E=\mathrm{D}_{s} A \triangle \mathrm{D}_{s} \Gamma^{+}$, so $W_{\mathcal{J}}(E) \leqslant W_{\mathcal{J}}(A)+$ $W_{\mathcal{J}}\left(\Gamma^{+}\right)$. We shall bound $W_{\mathcal{J}}\left(\Gamma^{+}\right)$using (88). Suppose that $q \in \mathrm{Q}(\mathcal{T})$, and let $R=\operatorname{diam} \mathrm{Q}(\mathcal{T})+4 \sigma(\mathrm{Q}(\mathcal{T}))$. Then $N_{4}(\mathrm{Q}(\mathcal{T})) \subseteq B_{R}(q)$. We claim that for all $s \in \mathbb{R}$, we have $W_{\mathcal{J}}\left(\Gamma^{+}\right)(s) \lesssim_{k, K}{\overline{{ }_{B}^{R}}}(q)\left(\Gamma^{+}\right)(s)$. Indeed, if $i=\lceil s\rceil$, then

$$
\begin{equation*}
W_{\mathcal{T}}\left(\Gamma^{+}\right)(s)=\sum_{Q \in \mathcal{T} \cap \Delta_{i}} 2^{-s} \mathscr{H}^{2 k+2}\left(N_{4}(Q) \cap \mathrm{D}_{s} \Gamma^{+}\right) . \tag{95}
\end{equation*}
$$

The sets $\left\{N_{4}(Q): Q \in \mathcal{T} \cap \Delta_{i}\right\}$ have multiplicity bounded by a quantity that depends only on $k$ and $K$, and they are all contained in $B_{R}(q)$. It therefore follows from (95) that

$$
W_{\mathcal{T}}\left(\Gamma^{+}\right)(s) \lesssim_{k, K} 2^{-s} \mathcal{H}^{2 k+2}\left(B_{R}(q) \cap \mathrm{D}_{s} \Gamma^{+}\right)=\overline{\mathrm{v}}_{B_{R}(q)}\left(\Gamma^{+}\right)(s),
$$

which is the required point-wise estimate. By the assumption of Proposition 55 we deduce that

$$
\begin{equation*}
\left\|W_{\mathcal{T}}\left(\Gamma^{+}\right)\right\|_{L_{p}(\mathbb{R})} \lesssim_{k, K}\left\|\overline{\mathrm{v}}_{B_{R}(q)}\left(\Gamma^{+}\right)\right\|_{L_{p}(\mathbb{R})} \stackrel{(88)}{\lesssim}_{\gtrsim_{k, K}} \Lambda R^{2 k+1} \asymp_{k, K} \Lambda \sigma(\mathrm{Q}(\mathcal{T}))^{2 k+1} . \tag{96}
\end{equation*}
$$

Due to (96), to establish (94) it suffices to show that $\left\|W_{\mathcal{J}}(A)\right\|_{L_{p}(\mathbb{R})} \lesssim_{k, K}$ $\sigma(\mathrm{Q}(\mathcal{T}))^{2 k+1}$. Let $\psi=\frac{1}{20}$ (thus $\theta<\frac{\psi}{12}$ ) and let $\left\{c_{i}\right\}_{i \in I} \subseteq 10 \mathrm{Q}(\mathcal{T})$ and $\left\{4 B_{i}=\right.$ $\left.B_{4 \psi d_{\mathcal{J}}\left(c_{i}\right)}\left(c_{i}\right)\right\}_{i \in I}$ be the covering constructed in Lemma 54. Let $A^{\prime}=A \cap$ $N_{4}(\mathrm{Q}(\mathscr{T}))$. By Lemma 56, $A^{\prime} \subseteq \bigcup_{i \in I} 4 B_{i} \cup \Gamma$.

For $s \in \mathbb{R}$ and $F \subseteq \mathbb{H}^{2 k+1}$, let $\mathrm{U}_{s} F \stackrel{\text { def }}{=} F \cup F Z^{2^{2 s}}$. For every $Q \in \Delta$, define

$$
\overline{\mathrm{v}}^{\prime}(F ; Q)(s) \stackrel{\text { def }}{=} \frac{\mathcal{H}^{2 k+2}\left(N_{4}(Q) \cap \mathrm{U}_{s} F\right)}{2^{s}} \mathbf{1}_{\left\{\frac{\sigma(Q)}{2}<2^{s} \leqslant \sigma(Q)\right\}}
$$

and let $W_{\mathcal{J}}^{\prime}(F) \stackrel{\text { def }}{=} \sum_{Q \in \mathcal{T}} \overline{\mathrm{v}}^{\prime}(F ; Q)$. Then $\mathrm{D}_{s} F \subseteq \mathrm{U}_{s} F$, and therefore we have the point-wise inequality $\overline{\mathrm{v}}^{\prime}(F ; Q) \leqslant \overline{\mathrm{v}}(F ; Q)$. Importantly, $\overline{\mathrm{v}}^{\prime}(\cdot ; Q)$ is subadditive, i.e., if $F, G \subseteq \mathbb{H}^{2 k+1}$, then $\overline{\mathrm{v}}^{\prime}(F \cup G ; Q) \leqslant \overline{\mathrm{v}}^{\prime}(F ; Q)+\overline{\mathrm{v}}^{\prime}(G ; Q)$. Consequently, $W_{\mathcal{J}}\left(A^{\prime}\right) \leqslant \sum_{i \in I} W_{\mathcal{J}}^{\prime}\left(4 B_{i}\right)$.

Let $y=c_{i} \in 10 \mathrm{Q}(\mathcal{T})$ for some $i \in I$, and let $B=4 B_{i}=B_{d_{\mathcal{J}}(y) / 5}(y)$. As $Q$ ranges over $\Delta_{j}$ for some $j \in \mathbb{Z}$, the sets $N_{4}(Q)$ cover $\mathrm{U}_{s} B$ with multiplicity
bounded by a quantity that depends only on $k, K$. Hence, for any $s \in \mathbb{R}$, we have

$$
\begin{equation*}
0 \leqslant W_{\mathcal{T}}^{\prime}(B)(s) \lesssim_{k, K} \frac{\mathscr{H}^{2 k+2}\left(\mathrm{U}_{s} B\right)}{2^{s}} \lesssim \frac{d_{\mathcal{J}}(y)^{2 k+2}}{2^{s}} . \tag{97}
\end{equation*}
$$

We claim that if $W_{\mathcal{J}}^{\prime}(B)(s)>0$, then $2^{s}>d_{\mathcal{T}}(y) / 24$. Indeed, suppose that $Q \in \mathcal{T}$ and $s \in \mathbb{R}$ are such that $\overline{\mathrm{v}}^{\prime}(B ; Q)(s)>0$. Then $\sigma(Q) / 2<2^{s} \leqslant \sigma(Q)$ and $\left(B \cup B Z^{2^{2 s}}\right) \cap N_{4}(Q) \neq \varnothing$. Let $x$ lie in this intersection. Since $x \in N_{4}(Q)$, we have $d_{\mathcal{T}}(x) \leqslant 5 \sigma(Q)$, and since $x \in B \cup B Z^{2^{2 s}}$, we have

$$
d(x, y) \leqslant \frac{d_{\mathcal{J}}(y)}{5}+4 \cdot 2^{s} \leqslant \frac{d_{\mathcal{J}}(y)}{5}+4 \sigma(Q) .
$$

Since $d_{\mathcal{S}}$ is 1-Lipschitz, it follows that

$$
d_{\mathcal{T}}(y) \leqslant d_{\mathcal{J}}(x)+d(x, y) \leqslant \frac{d_{\mathcal{T}}(y)}{5}+9 \sigma(Q) .
$$

Solving for $d_{\mathcal{T}}(y)$, we find that $d_{\mathcal{T}}(y)<12 \sigma(Q)$ and thus $2^{s}>\frac{d_{\mathcal{J}}(y)}{24}$. Consequently,

$$
\begin{aligned}
&\left\|W_{\mathcal{T}}^{\prime}(B)\right\|_{L_{p}(\mathbb{R})}=\left(\int_{\log _{2}}^{\infty} W_{\mathcal{J}(y)}^{24}\right. \\
&\left.W_{\mathcal{T}}^{\prime}(B)(s)^{p} \mathrm{~d} s\right)^{\frac{1}{p}} \\
&{ }^{(97)} \\
& k, K \\
&\left(\int_{\log _{2} \frac{d_{\mathcal{J}}(y)}{24}}^{\infty} \frac{d_{\mathcal{T}}(y)^{p(2 k+2)}}{2^{p s}} \mathrm{~d} s\right)^{\frac{1}{p}} \lesssim d_{\mathcal{J}}(y)^{2 k+1} .
\end{aligned}
$$

This proves that

$$
\left\|W_{\mathcal{J}}\left(A^{\prime}\right)\right\|_{L_{p}(\mathbb{R})} \leqslant \sum_{i \in I}\left\|W_{\mathcal{T}}^{\prime}\left(4 B_{i}\right)\right\|_{L_{p}(\mathbb{R})} \lesssim_{k, K} \sum_{i \in I} d_{\mathcal{J}}\left(c_{i}\right)^{2 k+1}
$$

The sets $\left\{\partial E \cap B_{i}\right\}_{i \in I}$ are disjoint, and for each $i, \mathscr{H}^{2 k+1}\left(\partial E \cap B_{i}\right) \asymp$ $d_{\mathcal{J}}\left(c_{i}\right)^{2 k+1}$. Since $c_{i} \in 10 \mathrm{Q}(\mathcal{T})$, we have $d_{\mathcal{T}}\left(c_{i}\right) \leqslant 11 \sigma(\mathrm{Q}(\mathcal{T}))$. Therefore, $\partial E \cap B_{i} \subseteq 12 \mathrm{Q}(\mathcal{T})$, and

$$
\begin{aligned}
&\left\|W_{\mathcal{J}}\left(A^{\prime}\right)\right\|_{L_{p}(\mathbb{R})} \lesssim k, K \\
& \sum_{i \in I} d_{\mathcal{J}}\left(c_{i}\right)^{2 k+1} \\
& \asymp \mathcal{H}^{2 k+1}\left(\partial E \cap \bigcup_{i \in I} B_{i}\right) \leqslant \mathcal{H}^{2 k+1}(12 \mathrm{Q}(\mathcal{T})) \lesssim \sigma(\mathrm{Q}(\mathcal{T}))^{2 k+1}
\end{aligned}
$$

Now we correct for the difference between $W_{\mathcal{T}}\left(A^{\prime}\right)$ and $W_{\mathcal{T}}(A)$. Suppose that $Q \in \mathcal{T}$ and $s \in \mathbb{R}$ satisfies $\sigma(Q) / 2<2^{s} \leqslant \sigma(Q)$. If in addition $Q \neq$ $\mathrm{Q}(\mathcal{T})$, then $\sigma(Q) \leqslant \sigma(\mathrm{Q}(\mathcal{T})) / 2$, and therefore $N_{4}(Q) \cup N_{4}(Q) Z^{-2^{2 s}} \subseteq N_{8}(Q) \subseteq$ $N_{4}(\mathrm{Q}(\mathcal{T}))$. Consequently, $N_{4}(Q) \cap \mathrm{D}_{s} A=N_{4}(Q) \cap \mathrm{D}_{s} A^{\prime}$. Hence, $\overline{\mathrm{v}}\left(A^{\prime} ; Q\right)=$
$\overline{\mathrm{v}}(Q ; A)$, and therefore $W_{\mathcal{J}}(A)=W_{\mathcal{T}}\left(A^{\prime}\right)+\overline{\mathrm{v}}(A ; \mathrm{Q}(\mathcal{T}))-\overline{\mathrm{v}}\left(A^{\prime} ; \mathrm{Q}(\mathcal{T})\right)$. So,

$$
\begin{align*}
\left\|W_{\mathcal{T}}(A)\right\|_{L_{p}(\mathbb{R})} & =\left\|W_{\mathcal{T}}\left(A^{\prime}\right)+\overline{\mathrm{v}}(A ; \mathrm{Q}(\mathcal{T}))-\overline{\mathrm{v}}\left(A^{\prime} ; \mathrm{Q}(\mathcal{T})\right)\right\|_{2}  \tag{98}\\
& \leqslant\left\|W_{\mathcal{J}}\left(A^{\prime}\right)\right\|_{L_{p}(\mathbb{R})}+2 \frac{\mathcal{H}^{2 k+2}\left(N_{4}(\mathrm{Q}(\mathcal{T}))\right)}{\sigma(\mathrm{Q}(\mathcal{T})) / 2} \lesssim \sigma(\mathrm{Q}(\mathcal{T}))^{2 k+1} .
\end{align*}
$$

Therefore, for any $\mathcal{T} \in \mathcal{F}$,

$$
\begin{aligned}
& \left\|W_{\mathcal{J}}(E)\right\|_{L_{p}(\mathbb{R})} \leqslant\left\|W_{\mathcal{T}}(A)\right\|_{L_{p}(\mathbb{R})}+\left\|W_{\mathcal{J}}\left(\Gamma^{+}\right)\right\|_{L_{p}(\mathbb{R})} \\
& \stackrel{(96) \wedge(98)}{\lesssim k, K}(\Lambda+1) \sigma(\mathrm{Q}(S))^{2 k+1} .
\end{aligned}
$$

This completes the justification of (94).
By combining (94) with the Carleson packing condition on $\mathscr{F}_{M}$, (91), and (93), for any $M \in \Delta$,

$$
\begin{aligned}
\left\|W_{\Delta_{M}}(E)\right\|_{L_{p}(\mathbb{R})} & \leqslant\left\|W_{\mathcal{B}_{M}}(E)\right\|_{L_{p}(\mathbb{R})}+\sum_{\mathcal{T} \in \mathscr{T}_{M}}\left\|W_{\mathcal{T}}(E)\right\|_{L_{p}(\mathbb{R})} \\
& \lesssim k, K \\
& \lesssim \sigma(M)^{2 k+1}+(\Lambda+1) \sum_{\mathcal{T} \in \mathscr{F}_{M}} \sigma(\mathrm{Q}(\mathcal{T}))^{2 k+1} \\
& \lesssim C(\Lambda+1) \sigma(M)^{2 k+1}
\end{aligned}
$$

thus completing the proof of (90).
Finally, we bound $\left\|\bar{v}_{B_{1}}(E)\right\|_{L_{p}(\mathbb{R})}$. Let $M_{1}, \ldots, M_{m} \in \Delta_{0}$ be the set of cubes such that $B_{5} \cap M_{i} \neq \varnothing$ for all $i \in\{1, \ldots, m\}$. The number of such cubes is bounded by a quantity that depends only on $k$ and $K$. If $Q \in \Delta$ and $N_{4}(Q) \cap$ $B_{1} \neq \varnothing$, then $Q \subseteq M_{i}$ for some $i \in\{1, \ldots, m\}$, so for every $s \leqslant 0$, we have $\overline{\mathrm{v}}_{B_{1}}(E)(s) \leqslant \sum_{i=1}^{m} W_{\Delta_{M_{i}}}(E)(s)$. Hence by (90) we have $\left\|\overline{\mathrm{v}}_{B_{1}}(E)\right\|_{L_{p}(-\infty, 0)} \lesssim_{k, K}$ $C(\Lambda+1)$. On the other hand, if $s>0$, then $\overline{\mathrm{v}}_{B_{1}}(E)(s) \leqslant 2^{-s} \mathcal{H}^{2 k+2}\left(B_{1}\right) \lesssim 2^{-s}$, so $\left\|\overline{\mathrm{v}}_{B_{1}}(E)\right\|_{L_{p}(0, \infty)} \lesssim 1$, and we conclude that $\left\|\overline{\mathrm{v}}_{B_{1}}(E)\right\|_{L_{p}(\mathbb{R})} \lesssim_{k, K} C(\Lambda+1)$, as desired.

## 8. Decomposing surfaces into Ahlfors regular pieces

In Section 7, we established Theorem 36 for sets admitting a corona decomposition. In this section and the next section, we will prove it for cellular sets and show that this implies Theorem 36 in full generality. The following theorem will be proven in Section 9.

Theorem 57. Fix $C, r \in(0, \infty)$. Let $E \subseteq \mathbb{H}^{2 k+1}$ be a set such that $\left(E, E^{c}, \partial E\right)$ is $(C, r)$-regular as in Definition 45. Then there exists $K=$ $K(k, C) \in(0, \infty)$ such that the pair $(E, \partial E)$ admits a $(K, r)$-corona decomposition as in Definition 53.

Theorem 57 asserts that local Ahlfors-regularity implies the existence of local intrinsic corona decompositions. In order to apply this structural statement, in this section we shall prove a further structural result that provides a decomposition of a general finite-perimeter cellular set into pieces that are locally Ahlfors-regular.

Lemma 58. There exists $C=C_{k}>0$ such that if $E \subseteq \mathbb{H}^{2 k+1}$ is a finiteperimeter cellular set, then one can find $n \in \mathbb{N},\left\{x_{i}\right\}_{i=0}^{n-1} \subseteq \mathbb{H}^{2 k+1},\left\{s_{i}\right\}_{i=0}^{n-1} \subseteq$ $[1, \infty)$ and a sequence $E_{0}=E, E_{1}, \ldots, E_{n}=\varnothing$ of subsets of $\mathbb{H}^{2 k+1}$ such that if we denote $B_{i}=B_{s_{i}}\left(x_{i}\right)$, then the following assertions hold true:

- $\left(E_{i}, E_{i}^{\subset}, \partial E_{i}\right)$ is $\left(C, s_{i}\right)$-regular for every $i \in\{1, \ldots, n\}$.
- $E_{i} \triangle E_{i+1} \subseteq B_{i}$ for every $i \in\{0, \ldots, n-1\}$.
- $\sum_{i=0}^{n-1} s_{i}^{2 k+1} \lesssim \mathscr{H}^{2 k+1}(\partial E)$.
- For every $p \in[1, \infty)$, we have

$$
\begin{equation*}
\|\overline{\mathrm{v}}(E)\|_{L_{p}(\mathbb{R})} \leqslant C \mathscr{H}^{2 k+1}(\partial E)+\sum_{i=0}^{n-1}\left\|\overline{\mathrm{v}}_{B_{i}}\left(E_{i}\right)\right\|_{L_{p}(\mathbb{R})} . \tag{99}
\end{equation*}
$$

Before proving Lemma 58, we now show that Theorem 36 follows quickly from Theorem 57 and Lemma 58, in combination with the results that we already proved in Sections 3, 5 and 7.

Proof of Theorem 36 assuming Theorem 57 and Lemma 58. In Section 3 we have already shown that Theorem 36 follows from Theorem 39, i.e., its special case for cellular sets. So, suppose that $E \subseteq \mathbb{H}^{2 k+1}$ is a cellular set with $\mathscr{H}^{2 k+1}(\partial E)<\infty$. Let $n \in \mathbb{N},\left\{s_{i}\right\}_{i=0}^{n-1},\left\{E_{i}\right\}_{i=0}^{n},\left\{B_{i}\right\}_{i=0}^{n-1}$ be as in Lemma 58. By Theorem 57 we know that for every $i \in\{1, \ldots, n\}$, the pair $\left(E_{i}, \partial E_{i}\right)$ admits a ( $K, s_{i}$ )-local corona decomposition with $K \lesssim 1$. By combining Proposition 41 and Proposition 55 we see that $\overline{\mathrm{v}}_{B_{i}}\left(E_{i}\right) \lesssim s_{i}^{2 k+1}$ for every $i \in\{0, \ldots, n-1\}$. Hence,

$$
\begin{aligned}
\|\overline{\mathrm{v}}(E)\|_{L_{2}(\mathbb{R})} & \lesssim \mathscr{H}^{2 k+1}(\partial E)+\sum_{i=0}^{n-1}\left\|\overline{\mathrm{v}}_{B_{i}}\left(E_{i}\right)\right\|_{L_{p}(\mathbb{R})} \\
& \lesssim \mathscr{H}^{2 k+1}(\partial E)+\sum_{i=0}^{n-1} s_{i}^{2 k+1} \lesssim \mathscr{H}^{2 k+1}(\partial E)
\end{aligned}
$$

where the first inequality is (99) and the final inequality uses the third assertion of Lemma 58.

The remainder of this section is devoted to the proof of Lemma 58. To do so, we will need some preparatory definitions and lemmas regarding cellular sets. Recall that $E \subseteq \mathbb{H}^{2 k+1}$ is cellular if it is a union of translates of the unit cube $C_{0}=\left[-\frac{1}{2}, \frac{1}{2}\right]^{2 k+1}$ by elements of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$. The set of all translates of $C_{0}$ by elements of $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ gives rise to the unit lattice $\tau^{2 k+1}$. If $C \in \tau^{2 k+1}$, then $C=h C_{0}$ for some $h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}$, which we denote $h_{C}$.

For $x, y, z>0$, let $C_{x, y, z}$ be the cuboid $[-x, x]^{k} \times[-y, y]^{k} \times[-z, z] \subseteq \mathbb{H}^{2 k+1}$. If $h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ and $r>0$ is an integer, we define the cellular ball $B_{r}^{\square}(h)$ to be the set

$$
B_{r}^{\square}(h) \stackrel{\text { def }}{=} h \cdot B_{r}^{\square} \stackrel{\text { def }}{=} h\left(\mathbb{H}_{\mathbb{Z}}^{2 k+1} \cap C_{r, r, r^{2}}\right) C_{0} .
$$

Lemma 59. There exists $c=c_{k} \in(0, \infty)$ such that for all $g \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ and all $r \in \mathbb{N}$ and $p \in(0, \infty)$,

- $B_{r / c}(g) \subseteq B_{r}^{\square}(g) \subseteq B_{c r}(g)$,
- $\mathscr{H}^{2 k+2}\left(B_{r}^{\square}(g)\right) \asymp r^{2 k+2}$,
- $\mathcal{H}^{2 k+1}\left(\partial B_{r}^{\square}(g)\right) \asymp r^{2 k+1}$,
- $\left\|\overline{\mathrm{v}}\left(B_{r}^{\square}(g)\right)\right\|_{L_{p}(\mathbb{R})} \asymp r^{2 k+1}$.

Proof. It suffices to consider the case that $g=\mathbf{0}$. The first and second assertions follow from (23) and the fact that $\mathcal{H}^{2 k+2}\left(B_{r}\right) \asymp r^{2 k+2}$.

If $D \in \tau^{2 k+1}$ is adjacent to the cell $C_{0}$, then there is a $v \in C_{0} \cap D$. Then $v, h_{D}^{-1} v \in C_{0}$, so $h_{D} \in C_{0} \cdot C_{0}^{-1}=C_{0}^{2}$. By (21), $C_{0}^{2} \subseteq C_{1,1,1+k / 4}$. Let $A \subseteq \tau^{2 k+1}$ be the set of $(2 k+1)$-cells contained in $B_{r}^{\square}$ that intersect $\partial B_{r}^{\square}$. Each element of $A$ is adjacent to some cell $D$ such that $h_{D} \notin C_{r, r, r^{2}}$, so if $C \in A$, then $h_{C} \notin C_{r-1, r-1, r^{2}-2 k r}$. Then $A \subseteq C_{r, r, r^{2}} \backslash C_{r-1, r-1, r^{2}-2 k r}$, so $|A| \lesssim r^{2 k+1}$. Since the boundary of $B_{r}^{\square}$ is contained in the boundaries of the cells of $A$, we have $\mathscr{H}^{2 k+1}\left(\partial B_{r}^{\square}(g)\right) \lesssim r^{2 k+1}$. On the other hand, Lemma 20 implies that $\mathcal{H}^{2 k+1}\left(\partial B_{r}^{\square}\right) \gtrsim r^{2 k+1}$. This proves the third assertion.

Finally, to bound the vertical perimeter of $B_{r}^{\square}$, note that $B_{r}^{\square}$ consists of the union of $(2 r+1)^{2 k}$ "stacks" of parallelepipeds, i.e., regions of the form $S_{h, n}=\bigcup_{i=0}^{n} h Z^{i} C_{0}$ for some $h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ and some height $n \lesssim r^{2}$. Each such stack is a single, very tall parallelepiped, so $\overline{\mathrm{v}}\left(S_{h, n}\right)(s)=\frac{\min \left\{2^{2 s}, n\right\}}{2^{s}}$ for every $s \in \mathbb{R}$. It follows that

$$
\overline{\mathrm{v}}\left(B_{r}^{\square}\right)(s) \asymp(2 r+1)^{2 k} \cdot \frac{\min \left\{2^{2 s}, r^{2}\right\}}{2^{s}} .
$$

By integrating these estimates we conclude that $\left\|\overline{\mathrm{v}}\left(B_{r}^{\square}\right)\right\|_{L_{p}(\mathbb{R})} \asymp r^{2 k+1}$.
We will need the following discrete version of Ahlfors regularity.
Definition 60. Fix $C, r, s>0$. A set $A \subseteq \mathbb{H}^{2 k+1}$ is discretely $(C, r)$-Ahlfors $s$-regular if

$$
\frac{\rho^{s}}{C} \leqslant \mathscr{H}^{s}\left(A \cap B_{\rho}^{\square}(x)\right) \leqslant C \rho^{s}
$$

for all $1 \leqslant \rho \leqslant r$ and all $x \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}$ such that $\mathscr{H}^{s}\left(x C_{0} \cap A\right)>0$.
We will use the same convention on dimension as in Definition 45. That is, if $E \subseteq \mathbb{H}^{2 k+1}$ and we say that $E$ is discretely $(C, r)$-regular, then it will be understood that $E$ is discretely $(C, r)$-Ahlfors $s$-regular with $s=2 k+2$. When
we say that $\partial E$ is discretely $(C, r)$-regular we mean that $\partial E$ is assumed to be discretely ( $C, r$ )-Ahlfors $s$-regular with $s=2 k+1$.

By Lemma 59, a ( $C, r$ )-regular set is discretely $\left(C^{\prime}, r\right)$-regular for some $C^{\prime}>0$ depending on $C$ and $k$. Conversely, any union of $s$-cells of $\tau^{2 k+1}$ is locally Ahlfors regular on sufficiently small scales, so if $E$ is cellular and $\left(E, E^{c}, \partial E\right)$ is discretely $(C, r)$-regular, then it is also $\left(C^{\prime}, r\right)$-regular for some $C^{\prime}>0$ depending on $C$ and $k$.

Recall that Lemma 20 is a relative isoperimetric inequality for subsets of balls in $\mathbb{H}^{2 k+1}$. By Lemma 59, Lemma 20 also holds (with different constants) for subsets of cellular balls; i.e., there exists $S=S_{k} \in(1, \infty)$ such that if $r \geqslant 1$ and $h \in \mathbb{H}^{2 k+1}$, then

$$
\begin{equation*}
\left(\frac{\mathscr{H}^{2 k+2}\left(E \cap h B_{r}^{\square}\right) \cdot \mathcal{H}^{2 k+2}\left(E^{\mathrm{c}} \cap h B_{r}^{\square}\right)}{\mathscr{H}^{2 k+2}\left(B_{r}^{\square}\right)^{2}}\right)^{\frac{2 k+1}{2 k+2}} \lesssim \frac{r \mathscr{H}^{2 k+1}\left(\partial E \cap h B_{S r}^{\square}\right)}{\mathcal{H}^{2 k+2}\left(h B_{S r}^{\square}\right)} . \tag{100}
\end{equation*}
$$

Proof of Lemma 58. We will construct the $\left\{E_{i}\right\}_{i=0}^{n}$ inductively, by choosing appropriate cellular balls $H_{i}=B_{r_{i}}^{\square}\left(x_{i}\right)$ with $r_{i} \geqslant 1$ and letting

$$
E_{i+1} \stackrel{\text { def }}{=} \begin{cases}E_{i} \cup H_{i} & \text { if } \mathcal{H}^{2 k+2}\left(E_{i} \cap H_{i}\right)>\frac{\mathcal{H}^{2 k+2}\left(H_{i}\right)}{2},  \tag{101}\\ \overline{E_{i} \backslash H_{i}} & \text { if } \mathscr{H}^{2 k+2}\left(E_{i} \cap H_{i}\right) \leqslant \frac{\mathcal{H}^{2 k+2}\left(H_{i}\right)}{2} .\end{cases}
$$

This ensures that $E_{i}$ will be a cellular set for all $i \in\{1, \ldots, n\}$. By the remark above, $\left(E_{i}, E_{i}^{\mathrm{c}}, \partial E_{i}\right)$ is regular if and only if it is discretely regular.

The most important factor in choosing the cellular ball $H_{i}$ is that the size of $\partial E_{i}$ must decrease quickly as $i$ increases. Specifically, we claim that for any finite-perimeter cellular set $E=E_{0}$, we can choose a cellular ball $H=H_{0}=B_{r_{0}}^{\square}\left(x_{0}\right)=B_{r}^{\square}(x)$ so that if $E_{1}$ is defined as above, then

$$
\begin{equation*}
\mathscr{H}^{2 k+1}(\partial E)-\mathscr{H}^{2 k+1}\left(\partial E_{1}\right)=\mathscr{H}^{2 k+1}(H \cap \partial E)-\mathscr{H}^{2 k+1}\left(H \cap \partial E_{1}\right) \gtrsim r^{2 k+1} . \tag{102}
\end{equation*}
$$

Let $c_{0}>0$ be a large constant to be determined later, and let $S$ be as in (100). Let $r \in \mathbb{N}$ be the largest integer such that $\partial E$ is discretely ( $c_{0}, r-1$ )regular and $E, E^{\mathrm{c}}$ are both discretely $\left(c_{0}^{(2 k+2) /(2 k+1)}, r-1\right)$-regular. Since $E$ is cellular, we can assume $r \geqslant 10 S$ by taking $c_{0}$ to be sufficiently large. Any ball of radius $r$ can be covered by boundedly (depending only on $k$ ) many balls of radius $r-1$, so ( $E, E^{\text {c }}, \partial E$ ) is discretely ( $c_{1}, r$ )-regular, where $c_{1} \asymp$ $c_{0}^{(2 k+2) /(2 k+1)}$.

If $r \geqslant \operatorname{diam} E$, then we can take $H$ to be a ball of radius $2 r$ containing $E$, so that $E_{1}=\varnothing$. Since $\partial E$ is discretely $\left(c_{0}, r-1\right)$-regular, we have $\mathscr{H}^{2 k+1}(\partial E) \gtrsim$ $r^{2 k+1}$ and thus (102) holds. Otherwise, suppose that $r<\operatorname{diam} E$. Then one of the following conditions holds:
(1) There is $x \in \partial E$ such that $\mathscr{H}^{2 k+1}\left(\partial E \cap B_{r}^{\square}(x)\right)>c_{0} r^{2 k+1}$.
(2) There is $x \in \partial E$ such that $\mathscr{H}^{2 k+1}\left(\partial E \cap B_{r}^{\square}(x)\right)<c_{0}^{-1} r^{2 k+1}$.
(3) There is $x \in E$ such that $\mathcal{H}^{2 k+2}\left(E \cap B_{r}^{\square}(x)\right)<c_{0}^{-(2 k+2) /(2 k+1)} r^{2 k+2}$.
(4) There is $x \in E^{\text {c }}$ such that $\mathscr{H}^{2 k+2}\left(E^{\text {c }} \cap B_{r}^{\square}(x)\right)<c_{0}^{-(2 k+2) /(2 k+1)} r^{2 k+2}$.

We will handle these cases one by one.
If Case (1) holds, take $H=B_{r}^{\square}(x)$. So, $\mathscr{H}^{2 k+1}\left(H \cap \partial E_{0}\right)>c_{0} r^{2 k+1}$ and $\mathcal{H}^{2 k+1}\left(H \cap \partial E_{1}\right) \lesssim r^{2 k+1}$. If $c_{0}$ is sufficiently large, then this implies (102).

Next, suppose that Case (2) holds. Then $\mathcal{H}^{2 k+1}\left(\partial E \cap B_{r}^{\square}(x)\right)<c_{0}^{-1} r^{2 k+1}$, so by (100) we have

$$
\begin{aligned}
& \frac{\min \left\{\mathscr{H}^{2 k+2}\left(E \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right), \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right)\right\}}{r^{2 k+2}} \\
& \\
& \lesssim\left(\frac{\mathscr{H}^{2 k+1}\left(\partial E \cap B_{r}^{\square}(x)\right)}{r^{2 k+1}}\right)^{\frac{2 k+2}{2 k+1}}<c_{0}^{-\frac{2 k+2}{2 k+1}} .
\end{aligned}
$$

Hence,
(103) $\min \left\{\mathcal{H}^{2 k+2}\left(E \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right), \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right)\right\} \lesssim c_{0}^{-\frac{2 k+2}{2 k+1}} r^{2 k+2}$.

Suppose that $\mathscr{H}^{2 k+2}\left(E \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right) \leqslant \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right)$. Then by (103),

$$
\begin{equation*}
\mathcal{H}^{2 k+2}\left(E \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right) \lesssim c_{0}^{-\frac{2 k+2}{2 k+1}} r^{2 k+2} \tag{104}
\end{equation*}
$$

By the pigeonhole principle, it follows that there is an integer $r_{*} \in\left[\frac{r}{2 S}, \frac{r}{S}\right]$ such that

$$
\mathscr{H}^{2 k+1}\left(E \cap \partial B_{r_{*}}^{\square}(x)\right) \lesssim c_{0}^{-\frac{2 k+2}{2 k+1}} r^{2 k+1} .
$$

Let $H=B_{r_{*}}^{\square}(x)$; note that $r_{*} \geqslant 5$. Due to (104), if $c_{0}$ is large enough, then $E_{1}=\overline{E \backslash H}$ by (101). By the regularity of $\partial E$, we have

$$
\begin{equation*}
\mathscr{H}^{2 k+1}(H \cap \partial E) \geqslant c_{0}^{-1} r_{*}^{2 k+1} . \tag{105}
\end{equation*}
$$

On the other hand, $H \cap \partial E_{1} \subseteq E \cap \partial H$, so

$$
\begin{equation*}
\mathscr{H}^{2 k+1}\left(H \cap \partial E_{1}\right) \lesssim c_{0}^{-\frac{2 k+2}{2 k+1}} r^{2 k+1} \leqslant c_{0}^{-\frac{2 k+2}{2 k+1}} r_{*}^{2 k+1} \tag{106}
\end{equation*}
$$

The implicit constants above are independent of $c_{0}$, so if $c_{0}$ is sufficiently large, then

$$
\mathscr{H}^{2 k+1}(H \cap \partial E)-\mathscr{H}^{2 k+1}\left(H \cap \partial E_{1}\right) \gtrsim c_{0}^{-1} r_{*}^{2 k+1},
$$

as desired.
We argue similarly when $\mathscr{H}^{2 k+2}\left(E \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right) \geqslant \mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right)$. In this case, by (103),

$$
\mathscr{H}^{2 k+2}\left(E^{\mathrm{c}} \cap B_{\lfloor r / S\rfloor}^{\square}(x)\right) \lesssim c_{0}^{-\frac{2 k+2}{2 k+1}} r^{2 k+2} .
$$

Choose $r_{*} \in\left[\frac{r}{2 S}, \frac{r}{S}\right]$ such that

$$
\mathscr{H}^{2 k+1}\left(\overline{E^{\mathrm{c}}} \cap \partial B_{r^{*}}^{\square}(x)\right) \lesssim c_{0}^{\frac{2 k+2}{2 k+1}} r^{2 k+1},
$$

and let $H=B_{r^{\prime}}^{\square}(x)$. Again, $r^{*} \geqslant 5$ and if $c_{0}$ is large enough, then by (101) we have $E_{1}=E \cup H$. Then (105) holds as before and $H \cap \partial E_{1} \subseteq \overline{E^{c}} \cap \partial H$, so (106) holds as above. If $c_{0}$ is sufficiently large, then (106) implies (102).

Finally, Cases (3) and (4) are symmetric, so we consider only the case that

$$
\begin{equation*}
\mathscr{H}^{2 k+2}\left(E \cap B_{r}^{\square}(x)\right)<c_{0}^{-\frac{2 k+2}{2 k+1}} r^{2 k+2} . \tag{107}
\end{equation*}
$$

Hence there exists $\rho \asymp c_{0}^{-1 /(2 k+1)} r$ for which

$$
\mathscr{H}^{2 k+2}\left(B_{\rho}^{\square}(x)\right)>\mathscr{H}^{2 k+2}\left(E \cap B_{r}^{\square}(x)\right) .
$$

This implies that $B_{\rho}^{\square}(x)$ intersects $\partial E$, say at $y \in B_{\rho}^{\square}(x) \cap \partial E$. If $c_{0}$ is sufficiently large, then $\rho<r / 4$, so $B_{r / 4}^{\square}(y) \subseteq B_{r / 2}^{\square}(x)$. By the $(r-1)$-local regularity of $\partial E$,

$$
\begin{equation*}
\mathcal{H}^{2 k+1}\left(B_{\frac{r}{2}}^{\square}(x) \cap \partial E\right) \geqslant \mathscr{H}^{2 k+1}\left(B_{\frac{r}{4}}^{\square}(y) \cap \partial E\right) \geqslant c_{0}^{-1}\left(\frac{r}{4}\right)^{2 k+1} . \tag{108}
\end{equation*}
$$

On the other hand, by the coarea formula, by (107) there is $r^{*} \in\left[\frac{r}{2}, r\right]$ such that if $H=B_{r}^{\square}(x)$, then

$$
\begin{equation*}
\mathscr{H}^{2 k+1}\left(H \cap \partial E_{1}\right) \lesssim c_{0}^{-\frac{2 k+2}{2 k+1}} r^{2 k+1} . \tag{109}
\end{equation*}
$$

Comparing (108) and (109), we see that (102) holds when $c_{0}$ is sufficiently large.

Applying the above procedure inductively, we obtain $\left\{E_{i}\right\}_{i=0}^{n}$ and $\left\{H_{i}=\right.$ $\left.B_{r_{i}}^{\square}\left(x_{i}\right)\right\}_{i=0}^{n-1}$ such that

$$
\begin{equation*}
\forall i \in\{0, \ldots, n-1\}, \quad \mathscr{H}^{2 k+1}\left(\partial E_{i}\right)-\mathscr{H}^{2 k+1}\left(\partial E_{i+1}\right) \gtrsim r_{i}^{2 k+1} . \tag{110}
\end{equation*}
$$

Let $s_{i}=\operatorname{diam} H_{i} \asymp r_{i} \gtrsim 1$ and $B_{i}=B_{s_{i}}\left(x_{i}\right)$ for all $i \in\{0, \ldots, n-1\}$. We claim that this sequence satisfies the conclusions of Lemma 58. Firstly, each step of this iteration decreases $\mathscr{H}^{2 k+1}\left(\partial E_{i}\right)$ by a definite amount, so the construction eventually halts. By (110), we have

$$
\mathscr{H}^{2 k+1}(\partial E)=\sum_{i=0}^{n-1}\left(\mathscr{H}^{2 k+1}\left(\partial E_{i}\right)-\mathscr{H}^{2 k+1}\left(\partial E_{i+1}\right)\right) \gtrsim \sum_{i=0}^{n-1} r_{i}^{2 k+1} \asymp \sum_{i=0}^{n-1} s_{i}^{2 k+1} .
$$

By construction, there exists $C=C_{k} \in(0, \infty)$ such that $\left(E_{i}, E_{i}^{c}, \partial E_{i}\right)$ is $\left(C, s_{i}\right)-$ regular and $E_{i} \triangle E_{i+1} \subseteq B_{i}$ for all $i \in\{0, \ldots, n-1\}$.

Finally, up to a set of measure zero, for every $i \in\{0, \ldots, n-1\}$, either $E_{i} \triangle E_{i+1}=H_{i} \cap E_{i}^{c}$ or $E_{i} \triangle E_{i+1}=H_{i} \cap E_{i}$. Since $\overline{\mathrm{v}}_{H_{i}}\left(E_{i}\right)=\overline{\mathrm{v}}_{H_{i}}\left(E_{i}^{c}\right)$,

Lemma 24 implies that in both cases we have

$$
\begin{aligned}
\left|\overline{\mathrm{v}}\left(E_{i}\right)(s)-\overline{\mathrm{v}}\left(E_{i+1}\right)(s)\right| & \leqslant \overline{\mathrm{v}}\left(E_{i} \triangle E_{i+1}\right)(s) \\
& \leqslant \overline{\mathrm{v}}_{H_{i}}\left(E_{i}\right)(s)+\overline{\mathrm{v}}\left(H_{i}\right)(s) \leqslant \overline{\mathrm{v}}_{B_{i}}\left(E_{i}\right)(s)+\overline{\mathrm{v}}\left(H_{i}\right)(s)
\end{aligned}
$$

for all $s \in \mathbb{R}$. Consequently,

$$
\forall s \in \mathbb{R}, \quad|\overline{\mathrm{v}}(E)(s)| \leqslant \sum_{i=0}^{n-1}\left(\overline{\mathrm{v}}_{B_{i}}\left(E_{i}\right)(s)+\overline{\mathrm{v}}\left(H_{i}\right)(s)\right)
$$

Using Lemma 59, we therefore conclude that there exists $C^{\prime} \in(0, \infty)$ such that

$$
\begin{aligned}
\|\overline{\mathrm{v}}(E)\|_{L_{p}(\mathbb{R})} & \leqslant \sum_{i=0}^{n-1}\left(\left\|{\overline{v_{B_{i}}}}\left(E_{i}\right)\right\|_{L_{p}(\mathbb{R})}+C^{\prime} s_{i}^{2 k+1}\right) \\
& \leqslant C^{\prime} \mathscr{H}^{2 k+1}(\partial E)+\sum_{i=0}^{n-1}\left\|\overline{\mathrm{v}}_{B_{i}}\left(E_{i}\right)\right\|_{L_{p}(\mathbb{R})} .
\end{aligned}
$$

## 9. Constructing corona decompositions

In this section, we prove Theorem 57. Recall that we are given $C, r \in$ $(0, \infty)$ and $E \subseteq \mathbb{R}^{2 k+1}$ such that $\left(E, E^{c}, \partial E\right)$ is $(C, r)$-regular. Our goal is, given $\lambda, \theta \in(0, \infty)$, to construct a $(K, \gamma, \lambda, \theta, r)$-corona decomposition for the pair $(E, \partial E)$, where $K$ is allowed to depend only on $k$ and $C$ and the Carleson packing constant $\gamma$ is allowed to depend only on $C, k, \lambda, \theta$.

To start, as we recalled in Section 6 , since $\partial E$ is $(C, r)$-regular, we may apply the classical construction of Christ cubes [27], [32] to obtain a ( $K, r$ )-cubical patchwork $\left\{\Delta_{i}\right\}_{i \leqslant\left\lfloor\log _{2} r\right\rfloor}$ of $\partial E$, where $K=K(C, k)$. The cubical patchwork $\Delta$ will be used below to build the desired coronization.

We will first give a brief outline of the argument, proceed to introduce the key concept of monotone and $\delta$-monotone sets, and then prove Theorem 57. As in Section 6, for every $Q \in \Delta$, we define $N_{\rho}(Q)=\operatorname{nbhd}_{\rho \sigma(Q)}(Q)$ and $\rho Q=\partial E \cap N_{\rho}(Q)$. The proof follows these steps:
(1) We show that for every $\varepsilon>0$, there are sets $\mathscr{G}_{0}, \mathscr{B}_{0} \subseteq \Delta$ such that $\Delta=$ $\mathcal{L}_{0} \cup \mathscr{B}_{0}, \mathscr{B}_{0}$ is Carleson, and for every $Q \in \mathcal{L}_{0}$, there is a half-space $P_{Q}^{+}$ with bounding plane $P_{Q}$ such that

$$
d_{N_{\frac{1}{\varepsilon}}(Q)}\left(P_{Q}^{+}, E\right) \leqslant \varepsilon \sigma(Q) .
$$

To prove this, we adapt the results on the stability of monotone sets proved in [26]. In that paper, quantitative stability results were proved for monotone sets in $\mathbb{H}^{3}$. Similar techniques would lead to a corresponding stability statement for $\mathbb{H}^{2 k+1}$, but to keep this paper self-contained we will give a much simpler nonquantitative proof of the stability of monotone sets in $\mathbb{H}^{2 k+1}$ that suffices for the present purposes.
(2) We improve the previous result by showing that for all $\varepsilon$, there are sets $\mathscr{\mathcal { L }}, \mathscr{B}$ such that $\mathcal{B}$ is Carleson and such that for all $Q \in \mathscr{L}$, there is a vertical half-space $V_{Q}^{+}$with bounding plane $V_{Q}$ such that

$$
d_{N_{\frac{1}{\varepsilon}}(Q)}\left(V_{Q}^{+}, E\right) \leqslant \varepsilon \sigma(Q)
$$

This step uses the fact that a horizontal plane in $\mathbb{H}^{2 k+1}$ is not homogeneous. A horizontal plane is the union of the horizontal lines through a point, so it has a natural center. The intersection of a horizontal plane with a ball far from its center is close to a vertical plane, so if $E$ is close to a horizontal plane $P$ on a ball, then $E$ is close to a vertical plane on most smaller balls.
(3) For any $\eta>0$, we construct a coronization $(\mathscr{B}, \mathcal{L}, \mathscr{F})$ of $\partial E$ based on the sets of good and bad cubes from the previous step. This coronization has the property that if $\mathcal{\delta} \in \mathcal{F}$ is a stopping-time region and $Q(\mathcal{S})$ is the maximal cube of $\mathcal{S}$, then all $Q \in \mathcal{F}$ satisfy the angle estimate $\angle\left(V_{Q}, V_{Q(\mathcal{S})}\right) \leqslant \eta$. As in Chapter 7 of [30], we construct the stopping-time regions by starting with a maximal cube, then adding descendants of that cube as long as they satisfy certain conditions. The main difficulty is proving that this construction does not produce too many stopping-time regions, i.e., that the maximal cubes of these regions satisfy a Carleson packing condition.
(4) Finally, we show that this coronization is in fact a corona decomposition by showing that for all $\mathcal{S} \in \mathcal{F}$, there is an intrinsic Lipschitz graph satisfying Definition 53. We construct this graph by using a partition of unity to glue together the different planes $V_{Q}$ as $Q$ ranges over $\mathcal{S}$. This is similar to the argument in Chapter 8 of [30], with some complications that arise from the Heisenberg group setting.
9.1. Monotone and $\delta$-monotone sets. One of the main tools in this section is the notion of nonmonotonicity used in [26]. We will review the definition of nonmonotonicity and state some results relating the nonmonotonicity of $E$ to the intersections of $E$ with horizontal lines. Note, however, that our versions of $w_{j}$ and $\widehat{w}_{j}$ count intervals of length between $2^{j-1}$ and $2^{j}$, while the corresponding measures in [26] count intervals of length between $\delta^{j+1}$ and $\delta^{j}$ for some $0<\delta<1$.

Definition 61 (monotone sets). A measurable subset $K \subseteq \mathbb{R}$ is said to be monotone if the characteristic function $\mathbf{1}_{K}$ is monotone up to a set of measure zero. Equivalently, up to a set of measure zero, $K$ and $\mathbb{R} \backslash K$ are both intervals. A subset $F \subseteq \mathbb{R}^{n}$ is said to be monotone if for almost every line $L \subseteq \mathbb{R}^{n}$, the intersection $F \cap L$ is monotone.

We say that a measurable set $E \subseteq \mathbb{H}^{2 k+1}$ is monotone if for almost every horizontal line $L$, the intersection $E \cap L$ is a monotone subset of $L$. We say that $E$ is monotone on an open subset $U \subseteq \mathbb{H}^{2 k+1}$ if for almost every horizontal
line $L$, the intersections $E \cap L \cap U$ and $E^{\mathrm{c}} \cap L \cap U$ are both, up to a set of measure zero, the intersections of $L \cap U$ with intervals.

We will see in the next section that monotone subsets of $\mathbb{R}^{n}$ and $\mathbb{H}^{2 k+1}$ are, up to a set of measure zero, half-spaces; this generalizes a result that Cheeger and Kleiner proved [24] for $\mathbb{H}^{3}$.

Let $\mathcal{L}$ be the space of horizontal lines in $\mathbb{H}^{2 k+1}$, and let $\mathcal{L}\left(B_{r}(x)\right)$ be the set of horizontal lines that intersect $B_{r}(x)$. Let $n$ be the unique (up to constants) measure on $\mathcal{L}$ that is invariant under the action of the isometry group. We normalize $n$ so that $n\left(\mathcal{L}\left(B_{r}(x)\right)\right)=r^{2 k+1}$.

If $U \subseteq \mathbb{H}^{2 k+1}$ and $L \in \mathcal{L}$, define

$$
\mathrm{NM}_{U}(E, L)=\inf \left\{\int_{L \cap U}\left|\mathbf{1}_{I}-\mathbf{1}_{E}\right| \mathrm{d} \mathscr{H}_{L}^{1}: \quad I \text { is a monotone subset of } L\right\} .
$$

The nonmonotonicity of $E$ on $B_{r}(x)$ is defined by

$$
\begin{equation*}
\mathrm{NM}_{B_{r}(x)}(E)=\frac{1}{r^{2 k+2}} \int_{\mathscr{L}\left(B_{r}(x)\right)} \mathrm{NM}_{B_{r}(x)}(E, L) \mathrm{d} \eta(L) . \tag{111}
\end{equation*}
$$

The normalization by $r^{2 k+2}$ makes $\mathrm{NM}_{B_{r}(x)}(E)$ be scale-invariant; i.e., we have

$$
\operatorname{NM}_{B_{t r}(\mathbf{0})}\left(\mathfrak{s}_{t}(E)\right)=\operatorname{NM}_{B_{r}(\mathbf{0})}(E)
$$

for all $t, r>0$. A set $E \subseteq B_{r}(x)$ is said to be $\delta$-monotone on $B_{r}(x)$ if $\mathrm{NM}_{B_{r}(x)}(E)<\delta$.

One can express the nonmonotonicity of $E$ using the kinematic formula. This formula writes the perimeter of a subset of $\mathbb{H}^{2 k+1}$ as an integral over the space $\mathcal{L}$ of horizontal lines in $\mathbb{H}^{2 k+1}$. Recall that, by Lemma 22, if $\mathscr{H}^{2 k+1}(\partial E)<\infty$, then $E$ has finite perimeter and $\operatorname{Per}_{E}(U) \lesssim \mathscr{H}^{2 k+1}(U \cap \partial E)$ for any open set $U \subseteq \mathbb{H}^{2 k+1}$. If $E$ is a set with finite perimeter, then for almost every horizontal line $L$, the intersection $F=L \cap E$ is a subset of $L$ with finite perimeter. Hence there exists a unique collection of disjoint closed intervals $\mathscr{C}(F)=\left\{I_{1}(F), I_{2}(F), \ldots\right\}$ such that $F \triangle \bigcup \mathscr{G}(F)$ has measure zero. The boundary $\partial \bigcup \mathscr{(}(F)$ is a discrete set, and the perimeter measure $\operatorname{Per}(F)$ is the counting measure on $\partial \bigcup \mathscr{C}(F)$; see [2] or [26] for these basic statements. This description leads to the following kinematic formula; see [88] or equation (6.1) in [26]. There is $c>0$ such that for any finite-perimeter set $E \subseteq \mathbb{H}^{2 k+1}$ and any open subset $U \subseteq \mathbb{H}^{2 k+1}$,

$$
\begin{equation*}
\operatorname{Per}(E)(U)=c \int_{\mathcal{L}} \operatorname{Per}(E \cap L)(U \cap L) \mathrm{d} \eta(L) . \tag{112}
\end{equation*}
$$

In [26] the perimeter measure is further decomposed as follows into measures based on the lengths of the intervals in $\mathscr{I}(E \cap L)$. For every $L \in \mathcal{L}$ and
$j \in \mathbb{Z}$, let

$$
C_{j}(E, L) \stackrel{\text { def }}{=}\left\{I \in \mathscr{I}(E \cap L): 2^{j-1} \leqslant \text { length }(I)<2^{j}\right\}
$$

and

$$
C_{\infty}(E, L)=\{I \in \mathscr{I}(E \cap L): \text { length }(I)=\infty\}
$$

For every $j \in \mathbb{Z} \cup\{\infty\}$, let $\varepsilon_{j}(E, L)$ denote the set of endpoints of the intervals in $C_{j}(E, L)$. Let $c$ be as in (112), and let $w_{j}(E, L)$ be $c$ times the counting measure on $\mathscr{E}_{j}(E, L)$. If we let

$$
w_{j}(E)(A) \stackrel{\text { def }}{=} \int_{\mathcal{L}} w_{j}(E, L)(A) \mathrm{d} \eta(L)
$$

and

$$
\widehat{w}_{j}(E)(A) \stackrel{\text { def }}{=} \frac{w_{j}(E)(A)+w_{j}\left(E^{\mathrm{c}}\right)(A)}{2},
$$

then

$$
\begin{equation*}
\operatorname{Per}(E)=\widehat{w}_{\infty}(E)+\sum_{j \in \mathbb{Z}} \widehat{w}_{j}(E) \tag{113}
\end{equation*}
$$

Furthermore, we can use the measures $\left\{\widehat{w}_{j}\right\}_{j \in \mathbb{Z} \cup\{\infty\}}$ to bound the nonmonotonicity of $E$. If $L \in \mathcal{L}$ and the intersection $E \cap L$ has finite perimeter, then after changing $E$ on a measure-zero subset, we can partition $L$ into intervals $J_{1}, \ldots, J_{n}$, arranged in ascending order and each with positive length, so that the $J_{i}$ 's are alternately contained in $E$ and disjoint from $E$. Suppose that $J_{1}, J_{3}, \ldots \subseteq E$. Then $J_{2}, \ldots, J_{n-1}$ all have finite length, and there is an interval $J=J_{1}$ or $J=L$ (depending on whether $n$ is even or odd) such that

$$
\begin{aligned}
\mathrm{NM}_{\mathbb{H}^{2 k+1}}(E, L) & \leqslant \int_{L \cap U}\left|\mathbf{1}_{J}-\mathbf{1}_{E}\right| \mathrm{d} \mathscr{H}_{L}^{1} \\
& \leqslant \sum_{i=2}^{n-1} \operatorname{length}\left(J_{i}\right) \asymp \sum_{j \in \mathbb{Z}} 2^{j} \cdot \widehat{w}_{j}(E, L)\left(\mathbb{H}^{2 k+1}\right) .
\end{aligned}
$$

(If $J_{2}, J_{4}, \ldots \subseteq E$, then the same inequality holds for either $J=J_{n}$ or $J=\varnothing$.)
A similar inequality holds for local nonmonotonicity. Let $r>0, x \in$ $\mathbb{H}^{2 k+1}, a=\inf \left(B_{r}(x) \cap L\right)$, and $b=\sup \left(B_{r}(x) \cap L\right)$. If $m \leqslant M$ are such that $a \in J_{m}, b \in J_{M}$, then

$$
\mathrm{NM}_{B_{r}(x)}(E, L) \leqslant \sum_{i=m+1}^{M-1} \operatorname{length}\left(J_{i}\right) .
$$

For each of these intervals, we have length $J_{i} \leqslant b-a \leqslant 2 r$, so

$$
\begin{equation*}
\mathrm{NM}_{B_{r}(x)}(E, L) \lesssim \sum_{i=-\infty}^{\left\lceil 1+\log _{2} r\right\rceil} 2^{i} \cdot \widehat{w}_{i}(E, L)(U) \tag{114}
\end{equation*}
$$

The following proposition is a simple consequence of the above discussion; see [26, Prop. 4.5].

Proposition 62. Fix $C, r>0$. Suppose that $E \subseteq \mathbb{H}^{2 k+1}$ and that $\partial E$ is $(C, r)$-regular. For every $x \in \mathbb{H}^{2 k+1}$, if we denote $B=B_{r}(x)$, then

$$
\begin{equation*}
\mathrm{NM}_{B}(E) \lesssim C \sum_{n=-\infty}^{\left\lceil 1+\log _{2} r\right\rceil} \frac{2^{n}}{r} \cdot \frac{\widehat{w}_{n}(E)(B)}{\operatorname{Per}(E)(B)} \tag{115}
\end{equation*}
$$

Proof. Integrating (114) and using the regularity of $\partial E$, we find that

$$
\mathrm{NM}_{B}(E) \lesssim \frac{\sum_{n=-\infty}^{\left\lceil 1+\log _{2} r\right\rceil} 2^{n} \cdot \widehat{w}_{n}(E)(B)}{\mathcal{H}^{2 k+2}(B)} \lesssim C \sum_{n=-\infty}^{\left\lceil 1+\log _{2} r\right\rceil} \frac{2^{n}}{r} \cdot \frac{\widehat{w}_{n}(E)(B)}{\operatorname{Per}(E)(B)}
$$

Due to (113), we have $\frac{\widehat{w}_{n}(E)(B)}{\operatorname{Per}(E)(B)} \leqslant 1$ for every $n \in \mathbb{Z}$, so the sum in (115) is typically dominated by the terms where $2^{n}$ is close to $r$.
9.2. $\delta$-monotone sets are close to planes. By [24], every monotone subset of $\mathbb{H}^{3}$ is, up to a set of measure zero, a half-space bounded by a horizontal or vertical plane. This statement was quantified in [26], showing that $\delta$-monotone sets are quantitatively close to half-spaces. Specifically,

Theorem 63 (Theorem 4.13 of [26]). There exists $a>0$ such that for all $\varepsilon>0$, if $E \subseteq B_{1} \subseteq \mathbb{H}^{3}$ is measurable and $\mathrm{NM}_{B_{1}}(E)<\varepsilon^{a}$, then there is a half-space $P^{+} \subseteq \mathbb{H}^{3}$ such that

$$
\frac{\mathscr{H}^{4}\left(\left(E \triangle P^{+}\right) \cap B_{\varepsilon^{3}}\right)}{\mathscr{H}^{4}\left(B_{\varepsilon^{3}}\right)} \lesssim \varepsilon .
$$

While we were motivated by this result, we will not use it in our proof. Instead, we will use a higher-dimensional qualitative version (Proposition 66 below) that states that as $\delta$ goes to $0, \delta$-monotone sets approach half-spaces. We will not need a quantitative estimate on the rate of convergence, so we can avoid the careful estimates necessary to prove Theorem 63.

We will first classify monotone subsets of $\mathbb{R}^{n}$ and $\mathbb{H}^{2 k+1}$. Then, using a compactness result we will conclude that $\delta$-monotone sets are close to monotone sets. Starting with the Euclidean setting, the following lemma asserts that monotone subsets of $\mathbb{R}^{n}$ are half-spaces.

Lemma 64. If $F \subseteq \mathbb{R}^{n}$ is measurable and monotone, either $\mathscr{H}^{n}(F)=0$, $\mathscr{H}^{n}\left(\mathbb{R}^{n} \backslash F\right)=0$, or there is a half-space $P^{+} \subseteq \mathbb{R}^{n}$ such that $\mathcal{H}^{n}\left(F \triangle P^{+}\right)=0$.

Proof. Breaking (locally) from notational conventions in the rest of this paper, in the ensuing proof of Lemma 64 all the distances and balls are with respect to the Euclidean metric on $\mathbb{R}^{n}$.

Suppose that $x_{1}, x_{2} \in \mathbb{R}^{n}$ are density points of $F$ and that $y$ is on the line segment between them. We claim that $y$ is a density point of $F$. Let $\varepsilon>0$, and let $0<r<\min \left\{d\left(x_{1}, y\right), d\left(x_{2}, y\right)\right\}$ satisfy $\mathscr{H}^{n}\left(B_{r}\left(x_{i}\right) \cap F\right)>(1-\varepsilon) \mathscr{H}^{n}\left(B_{r}\left(x_{i}\right)\right)$.

Let $\ell_{0}$ be a line passing through $B_{r / 3}\left(x_{1}\right)$ and $B_{r / 3}\left(x_{2}\right)$ such that for almost every line $\ell$ parallel to $\ell_{0}$, the set $F \cap \ell$ is monotone.

Let $A$ be the set of lines parallel to $\ell_{0}$ that pass through $B_{r / 3}(y)$. Let $\mu$ be the Lebesgue measure on $A$, normalized so that $\mu(A)=1$. Since $d\left(\ell_{0}, y\right) \leqslant r / 3$, every line $\ell \in A$ intersects $B_{2 r / 3}\left(x_{i}\right)$ for $i=1,2$. It follows that every $\ell \in A$ intersects $B_{r}\left(x_{i}\right)$ in an interval of length at least $2 r / 3$. Let

$$
\begin{aligned}
& X \stackrel{\text { def }}{=}\left\{\ell \in A: F \cap \ell \text { is monotone, } \mathscr{H}^{1}\left(\ell \cap F \cap B_{r}\left(x_{1}\right)\right)>0\right. \\
&\text { and } \left.\mathscr{H}^{1}\left(\ell \cap F \cap B_{r}\left(x_{2}\right)\right)>0\right\} .
\end{aligned}
$$

If $\ell \in X$, then the interval $\ell \cap B_{r / 3}(y)$ lies between $\ell \cap B_{r}\left(x_{1}\right)$ and $\ell \cap B_{r}\left(x_{1}\right)$, so the monotonicity of $F \cap \ell$ implies that $\mathscr{H}^{1}\left(\ell \cap B_{r}(y) \cap F^{c}\right)=0$. On the other hand, if $\ell \notin X$, then either $\ell \cap B_{r}\left(x_{1}\right)$ or $\ell \cap B_{r}\left(x_{2}\right)$ is an interval disjoint from $F$, and therefore

$$
\mathscr{H}^{1}\left(\ell \cap F^{\mathrm{c}} \cap\left(B_{r}\left(x_{1}\right) \cup B_{r}\left(x_{2}\right)\right)\right) \geqslant \frac{2 r}{3} .
$$

By integrating over $A$, we find that

$$
\mathscr{H}^{2 k+2}\left(B_{r}\left(x_{1}\right) \cap F^{c}\right)+\mathscr{H}^{2 k+2}\left(B_{r}\left(x_{2}\right) \cap F^{c}\right) \gtrsim \frac{2 r}{3} \mu(A \backslash X) r^{n-1},
$$

so $\mu(A \backslash X) \lesssim \varepsilon$. It follows that

$$
\frac{\mathscr{H}^{n}\left(B_{r / 3}(y) \cap F^{\mathrm{c}}\right)}{\mathscr{H}^{n}\left(B_{r / 3}(y)\right)} \lesssim \mu(A \backslash X) \lesssim \varepsilon
$$

We have shown that the set of density points of $F$ is convex. By symmetry, the set of density points of $F^{c}$ is convex. Since these sets are disjoint and their union is all of $\mathbb{R}^{n}$ except for a set of measure zero, the sets of density points of $F$ and of $F^{c}$ are either empty, $\mathbb{R}^{n}$, or half-spaces.

Theorem 5.1 of [24] asserts that monotone subsets of $\mathbb{H}^{3}$ are half-spaces. By combining, this with Lemma 64, we will show that the monotone subsets of $\mathbb{H}^{2 k+1}$ are half-spaces.

Proposition 65. If $F \subseteq \mathbb{H}^{2 k+1}$ is monotone, then, up to a set of measure zero, we have $F=\varnothing, F=\mathbb{H}^{2 k+1}$, or $F$ is a half-space bounded by a plane.

Proof. Recall that multiplication on $\mathbb{H}^{2 k+1}$ is based on a symplectic form $\omega$ on $\mathbb{R}^{2 k}$. We identify $\mathbb{R}^{2 k}$ with $\mathbb{C}^{k}$ so that $\omega(z, w)=\operatorname{Im}\left(\sum_{i=1}^{k} \overline{z_{i}} w_{i}\right)$. Let $\pi: \mathbb{H}^{2 k+1} \rightarrow \mathbb{C}^{2 k}$ be the abelianization map. We claim that up to measure zero, $F$ is a half-space of $\mathbb{H}^{2 k+1}$. (Recall that half-spaces of $\mathbb{H}^{2 k+1}$ coincide with half-spaces of $\mathbb{R}^{2 k+1}$.)

If $v \in \mathbb{C}^{k} \backslash\{\mathbf{0}\}$, let $\mathbb{H}_{v}=\pi^{-1}(\mathbb{C} v)$. This is a subgroup of $\mathbb{H}^{2 k+1}$ that is isometrically isomorphic to $\mathbb{H}^{3}$. If $h \in \mathbb{H}^{2 k+1}$ and $v \in \mathbb{C}^{k}$ is a unit vector, let $\mathbb{H}(h, v)=h \mathbb{H}_{v}$. By the left-invariance of the Carnot-Carathéodory metric,
each such coset is also isometric to $\mathbb{H}^{3}$. Let $\mathcal{S}$ be the set of such cosets. Any horizontal line $L \in \mathcal{L}$ is contained in a unique element of $\mathcal{S}$, which can be constructed by letting $v$ be parallel to the line $\pi(L)$. The set $\mathcal{L}$ of horizontal lines in $\mathbb{H}^{2 k+1}$ thus fibers over $\mathcal{S}$, and the fiber is the set $\mathcal{L}_{1}$ of horizontal lines in $\mathbb{H}^{3}$.

By Fubini's theorem, if $F \subseteq \mathbb{H}^{2 k+1}$ is a monotone set, then for almost every $\mathbb{H}(h, v) \in \mathcal{S}$, the intersection $F \cap \mathbb{H}(h, v)$ is monotone. If $F \cap \mathbb{H}(h, v)$ is monotone, then by $[24, \mathrm{Th} .5 .1]$, up to a measure zero set, $F \cap \mathbb{H}(h, v) \in\{\varnothing, \mathbb{H}(h, v)\}$, or $F \cap \mathbb{H}(h, v)$ is a half-space in $\mathbb{H}(h, v)$. In each case, $\mathbb{H}(h, v)$ is a 3 -plane in $\mathbb{H}^{2 k+1}$ and $F \cap \mathbb{H}(h, v)$ is monotone not just along horizontal lines but along all lines in $\mathbb{H}(h, v)$.

For almost every line $\ell \subseteq \mathbb{H}^{2 k+1}$ (not necessarily horizontal), the projection $\pi(\ell)$ is a line in $\mathbb{C}^{k}$, and if $v \in \mathbb{C}^{k}$ is parallel to $\pi(\ell)$, then $\ell \subseteq \mathbb{H}(h, v)$ for some $h \in \mathbb{H}^{2 k+1}$. For almost every such $\ell$, the intersection $F \cap \mathbb{H}(h, v)$ is monotone, so $F \cap \ell$ is monotone. By Lemma 64, $F$ is, up to a measure-zero set, either empty, all of $\mathbb{H}^{2 k+1}$, or a half-space, as desired.

Proposition 66. For any $\varepsilon, \rho, c>0$, there exists $\delta>0$ such that if $F \subseteq \mathbb{H}^{2 k+1}$ is measurable, $\partial F$ is $(c, \rho)$-regular, and $\mathrm{NM}_{B_{\rho}}(F)<\delta^{2 k+3}$, then there is a half-space $P^{+} \subseteq \mathbb{H}^{2 k+1}$ such that

$$
\frac{\mathcal{H}^{2 k+2}\left(\left(F \triangle P^{+}\right) \cap B_{\delta_{\rho}}\right)}{\mathscr{H}^{2 k+2}\left(B_{\delta \rho}\right)}<\varepsilon .
$$

Proof. After rescaling, we may suppose that $\rho=1$. If the assertion of Proposition 66 does not hold true, then there is a sequence $\left\{F_{i}\right\}_{i=1}^{\infty}$ of subsets of $\mathbb{H}^{2 k+1}$ such that $\partial F_{i}$ is $(c, \rho)$-regular and such that for every $i \in \mathbb{N}$, we have $F_{i} \subseteq B_{1}, \mathrm{NM}_{B_{1}}\left(F_{i}\right)<i^{-2 k-3}$, and such that for any half-space $P^{+} \subseteq \mathbb{H}^{2 k+1}$,

$$
\begin{equation*}
\frac{\mathscr{H}^{2 k+2}\left(\left(F \triangle P^{+}\right) \cap B_{1 / i}\right)}{\mathscr{H}^{2 k+2}\left(B_{1 / i}\right)} \geqslant \varepsilon . \tag{116}
\end{equation*}
$$

For every $i \in \mathbb{N}$, denote $S_{i}=s_{i}\left(F_{i}\right)$. Then $S_{i}$ is $(c, i)$-regular, so if $r>0$, $x \in \mathbb{H}^{2 k+1}$, and $i>r$, then $\mathscr{H}^{2 k+1}\left(B_{r}(x) \cap \partial S_{i}\right) \leqslant c r^{2 k+1}$. By Lemma 22, this implies that the sets $S_{i}$ have uniformly locally finite perimeter, so by Proposition 19, we can pass to a subsequence such that $\mathbf{1}_{S_{i}}$ converges in $L_{1}^{\text {loc }}$ to the characteristic function of some $S \subseteq \mathbb{H}^{2 k+1}$. We claim that $S$ is monotone.

We have $\mathrm{NM}_{B_{i}}\left(S_{i}\right)=\mathrm{NM}_{B_{1}}\left(F_{i}\right)<i^{-2 k-3}$, so by (111),

$$
\int_{\mathcal{L}\left(B_{i}\right)} \mathrm{NM}_{B_{i}}\left(S_{i}, L\right) \mathrm{d} \eta(L)=i^{2 k+2} \mathrm{NM}_{B_{i}}\left(S_{i}\right)<\frac{1}{i}
$$

Again passing to a subsequence, we suppose that for almost every line $L \in \mathcal{L}$, the intersection $L \cap S_{i}$ converges in $L_{1}^{\text {loc }}$ to $L \cap S$ and $\lim _{i \rightarrow \infty} \mathrm{NM}_{B_{i}}\left(S_{i}, L\right)=0$. That is, there is a sequence $I_{i} \subseteq L$ of monotone sets such that

$$
\begin{equation*}
\lim _{i \rightarrow \infty}\left\|\mathbf{1}_{B_{i} \cap I_{i}}-\mathbf{1}_{B_{i} \cap S_{i}}\right\|_{L_{1}(L)}=0 \tag{117}
\end{equation*}
$$

For any $r>0$, the sets $L \cap B_{r} \cap S_{i}$ converge (in $L_{1}$ ) to $L \cap B_{r} \cap S$. By (117), the sets $L \cap B_{r} \cap I_{i}$ also converge to $L \cap B_{r} \cap S$. A limit of monotone sets is monotone, so $L \cap S$ is monotone. Since this is true for almost every $L$, we conclude that $S$ is monotone and thus there is a plane $P$ such that $\mathscr{H}^{2 k+2}\left(S \triangle P^{+}\right)=0$. But by (116), we have $\mathscr{H}^{2 k+2}\left(\left(S_{i} \triangle P^{+}\right) \cap B_{1}\right) \geqslant \varepsilon \mathscr{H}^{2 k+2}\left(B_{1}\right)$ for all $i$. This contradiction completes the proof of Proposition 66.

Lemma 67 below shows that when $\left(F, F^{\mathrm{c}}, \partial F\right)$ is regular and $F$ is $\delta$-monotone, Proposition 66 also implies that $F$ and $P^{+}$are close in the local distance (Definition 51 ), which will be necessary to construct a corona decomposition.

Lemma 67. Fix $C>0$. For any $\theta>0$, there exists $\varepsilon>0$ such that if $r>0,\left(F, F^{c}, \partial F\right)$ is $(C, r)$-regular, $x \in \mathbb{H}^{2 k+1}$, and $P^{+} \subseteq \mathbb{H}^{2 k+1}$ is a half-space satisfying

$$
\frac{\mathscr{H}^{2 k+2}\left(\left(F \triangle P^{+}\right) \cap B_{r}(x)\right)}{\mathcal{H}^{2 k+2}\left(B_{r}(x)\right)} \leqslant \varepsilon,
$$

then $d_{B_{r / 2}(x)}\left(F, P^{+}\right)<\frac{\theta r}{2}$.
Proof. Without loss of generality, we suppose that $\theta<\frac{1}{4}$. Take any $\varepsilon \in \mathbb{R}$ that satisfies

$$
\begin{equation*}
0<\varepsilon<\min \left\{\frac{\theta^{2 k+1}}{C \mathscr{H}^{2 k+2}\left(B_{1}(x)\right)}, \frac{\theta^{2 k+1}}{2}\right\} . \tag{118}
\end{equation*}
$$

Suppose that $y \in\left(P^{+} \backslash F\right) \cap B_{r / 2}(x)$. If $d(y, P)>\theta r$, then $B_{\theta r}(y) \subseteq P^{+}$and thus

$$
\begin{aligned}
\mathscr{H}^{2 k+2}\left(\left(F \triangle P^{+}\right) \cap B_{r}(x)\right) & \geqslant \mathscr{H}^{2 k+2}\left(B_{\theta r}(y) \backslash F\right) \\
& =\frac{(\theta r)^{2 k+2}}{C} \stackrel{(118)}{>} \varepsilon \mathscr{H}^{2 k+2}\left(B_{r}(x)\right) .
\end{aligned}
$$

This is a contradiction, so $d(y, P) \leqslant \theta r$. Replacing $F$ by $\mathbb{H}^{2 k+1} \backslash F$ and $P^{+}$ by $P^{-}$, we find that $d(y, P) \leqslant \theta r$ for all $y \in\left(F \backslash P^{+}\right) \cap B_{r / 2}(x)$.

Now suppose that $y \in\left(F \triangle P^{+}\right) \cap B_{r / 2}(x)$ and $d(y, \partial F)>\theta r$. Then denote

$$
A= \begin{cases}P^{+} \cap B_{\theta r}(y) & \text { if } y \in P^{+} \\ \overline{P^{-}} \cap B_{\theta r}(y) & \text { if } y \in \overline{P^{-}}\end{cases}
$$

By applying a translation, we may suppose that $y=\mathbf{0}$. Let $R$ be the plane through $\mathbf{0}$ that is parallel (in the Euclidean sense) to $P$. Orient $R$ so that $R^{+} \cap B_{\theta r} \subseteq A$. Recall that for all $x \in \mathbb{H}^{2 k+1}$, we have $x^{-1}=-x$, so $\left(R^{+}\right)^{-1}=R^{-}$. By the left-invariance of the metric, we have $d(\mathbf{0}, g)=d\left(g^{-1}, \mathbf{0}\right)$, so $\left(R^{+} \cap B_{\theta r}\right)^{-1}=R^{-} \cap B_{\theta r}$. The inverse map is measure-preserving, so

$$
\begin{aligned}
& \mathscr{H}^{2 k+2}\left(R^{+} \cap B_{\theta r}\right)=\mathscr{H}^{2 k+2}\left(R^{-} \cap B_{\theta r}\right)=\frac{1}{2} \mathscr{H}^{2 k+2}\left(B_{\theta r}\right) \text { and thus } \\
& \mathscr{H}^{2 k+2}\left(\left(F \triangle P^{+}\right) \cap B_{r}(x)\right) \geqslant \mathscr{H}^{2 k+2}(A) \geqslant \mathscr{H}^{2 k+2}\left(R^{+} \cap B_{\theta r}\right) \\
& \geqslant \frac{1}{2} \mathscr{H}^{2 k+2}\left(B_{\theta r}\right) \stackrel{(118)}{>} \varepsilon \mathscr{H}^{2 k+2}\left(B_{r}(x)\right) .
\end{aligned}
$$

Again, this is a contradiction, so $d(y, \partial F) \leqslant \theta r$.
9.3. Most cubes are close to a plane. Now we begin the proof of Theorem 57. The first step is to show that if $\left(E, E^{\mathrm{c}}, \partial E\right)$ is regular, then most parts of $\partial E$ are close to a plane. Specifically,

Proposition 68. For every $C, \varepsilon>0$, there exists $K>0$ with the following property. Suppose that $r>0$, and let $E \subseteq \mathbb{H}^{2 k+1}$ be a set such that $\left(E, E^{c}, \partial E\right)$ is $(C, r)$-regular. Let $\Delta=\left\{\Delta_{i}\right\}_{i=-\infty}^{\left\lfloor\log _{2} r\right\rfloor}$ be a $(C, r)$-cubical patchwork for $\partial E$. Denote
$\mathcal{C}_{0}(\varepsilon) \stackrel{\text { def }}{=}\left\{Q \in \Delta: \exists\right.$ a plane $P_{Q}$ such that $\left.d_{N_{\frac{1}{\varepsilon}}(Q)}\left(P_{Q}^{+}, E\right) \leqslant \varepsilon \sigma(Q)\right\}$. ( $P_{Q}^{+}$is one of the half-spaces bounded by $P_{Q}$.) If $\mathscr{B}_{0}(\varepsilon) \stackrel{\text { def }}{=} \Delta \backslash \mathscr{G}_{0}(\varepsilon)$, then $\mathcal{B}_{0}(\varepsilon)$ is $K$-Carleson.

Note that by rescaling, we may suppose that $r=1$. We will prove the Proposition 68 by showing that most cubes in $\Delta$ have small nonmonotonicity, then using Proposition 66 to conclude that most cubes are close to planes. For simplicity of notation, all the implicit constants in this section and the following sections will depend on $C$.

First, we decompose $\operatorname{NM}(E)$ into pieces based on cubes. Namely, for every integer $j \leqslant\left\lfloor\log _{2} r\right\rfloor$ and $Q \in \Delta_{j}$, denote $w(Q)=\widehat{w}_{j}(E)(Q)$. This counts segments with an endpoint in $Q$ that have lengths between $\sigma(Q) / 2$ and $\sigma(Q)$.

Lemma 69. For any $Q \in \Delta$, we have

$$
\sum_{\substack{Q^{\prime} \in \Delta \\ Q^{\prime} \subseteq Q}} w\left(Q^{\prime}\right) \lesssim \sigma(Q)^{2 k+1} .
$$

It follows that for any $\varepsilon>0$, the set $\mathcal{B}(\varepsilon) \stackrel{\text { def }}{=}\left\{Q \in \Delta: w(Q)>\varepsilon \sigma(Q)^{2 k+1}\right\}$ is $O(1 / \varepsilon)$-Carleson.

Proof. Suppose $Q \in \Delta_{i}$ for some integer $i \leqslant\left\lfloor\log _{2} r\right\rfloor$. Then, using the regularity of $\partial E$,

$$
\begin{aligned}
\sum_{\substack{Q^{\prime} \in \Delta \\
Q^{\prime} \subseteq Q}} w\left(Q^{\prime}\right) & =\sum_{j=-\infty}^{i} \sum_{\substack{Q^{\prime} \in \Delta_{j} \\
Q^{\prime} \subseteq Q}} \widehat{w}_{j}(E)\left(Q^{\prime}\right) \\
& =\sum_{j=-\infty}^{i} \widehat{w}_{j}(E)(Q) \stackrel{(113)}{\leqslant} \operatorname{Per}(E)(Q) \lesssim \sigma(Q)^{2 k+1} .
\end{aligned}
$$

This implies that $\mathcal{B}(\varepsilon)$ is $O(1 / \varepsilon)$-Carleson because

$$
\sum_{\substack{Q^{\prime} \in \mathscr{B}(\varepsilon) \\ Q^{\prime} \subseteq Q}} \sigma\left(Q^{\prime}\right)^{2 k+1} \leqslant \sum_{\substack{Q^{\prime} \in \mathscr{B}(\varepsilon) \\ Q^{\prime} \subseteq Q}} \frac{w\left(Q^{\prime}\right)}{\varepsilon} \leqslant \frac{1}{\varepsilon} \sum_{\substack{Q^{\prime} \in \Delta \\ Q^{\prime} \subseteq Q}} w\left(Q^{\prime}\right) \lesssim \frac{\sigma(Q)^{2 k+1}}{\varepsilon} .
$$

If $Q, Q^{\prime} \in \Delta$ and $n>1$, we say that $Q$ and $Q^{\prime}$ are $n$-close if

$$
\begin{equation*}
\frac{\sigma(Q)}{\sigma\left(Q^{\prime}\right)} \in\left[2^{-n}, 2^{n}\right] \quad \text { and } \quad d\left(Q, Q^{\prime}\right) \leqslant 2^{n} \max \left\{\sigma(Q), \sigma\left(Q^{\prime}\right)\right\} \tag{119}
\end{equation*}
$$

If $\mathcal{S} \subseteq \Delta$, denote

$$
\mathcal{S}^{(n)} \stackrel{\text { def }}{=}\{Q \in \Delta: Q \text { is } n \text {-close to an element of } \mathcal{S}\} .
$$

Any cube that is $n$-close to $Q$ intersects $N_{2^{2 n}}(Q)$ and satisfies (119). By Definition 46, if $Q \in \Delta_{i}$ and $j<i$, then

$$
\begin{align*}
\left|\left\{Q^{\prime} \in \Delta_{j}: Q^{\prime} \cap N_{2^{2 n}}(Q) \neq \varnothing\right\}\right| & \lesssim \frac{\mathcal{H}^{2 k+1}\left(N_{2^{2 n}}(Q)\right)}{2^{j(2 k+1)}}  \tag{120}\\
& \asymp 2^{(2 k+1)(i-j+2 n)} \lesssim 2^{3 n(2 k+1)} .
\end{align*}
$$

If $j \geqslant i$, then

$$
\begin{equation*}
\left|\left\{Q^{\prime} \in \Delta_{j}: Q^{\prime} \cap N_{2^{2 n}}(Q) \neq \varnothing\right\}\right| \lesssim 2^{2 n(2 k+1)} . \tag{121}
\end{equation*}
$$

It follows that

$$
\begin{equation*}
\left|\{Q\}^{(n)}\right| \lesssim \sum_{j=i-n}^{i+n} 2^{3 n(2 k+1)} \lesssim n 2^{3 n(2 k+1)} \tag{122}
\end{equation*}
$$

that is, the number of cubes of $\Delta_{i}$ that are $n$-close to $Q$ is bounded by a function of $n$. Furthermore,

Lemma 70. For any $K, n>0$, there exists $L>0$ such that if $\mathcal{S} \subseteq \Delta$ is $K$-Carleson, then $\mathcal{S}^{(n)}$ is L-Carleson.

Proof. If $D_{1}, D_{2} \in \Delta$ are $n$-close and if $A_{1}, A_{2} \in \Delta$ are ancestors of $D_{1}, D_{2}$, respectively, such that $\sigma\left(A_{1}\right) / \sigma\left(A_{2}\right) \in\left[2^{-n}, 2^{n}\right]$, then $A_{1}$ and $A_{2}$ are also $n$-close. It follows that if $M \in \Delta$ and $Q \subseteq M$ is such that $Q \in \mathcal{S}^{(n)}$, then there is a cube $R \in \mathcal{S}$ such that some ancestor $A$ of $R$ is $n$-close to $M$.

Consequently, we can write

$$
\begin{equation*}
\sum_{\substack{Q \in \mathcal{S}^{(n)} \\ Q \subseteq M}} \sigma(Q)^{2 k+1} \leqslant \sum_{A \in\{M\}^{(n)}} \sum_{\substack{R \in \mathcal{S} \\ R \subseteq A}} \sum_{Q \in\{R\}^{(n)}} \sigma(Q)^{2 k+1} . \tag{123}
\end{equation*}
$$

But if $R \in \Delta$, then

$$
\begin{equation*}
\sum_{Q \in\{R\}^{(n)}} \sigma(Q)^{2 k+1} \stackrel{(122)}{\lesssim} n 2^{4 n(2 k+1)} \sigma(R)^{2 k+1} . \tag{124}
\end{equation*}
$$

Since $\mathcal{S}$ is $K$-Carleson,

$$
\begin{aligned}
\sum_{\substack{Q \in \mathcal{S}^{(n)} \\
Q \subseteq M}} \sigma(Q)^{2 k+1} & \stackrel{(123) \wedge(124)}{\lesssim} \sum_{A \subseteq\{M\}^{(n)}} \sum_{\substack{R \in \mathcal{S} \\
R \subseteq A}} n 2^{4 n(2 k+1)} \sigma(R)^{2 k+1} \\
& \lesssim n 2^{4 n(2 k+1)} \sum_{A \subseteq\{M\}^{(n)}} K \sigma(A)^{2 k+1} \\
& \stackrel{(122)}{\lesssim} n^{2} 2^{(2 k+1) 8 N} K \sigma(M)^{2 k+1} .
\end{aligned}
$$

By Lemmas 69 and $70, \mathcal{B}(\varepsilon)^{(n)}$ satisfies a Carleson packing condition. Let

$$
\mathcal{L}(\varepsilon, n) \stackrel{\text { def }}{=}\left\{Q \in \Delta \backslash \mathscr{B}(\varepsilon)^{(n)}: \sigma(Q) \leqslant 2^{-n}\right\} .
$$

Then $\Delta \backslash \mathscr{L}(\varepsilon, n)=\bigcup_{i=-n}^{0} \Delta_{i} \cup \mathscr{B}(\varepsilon)^{(n)}$ is $L$-Carleson (with $L$ allowed to depend on $n$ ), and every cube in $\mathcal{L}(\varepsilon, n)$ has small nonmonotonicity. The following lemma is inspired by [26, Prop. 4.5].

Lemma 71. For all $0<\varepsilon<1$, there are $\delta>0$ and $n \in \mathbb{N}$ such that if $Q \in \mathcal{L}(\delta, n)$ and $x \in Q$,

$$
\mathrm{NM}_{B_{\frac{\sigma(Q)}{\varepsilon}}(x)}(E)<\varepsilon
$$

Proof. We take $n \geqslant\left\lceil\log _{2}(1 / \varepsilon)\right\rceil+2$, so that $2^{-n} \leqslant \varepsilon / 4$. In this case, we shall bound $\operatorname{NM}_{B_{\sigma(Q) / \varepsilon}(x)}(E)$ in terms of $n$ and $\delta$. Suppose that $Q \in \mathcal{L}(\delta, n)$ and $x \in Q$, and denote $B=B_{\rho}(x)$, where $\rho=\sigma(Q) / \varepsilon \leqslant 1 / 4$. Let $m=$ $\left\lceil\log _{2} \rho\right\rceil+1 \leqslant-1$. Then $2^{m} \leqslant 2^{n} \sigma(Q)$. Proposition 62 states that

$$
\begin{equation*}
\operatorname{NM}_{B}(E) \lesssim \sum_{j=-\infty}^{m} \frac{2^{j}}{\rho} \cdot q_{j} \tag{125}
\end{equation*}
$$

where $q_{j} \stackrel{\text { def }}{=} \frac{\widehat{w}_{n}(E)(B)}{\operatorname{Per}(E)(B)}$. Note that for any $j \leqslant m,(113)$ implies that $q_{j} \leqslant 1$.
Suppose that $j \in[m-n, m]$. Then

$$
\widehat{w}_{j}(E)(B) \leqslant \sum_{\substack{Q^{\prime} \in \Delta_{j} \\ Q^{\prime} \cap B \neq \varnothing}} \widehat{w}_{j}(E)\left(Q^{\prime}\right) .
$$

If $Q^{\prime} \in \Delta_{j}$ and $Q^{\prime} \cap B \neq \varnothing$, then $\operatorname{diam} Q^{\prime} \leqslant C \sigma\left(Q^{\prime}\right) \leqslant C 2^{m} \leqslant 4 C \rho$ (where $C$ is the patchwork constant for $\Delta$ ). Thus $Q^{\prime} \subseteq B_{(4 C+1) \rho}(x)$. Furthermore, $Q^{\prime}$ intersects $N_{1 / \varepsilon}(Q)$, so $Q^{\prime}$ is $n$-close to $Q$. Since $Q \in \mathcal{L}(\delta, N)$, it follows that $\widehat{w}_{j}(E)\left(Q^{\prime}\right) \leqslant \delta \sigma\left(Q^{\prime}\right)^{2 k+1} \asymp \delta \operatorname{Per}(E)\left(Q^{\prime}\right)$. Consequently,

$$
\begin{aligned}
\widehat{w}_{j}(E)(B) & \sum \sum_{\substack{Q^{\prime} \in \Delta_{j} \\
Q^{\prime} \subseteq B_{(4 C+1) \rho}(x)}} \delta \operatorname{Per}(E)\left(Q^{\prime}\right) \\
& \leqslant \delta \operatorname{Per}(E)\left(B_{(4 C+1) \rho}(x)\right) \asymp \delta \operatorname{Per}(E)(B) .
\end{aligned}
$$

We have thus shown that $q_{j} \lesssim \delta$ for all $j \in[m-n, m]$. By (125) we therefore have

$$
\mathrm{NM}_{B}(E) \lesssim \sum_{j=-\infty}^{m-n-1} \frac{2^{j}}{\rho} \cdot q_{j}+\sum_{j=m-n}^{m} \frac{2^{j}}{\rho} \cdot q_{j} \lesssim \frac{2^{m-n}}{2^{m}}+\frac{\delta 2^{m}}{2^{m}}=\frac{1}{2^{n}}+\delta
$$

Hence, if $n$ is sufficiently large and $\delta$ is sufficiently small, then $\operatorname{NM}_{B}(E) \leqslant \varepsilon$.

Proposition 68 now follows from Lemma 71 and the results on monotone sets in Section 9.2.

Proof of Proposition 68. Set $\theta=\frac{\varepsilon}{C+1 / \varepsilon}$. By Proposition 66 and Lemma 67, there is $\alpha>0$ such that if $x \in \partial E, 0<r<1$, and $\operatorname{NM}_{B_{r}(x)}(E)<\alpha^{2 k+3}$, then there is a half-space $P^{+} \subseteq \mathbb{H}^{2 k+1}$ with

$$
d_{B_{\frac{\alpha r}{2}}(x)}\left(F, P^{+}\right)<\frac{\theta \alpha r}{2} .
$$

Denote $\gamma=\frac{2(C+1 / \varepsilon)}{\alpha}$. By Lemma 71, there are $\delta>0$ and $n>0$ such that if $Q \in \mathscr{L}(\delta, n)$ and $x \in Q$, then $\operatorname{NM}_{B_{\gamma \sigma(Q)}(x)}(E)<\alpha^{2 k+3}$. Thus there is a plane $P$ such that

$$
d_{B_{\frac{\alpha \gamma \sigma(Q)}{2}}(x)}\left(E, P^{+}\right)=d_{B_{\left(C+\frac{1}{\varepsilon}\right) \sigma(Q)}(x)}\left(E, P^{+}\right) \leqslant\left(C+\frac{1}{\varepsilon}\right) \theta \sigma(Q)=\varepsilon \sigma(Q) .
$$

Since $N_{1 / \varepsilon}(Q) \subseteq B_{(C+1 / \varepsilon) \sigma(Q)}(x)$, this implies that $Q \in \mathcal{C}_{00}(\varepsilon)$.
Let $\mathscr{\mathcal { L }}=\mathscr{\mathcal { L }}(\delta, n), \mathscr{B}=\Delta \backslash \mathscr{L}$. By Lemmas 69 and $70, \mathscr{B}$ is $K$-Carleson for some $K$ depending on $C$. We have seen that $\mathscr{L} \subseteq \mathcal{L}_{0}(\varepsilon)$, so $\mathscr{B}_{0}(\varepsilon) \subseteq \mathscr{B}$ and thus $\mathscr{B}_{0}(\varepsilon)$ is $K$-Carleson.
9.4. Most cubes are close to a vertical plane. Having shown that most cubes are close to planes, we now show that most of these planes are vertical or nearly vertical.
9.4.1. Planes and angles. Before the next step, we recall some facts about planes and angles in the Heisenberg group. Recall that $\mathrm{H} \subseteq \mathbb{H}^{2 k+1}$, the subspace of horizontal vectors, is the subspace spanned by $\left\{X_{i}\right\}_{i=1}^{k}$ and $\left\{Y_{i}\right\}_{i=1}^{k}$. Recall also that a plane in $\mathbb{H}^{2 k+1}$ is either a vertical plane or the union $P_{y}=y \mathbf{H}$ of all of the horizontal lines passing through a point $y \in \mathbb{H}^{2 k+1}$. The planes in $\mathbb{H}^{2 k+1}$ are exactly the planes in $\mathbb{R}^{2 k+1}$ under the identification with $\mathbb{H}^{2 k+1}$.

If $P, Q \subseteq \mathbb{H}^{2 k+1}$ are planes, we define $\angle(P, Q)$ to be the angle between $P$ and $Q$ considered as planes in $\mathbb{R}^{2 k+1}$ (i.e., the Euclidean angle between their normals). If $P$ is a plane and $L$ is a line (a coset of a one-parameter subgroup, or equivalently, a line in $\mathbb{R}^{2 k+1}$ ) that intersects $P$ at a point $y$, we define $\angle(P, L)$ to be the minimum Euclidean angle between $L$ and a line in $P$ that passes through $y$; this quantity varies between 0 and $\frac{\pi}{2}$.

Importantly, the angle between planes is neither translation-invariant nor scale-invariant. For example, if $y \in \mathrm{H} \backslash\{\mathbf{0}\}$, then $P_{y}=y \mathrm{H}$ is a plane through $\mathbf{0}$.

If $L_{Z}=\langle Z\rangle$ is a vertical line, then $\angle\left(s_{t}\left(P_{y}\right), s_{t}\left(L_{Z}\right)\right)=\angle\left(s_{t}\left(P_{y}\right), L_{Z}\right)$ converges to zero as $t \rightarrow \infty$.

Since horizontal planes play an important role in this section, we give two formulas expressing a horizontal plane in coordinate form. Let $\pi: \mathbb{H}^{2 k+1} \rightarrow H$ be the abelianization map. If $u, p \in \mathbb{H}^{2 k+1}$, then $p \in P_{u}$ if and only if $p=$ $u(\pi(p)-\pi(u))$. That is, we can write $P_{u}$ as a graph over H:

$$
\begin{equation*}
P_{u}=\left\{p \in \mathbb{H}^{2 k+1}: z(p)=z(u)+\frac{\omega(u, p)}{2}\right\}=\left\{h Z^{z(u)+\frac{\omega(u, h)}{2}}: h \in \mathrm{H}\right\} . \tag{126}
\end{equation*}
$$

If $r>0$ and $u=Y_{1}^{r}$, then we can write

$$
\begin{equation*}
P_{u}=\left\{p: z(p)=\frac{\omega(u, p)}{2}\right\}=\left\{p: x_{1}(p)=\frac{-2 z(p)}{r}\right\} . \tag{127}
\end{equation*}
$$

If $P$ is a plane and $x \in \mathbb{H}^{2 k+1}$, denote

$$
\alpha_{x}(P) \stackrel{\text { def }}{=} \begin{cases}d_{\mathrm{Euc}}(\pi(x), \pi(y)) & \text { if } P=P_{y}, \\ \infty & \text { if } P \text { is vertical. }\end{cases}
$$

By (126), $\alpha_{0}(P)$ is determined by the angle between $P$ and the vertical line $L_{Z} \stackrel{\text { def }}{=}\langle Z\rangle$; specifically,

$$
\angle\left(L_{Z}, P\right)=\arctan \left(\frac{2}{\alpha_{\mathbf{0}}(P)}\right) .
$$

Lemma 72. Let $P^{+} \subseteq \mathbb{H}^{2 k+1}$ be a half-space with boundary $P$, and let $x \in \mathbb{H}^{2 k+1}$. There is a half-space $V^{+}$with vertical boundary such that

$$
d_{B_{1}(x)}\left(P^{+}, V^{+}\right) \lesssim \frac{1}{\alpha_{x}(P)}
$$

Proof. If $P$ is vertical, then the lemma is satisfied by letting $V^{+}=P^{+}$. If $P$ does not intersect $B_{1}(x)$, then either $B_{1}(x) \subseteq P^{+}$or $B_{1}(x) \cap P^{+}=\varnothing$. In this case, we can choose $V$ to be a vertical plane disjoint from $B_{1}(x)$ and choose $V^{+}$to be one of the sides of $V$.

We thus suppose that $P=P_{y}$ for some $y \in \mathbb{H}^{2 k+1}$ and that $P$ intersects $B_{1}(x)$, say at $w$. By translating and rotating, we may suppose that $w=\mathbf{0}$, that $x \in B_{1}$, and $y=Y_{1}^{r}$, where $r=\alpha_{w}(P)$.

If $d(x, y) \leqslant 2$, choose $V$ to be an arbitrary vertical plane through $\mathbf{0}$. Then the lemma is satisfied because $d_{B_{1}(x)}\left(P^{+}, V^{+}\right) \leqslant \operatorname{diam}\left(B_{1}(x)\right) \lesssim 1$. So, suppose that $d(x, y)>2$, thus $r \asymp \alpha_{x}(P)$. Let $V=\left\{h \in \mathbb{H}^{2 k+1}: x_{1}(h)=0\right\}$. By (127) we have $P=\left\{p: x_{1}(p)=-2 z(p) / r\right\}$.

If $P^{+}=\left\{p \in \mathbb{H}^{2 k+1}: x_{1}(p) \geqslant-2 z(p) / r\right\}$, let $V^{+}=\left\{p \in \mathbb{H}^{2 k+1}:\right.$ $\left.x_{1}(p) \geqslant 0\right\}$. We claim that $d_{B_{2}}\left(P^{+}, V^{+}\right)$is small. If $q \in\left(P^{+} \triangle V^{+}\right) \cap B_{2}$, then $|z(q)| \lesssim 1$ and $\left|x_{1}(q)\right| \leqslant 2|z(q)| / r \lesssim 1 / r$. It follows that $d(q, V) \lesssim 1 / r$ and $d(q, P) \lesssim 1 / r$, so $d_{B_{1}(x)}\left(P^{+}, V^{+}\right) \leqslant d_{B_{2}}\left(P^{+}, V^{+}\right) \lesssim 1 / r$, as desired. Likewise,
if $P^{+}=\left\{p \in \mathbb{H}^{2 k+1}: x_{1}(p) \leqslant-2 z(p) / r\right\}$, then choose $V^{+}=\left\{p \in \mathbb{H}^{2 k+1}:\right.$ $\left.x_{1}(p) \leqslant 0\right\}$ and again $d_{B_{1}(x)}\left(P^{+}, V^{+}\right) \lesssim 1 / r$.

By applying a scaling to Lemma 72, we find that
Corollary 73. For any $\sigma>0$, if $P^{+} \subseteq \mathbb{H}^{2 k+1}$ is a half-space and $x \in \mathbb{H}^{2 k+1}$, then there is a half-space $V^{+}$with vertical boundary such that

$$
\frac{d_{B_{\sigma}(x)}\left(P^{+}, V^{+}\right)}{\sigma} \lesssim \frac{\sigma}{\alpha_{x}(P)}
$$

Finally, we will need the following simple lemma.
Lemma 74. For any $\rho>1$, there exists $0<\delta<\frac{1}{2}$ such that for any $x \in \mathbb{H}^{2 k+1}$ and any $\sigma>0$, if $P, R \subseteq \mathbb{H}^{2 k+1}$ are planes such that $d(x, P)<\delta \sigma$, $d_{B_{\sigma}(x)}\left(P^{+}, R^{+}\right)<\delta \sigma$, and $\alpha_{x}(P)<\rho \sigma$, then $\alpha_{x}(R)<2 \rho \sigma$.

Proof. After rescaling and translating, we may suppose that $\sigma=1$ and $x=\mathbf{0}$. Then $\alpha_{x}(P)<\rho$ and $\angle\left(P, L_{Z}\right)>\arctan (2 / \rho)$. If $\delta$ is sufficiently small, then $\angle(P, R)$ is small, so $\angle\left(R, L_{Z}\right)>\arctan (1 / \rho)$ and thus $\alpha_{x}(R)<2 \rho$, as desired.
9.4.2. Approximating cubes by vertical planes. The goal here is to prove the following proposition.

Proposition 75. For every $0<\varepsilon<1$, there is a set $\mathcal{L}_{\mathcal{L}}(\varepsilon) \subseteq \Delta$ such that for every $Q \in \mathcal{L}_{\mathcal{L}}(\varepsilon)$, there is a vertical half-space $V_{Q}^{+}$with bounding plane $V_{Q}$ such that

$$
\begin{equation*}
d_{N_{\frac{1}{\varepsilon}}(Q)}\left(V_{Q}^{+}, E\right) \leqslant \varepsilon \sigma(Q) \tag{128}
\end{equation*}
$$

Moreover, $\mathcal{B} \stackrel{\text { def }}{=} \Delta \backslash \mathscr{C}(\varepsilon)$ satisfies a Carleson packing condition with constant that depends on $\varepsilon$.

Let $\eta>0$ be a small number to be determined later. By Proposition 68, there are a set of cubes $\mathcal{C}_{0}=\mathcal{G}_{\mathcal{L}}(\eta) \subseteq \Delta$ and a collection of planes $\left\{P_{Q}\right\}_{Q \in \Delta}$ such that $d_{N_{1 / \eta}(Q)}\left(P_{Q}^{+}, E\right) \leqslant \eta \sigma(Q)$ for all $Q \in \mathcal{L}_{0}$. The set $\mathcal{G}_{\mathcal{L}}=\mathcal{L}_{\mathcal{L}}(\varepsilon)$ will consist of the cubes $Q \in \mathcal{L}_{\mathcal{L}}$ for which $P_{Q}$ is close to vertical. Specifically, if $Q \in \mathcal{L}_{0}$, denote

$$
\alpha(Q) \stackrel{\text { def }}{=} \frac{\min _{x \in Q} \alpha_{x}\left(P_{Q}\right)}{\sigma(Q)}
$$

Corollary 73 implies that if $x \in Q$, then there is a vertical plane $V$ such that

$$
\begin{equation*}
d_{N_{\frac{2}{\varepsilon}}(Q)}\left(P_{Q}^{+}, V^{+}\right) \leqslant d_{\left(\frac{2}{\varepsilon}+C\right) \sigma(Q)}\left(P_{Q}^{+}, V^{+}\right) \lesssim \frac{\left(\frac{2}{\varepsilon}+C\right)^{2} \sigma(Q)^{2}}{\alpha_{x}\left(P_{Q}\right)} \lesssim \frac{\sigma(Q)}{\varepsilon^{2} \alpha(Q)} \tag{129}
\end{equation*}
$$

We will define $\mathcal{L}_{\mathcal{L}}=\mathcal{L}_{0} \backslash \mathcal{B}_{1}$, where $\mathscr{B}_{1} \stackrel{\text { def }}{=}\left\{Q \in \mathcal{L}_{0}: \alpha(Q)<r\right\}$ are $r>0$ will be a quantity depending on $\varepsilon$ that will be determined later.

The main difficulty is to prove that $\mathscr{B}_{1}$ is Carleson. If $Q, R \in \mathcal{G}_{0}$, we say that $R$ is a good descendant of $Q$ if $R \subseteq Q$ and $A \in \mathcal{L}_{0}$ for all $A \in \Delta$ such that $R \subseteq A \subseteq Q$. Let $G(Q)$ be the set of good descendants of $Q$, and for every $n \in \mathbb{N} \cup\{0\}$, define

$$
G(Q, r) \stackrel{\text { def }}{=}\{R \in G(Q): \alpha(R)<r\}
$$

and

$$
G_{n}(Q, r) \stackrel{\text { def }}{=}\left\{R \in G(Q, r): \sigma(R)=2^{-n} \sigma(Q)\right\}
$$

Our goal is to bound the size of $G(Q, r)$.
Lemma 76. Fix $r>0$. If $\eta$ is sufficiently small, then for all $Q \in \mathcal{G}_{0}$ and all $R \in G(Q, r)$, every $A \in \Delta$ such that $R \subseteq A \subseteq Q$ is an element of $G(Q, r)$. Consequently, for all integers $m, n \geqslant 0$,

$$
\begin{equation*}
G_{m+n}(Q, r)=\bigcup_{R \in G_{m}(Q, r)} G_{n}(R, r) . \tag{130}
\end{equation*}
$$

In fact, there exists $0<\lambda<1$ depending on $r$ such that for every integer $n \geqslant 0$ we have

$$
\begin{equation*}
\sum_{R \in G_{n}(Q, r)} \sigma(R)^{2 k+1} \lesssim \lambda^{n} \sigma(Q)^{2 k+1} \tag{131}
\end{equation*}
$$

and thus

$$
\sum_{R \in G(Q, r)} \sigma(R)^{2 k+1} \lesssim_{r} \sigma(Q)^{2 k+1}
$$

Proof. To prove the first part of the lemma, let $R \in G(Q, r)$ and let $A \in G(Q)$ be the parent of $R$. We need to show that if $\eta$ is sufficiently small, then $A \in G(Q, r)$, i.e., $\alpha(A)<r$.

By part (4) of Lemma 52 , if $x \in R$ and $\eta<\frac{1}{2}$, then

$$
\begin{aligned}
d_{B_{\sigma(R)}(x)}\left(P_{A}^{+}, P_{R}^{+}\right) & \leqslant d_{B_{2 \sigma(R)}(x)}\left(P_{A}^{+}, E\right)+d_{B_{2 \sigma(R)}(x)}\left(P_{R}^{+}, E\right) \\
& \leqslant \eta \sigma(A)+\eta \sigma(R)=3 \eta \sigma(R) .
\end{aligned}
$$

Since $\alpha(R)<r$, there exists $x \in R$ such that $\alpha_{x}\left(P_{R}\right)<r \sigma(R)$. We have $d\left(x, P_{R}\right) \leqslant \eta \sigma(R)$. By Lemma 74, if $\eta$ is sufficiently small, then $\alpha_{x}\left(P_{A}\right)<$ $2 r \sigma(R)$ and $\alpha(A) \leqslant \alpha_{x}\left(P_{A}\right) / \sigma(A)<r$.

For every $Q \in \mathcal{L}_{0}$ and $n \geqslant 0$, denote $g_{Q, n} \stackrel{\text { def }}{=} \sum_{R \in G_{n}(Q, r)} \sigma(R)^{2 k+1}$. To prove the second part, we will first show that there is some $n_{0}$ and some $\eta$ such that for all $Q \in \mathcal{G}_{\mathcal{L}}$, we have

$$
\begin{equation*}
g_{Q, n_{0}} \leqslant \frac{\sigma(Q)^{2 k+1}}{2} \tag{132}
\end{equation*}
$$

This estimate implies the desired conclusion (131) because for any $n \geqslant 0$, we have

$$
g_{Q, n+n_{0}} \stackrel{(130)}{=} \sum_{A \in G_{n}(Q, r)} g_{A, n_{0}} \stackrel{(132)}{\lessgtr} \sum_{A \in G_{n}(Q, r)} \frac{\sigma(A)^{2 k+1}}{2}=\frac{g_{Q, n}}{2}
$$

Furthermore, for all $n \geqslant 0$,

$$
\begin{aligned}
g_{Q, n}=\sum_{R \in G_{n}(Q, r)} \sigma(R)^{2 k+1} & \leqslant C \sum_{R \in G_{n}(Q, r)} \mathcal{H}^{2 k+1}(R) \\
& \leqslant C \mathcal{H}^{2 k+1}(Q) \leqslant C^{2} \sigma(Q)^{2 k+1} .
\end{aligned}
$$

For every integer $n \geqslant 0$, write $n=m n_{0}+i$, where $m=\left\lfloor n / n_{0}\right\rfloor$. Then $i \geqslant 0$ and it follows from the above estimates that $g_{Q, n} \leqslant 2^{-m} g_{Q, i} \leqslant 2 C^{2} 2^{-n / n_{0}} \sigma(Q)^{2 k+1}$. Hence (131) holds with $\lambda=2^{-1 / n_{0}}$. It therefore remains to justify (132).

Let $Q \in \mathcal{L}_{0}$. If $\alpha(Q) \geqslant r$, then $G_{0}(Q, r)=\varnothing$ and therefore (130) implies that $G_{n}(Q, r)=\varnothing$. So, suppose that $\alpha(Q)<r$. Let $y$ be the center of $P_{Q}$, i.e., $P_{y}=P_{Q}$. Suppose that $\delta>0$ is such that Lemma 74 holds true for $\rho=r$. Let $n \in \mathbb{N}$ and $R \in G_{n}(Q, r)$. We claim that $R$ is close to $y$ if $\eta<2^{-n-1} \delta$. Indeed, let $x \in R$ be such that $\alpha_{x}(R)<r \sigma(R)$. Then $d\left(x, P_{Q}\right)<\eta \sigma(Q)<\delta \sigma(R)$ and

$$
\begin{aligned}
d_{B_{\sigma(R)}(x)}\left(P_{Q}^{+}, P_{R}^{+}\right) & \leqslant d_{B_{2 \sigma(R)}(x)}\left(P_{Q}^{+}, E\right)+d_{B_{2 \sigma(R)}(x)}\left(P_{R}^{+}, E\right) \\
& \leqslant \eta(\sigma(Q)+\sigma(R))<\delta \sigma(R) .
\end{aligned}
$$

Since $\alpha_{x}\left(P_{R}\right)<r \sigma(R)$, Lemma 74 implies that $\alpha_{x}\left(P_{Q}\right)<2 r \sigma(R)$. Let $z \in P_{Q}$ be a point such that $d(x, z) \leqslant 2 d\left(x, P_{Q}\right) \leqslant 2 \eta \sigma(Q)$. Then $d(z, y)=$ $d(\pi(z), \pi(y)) \leqslant \alpha_{x}\left(P_{Q}\right)+2 \eta \sigma(Q)$ and

$$
\begin{aligned}
d(x, y) & \leqslant d(x, z)+d(z, y) \\
& \leqslant \alpha_{x}\left(P_{Q}\right)+4 \eta \sigma(Q) \leqslant(4 \delta+2 r) \sigma(R)=2^{-n}(4 \delta+2 r) \sigma(Q) .
\end{aligned}
$$

Consequently, if $R \in G_{n_{0}}(Q, r)$, then $R \subseteq B_{2^{-n_{0}}(4 \delta+2 r+C) \sigma(Q)}(y)$, where $\delta$ depends on $r$. Choose $n_{0}$ (depending on $r$ and $C$ ) sufficiently large that $2^{-n_{0}}(4 \delta+2 r+C)<1 /(2 C)$. Then

$$
\begin{aligned}
g_{Q, n_{0}} & \leqslant C \sum_{R \in G_{n}(Q, r)} \mathscr{H}^{2 k+1}(R) \leqslant C \mathscr{H}^{2 k+1}\left(Q \cap B_{\frac{\sigma(Q)}{2 C}}(y)\right) \\
& \leqslant C^{2} \cdot \frac{\sigma(Q)^{2 k+1}}{(2 C)^{2 k+1}} \leqslant \frac{\sigma(Q)^{2 k+1}}{2},
\end{aligned}
$$

which is the desired estimate (132).
Proof of Proposition 75. By (129), there exists $r \asymp 1 / \varepsilon^{3}$ so that if $Q \in \Delta$, $x \in Q$, and $P$ is a plane such that $\alpha_{x}(P) \geqslant r \sigma(Q)$, then there is a vertical plane $V$ such that

$$
d_{N_{\frac{2}{\varepsilon}}(Q)}\left(P^{+}, V^{+}\right) \leqslant \frac{\varepsilon \sigma(Q)}{2} .
$$

Let $0<\eta<\varepsilon / 2$ be such that Lemma 76 holds for this choice of $r$. Let $\mathcal{L}_{0}=\mathcal{L}_{0}(\eta)$ and $\mathscr{B}_{0}=\mathscr{B}_{0}(\eta)$ be as in Proposition 68. Denote $\mathscr{B}_{1}=\left\{Q \in \mathcal{L}_{0}\right.$ : $\alpha(Q)<r\}$ and $\mathscr{L}=\mathcal{L}_{0} \backslash \mathscr{B}_{1}$.

If $Q \in \mathcal{L}_{\mathcal{L}}$, then by part (4) of Lemma 52 we have
$d_{N_{\frac{1}{\varepsilon}}}\left(E, V_{Q}^{+}\right) \leqslant d_{N_{\frac{2}{\varepsilon}}(Q)}\left(E, P_{Q}^{+}\right)+d_{N_{\frac{2}{\varepsilon}}(Q)}\left(P_{Q}^{+}, V_{Q}^{+}\right) \leqslant \frac{\varepsilon \sigma(Q)}{2}+\frac{\varepsilon \sigma(Q)}{2}=\varepsilon \sigma(Q)$.
We claim that $\Delta \backslash \mathcal{L}=\mathcal{B}_{0} \cup \mathscr{B}_{1}$ satisfies a Carleson packing condition with constant that depends on $\varepsilon$. Since this is true for $\mathscr{B}_{0}$ by Proposition 68 , it remains only to prove that $\mathcal{B}_{1}$ is Carleson.

For all $R \in \Delta$ such that $\sigma(R)<1$, let $F(R)$ be the parent of $R$. Recall that $\Delta_{0} \subseteq \mathscr{B}_{0}$, so if $R \in \mathscr{B}_{1}$, then there is a smallest natural number $i$ such that $F^{i}(R) \in \mathscr{B}_{0}$. Lemma 76 implies that if $F^{j}(R) \in \mathcal{B}_{1}$, then either $F^{j+1}(R) \in \mathcal{B}_{0}$ or $F^{j+1}(R) \in \mathscr{B}_{1}$, so, by induction, $F^{j}(R) \in \mathscr{B}_{1}$ for all $j=0, \ldots, i-1$. That is, if $A=F^{i-1}(R)$, then $F(A) \in \mathcal{B}_{0}$ and $R \in G(A, r)$.

Let $M=\left\{A \in \Delta: F(A) \in \mathcal{B}_{0}\right\}$. Then $\mathcal{B}_{1}=\bigcup_{A \in M} G(A, r)$. By Lemma $70, M$ is $L$-Carleson for some $L$ depending on $\varepsilon$. If $Q \in \Delta, R \subseteq Q$, and $R \in \mathscr{B}_{1}$, then either $R \in G(Q, r)$ (where we set $G(Q, r)=\varnothing$ if $Q \in \mathscr{B}_{0}$ ) or $R \in G(A, r)$ for some $A \in M$ such that $A \subseteq Q$. Hence, by Lemma 76,

$$
\begin{aligned}
\sum_{\substack{R \in \mathcal{B}_{1} \\
R \subseteq Q}} \sigma(R)^{2 k+1} & =\sum_{R \in G(Q, r)} \sigma(R)^{2 k+1}+\sum_{\substack{A \in M \\
A \subseteq Q}} \sum_{R \in G(A, r)} \sigma(R)^{2 k+1} \\
& \lesssim \varepsilon \sigma(Q)^{2 k+1}+\sum_{\substack{A \in M \\
A \subseteq Q}} \sigma(A)^{2 k+1} \leqslant(L+1) \sigma(Q)^{2 k+1} .
\end{aligned}
$$

9.5. Constructing stopping-time regions. In this section, we group the cubes in $\mathscr{L}_{\mathcal{L}}$ into stopping-time regions. We will prove the following proposition. The arguments in this section are based on arguments in Chapter 7 of [30], except for Lemma 79, which uses nonmonotonicity.

Proposition 77. For every $0<\lambda<\frac{\pi}{2}$, if $\varepsilon>0$ is sufficiently small (depending on $\lambda$ ) and if $\mathcal{L}=\mathcal{L}(\varepsilon) \subseteq \Delta$ and $\mathscr{B}=\Delta \backslash \mathcal{G}$ are as in Proposition 75 , then we can partition $\mathcal{G}$ into a set $\mathcal{F}$ of stopping-time regions such that $(\mathscr{B}, \mathcal{G}, \mathcal{F})$ is a coronization (Definition 50) of $\partial E$ and

$$
\begin{equation*}
\forall \mathcal{S} \in \mathcal{F}, \forall Q \in \mathcal{S} \quad \angle\left(V_{Q}, V_{Q(\mathcal{S})}\right)<\lambda . \tag{133}
\end{equation*}
$$

The Carleson constants of the coronization will depend only on $\lambda, \varepsilon, C$, and $k$.
For every $Q \in \Delta$, let $\mathcal{S}_{Q}$ be the maximal coherent set such that $Q$ is the maximal cube in $\mathcal{S}_{Q}$ and such that for all $R \in \mathcal{S}_{Q}$, we have $\angle\left(V_{R}, V_{Q}\right)<\lambda$. We construct such a set by the usual inductive stopping-time argument; we start by adding $Q$, then for each $R \in \mathcal{S}_{Q}$, we add the children of $R$ to $\mathcal{S}_{Q}$ if and only if they are all in $\mathcal{L}_{\mathcal{L}}$ and all of their approximating vertical planes are $\lambda$-close to $V_{Q}$.

We partition $\mathcal{G}$ into such regions by applying this construction repeatedly. First, let $Q$ be a maximal element of $\mathcal{G}$. We add $\mathcal{S}_{Q}$ to $\mathcal{F}$ and remove the
elements of $\mathcal{S}_{Q}$ from $\mathcal{G}$, then take a maximal element from the remaining cubes and repeat. By design, this results in a partition of $\mathcal{G}$ into coherent regions that satisfies (133).

The main difficulty is to show that the maximal elements of the stoppingtime regions satisfy a Carleson packing condition. To prove this, we divide $\mathcal{F}$ into three parts. For $\mathcal{S} \in \mathcal{F}$, let $m(\mathcal{S}) \subseteq \mathcal{S}$ be the set of minimal cubes in $\mathcal{S}$. These cubes are disjoint, but need not cover $\mathrm{Q}(\mathcal{S})$ because there may be points that are contained in arbitrarily small cubes. Each element of $m(\mathcal{S})$ has a child $Q$ such that either $Q \in \mathscr{B}$ or $\angle\left(V_{Q}, V_{\mathrm{Q}(\mathcal{S})}\right)>\lambda$. Let $m_{0}(\mathcal{S})$ be the set of minimal cubes with at least one child in $\mathscr{B}$ and let $m_{1}(\mathcal{S})$ be the set of minimal cubes with at least one child such that $\angle\left(V_{Q}, V_{\mathrm{Q}(\delta)}\right)>\lambda$. Then $m(\mathcal{S})=m_{0}(\mathcal{S}) \cup m_{1}(\mathcal{S})$. Let

$$
\begin{align*}
& \mathscr{F}_{0}=\left\{\mathcal{S} \in \mathcal{F}: \mathscr{H}^{2 k+1}\left(\bigcup_{Q \in m_{0}(\mathcal{S})} Q\right) \geqslant \frac{\mathscr{H}^{2 k+1}(\mathrm{Q}(\mathcal{S}))}{4}\right\} \\
& \mathcal{F}_{1}=\left\{\mathcal{S} \in \mathcal{F}: \mathscr{H}^{2 k+1}\left(\bigcup_{Q \in m_{1}(\mathcal{S})} Q\right) \geqslant \frac{\mathscr{H}^{2 k+1}(\mathrm{Q}(\mathcal{S}))}{2}\right\},  \tag{134}\\
& \mathcal{F}_{2}=\left\{\mathcal{S} \in \mathcal{F}: \mathscr{H}^{2 k+1}\left(\mathrm{Q}(\mathcal{S}) \backslash \bigcup_{Q \in m(\mathcal{S})} Q\right) \geqslant \frac{\mathscr{H}^{2 k+1}(\mathrm{Q}(\mathcal{S}))}{4}\right\} .
\end{align*}
$$

Then $\mathscr{F}=\mathscr{F}_{0} \cup \mathscr{F}_{1} \cup \mathscr{F}_{2}$. For every $i \in\{0,1,2\}$ let $M_{i}=\left\{\mathrm{Q}(\mathcal{S}): \mathcal{S} \in \mathscr{F}_{i}\right\}$, and if $R \in \Delta$, let $M_{i}(R)=\left\{Q \in M_{i}: Q \subseteq R\right\}$ and $\mathscr{F}_{i}(R)=\left\{\mathcal{\delta} \in \mathscr{F}_{i}: Q(\mathcal{S}) \subseteq R\right\}$. We claim that there are $K_{0}, K_{1}, K_{2}>0$ depending on $\lambda$ such that $M_{i}$ is $K_{i}$-Carleson for each $i \in\{0,1,2\}$.

Firstly, we consider $M_{0}$. For every $Q \in \Delta$, let $F(Q)$ be the parent of $Q$. For each $\mathcal{S} \in \mathcal{F}$, let $Y_{\mathcal{S}}=\{Q \in \mathcal{B}: F(Q) \in \mathcal{S}\}$ be the set of bad children of elements of $m_{0}(\mathcal{S})$. If $Q \in m_{0}(\mathcal{S})$, then $Q$ has a child $D \in Y_{\mathcal{S}}$, so

$$
\sum_{Q \in m_{0}(\delta)} \sigma(Q)^{2 k+1} \asymp \sum_{Q \in Y_{S}} \sigma(Q)^{2 k+1} .
$$

Consequently, if $\mathcal{S} \in \mathscr{F}_{0}$, then we have

$$
\sum_{Q \in Y_{\mathcal{S}}} \sigma(Q)^{2 k+1} \asymp \sigma(\mathrm{Q}(\mathcal{S}))^{2 k+1}
$$

Because the elements of $\mathscr{F}$ are disjoint, the sets $\left\{Y_{\mathcal{S}}\right\}_{\mathcal{S} \in \mathcal{F}}$ are disjoint subsets of $\mathscr{B}$, so for $R \in \Delta$,

$$
\sum_{Q \in M_{0}(R)} \sigma(Q)^{2 k+1} \asymp \sum_{\delta \in \mathscr{F}_{0}(R)} \sum_{Q \in Y_{\mathcal{S}}} \sigma(Q)^{2 k+1} \leqslant \sum_{\substack{Q \in \mathscr{B} \\ Q \subseteq R}} \sigma(Q)^{2 k+1} \lesssim \sigma(R)^{2 k+1} .
$$

That is, the maximal cubes of $\mathscr{F}_{0}$ satisfy the desired Carleson packing condition.

Next, we consider $M_{2}$. The sets $\left\{U_{\mathcal{S}}=\mathrm{Q}(\mathcal{S}) \backslash \bigcup_{Q \in m(\mathcal{S})} Q\right\}_{\mathcal{S} \in \mathscr{F}}$ are pairwise disjoint and satisfy $\mathscr{H}^{2 k+1}\left(U_{\mathcal{S}}\right) \gtrsim \sigma(\mathrm{Q}(\mathcal{S}))^{2 k+1}$ for all $\mathcal{S} \in \mathscr{F}_{2}$. If $Q \in \Delta$, then

$$
\begin{aligned}
\sum_{\mathcal{S} \in \mathscr{F}_{2}(Q)} \sigma(\mathrm{Q}(\mathcal{S}))^{2 k+1} & \lesssim \sum_{\mathcal{S} \in \mathscr{F}_{2}(Q)} \mathscr{H}^{2 k+1}\left(U_{\mathcal{S}}\right) \\
& =\mathscr{H}^{2 k+1}\left(\bigcup_{\mathcal{S} \in \mathcal{F}_{2}(Q)} U_{\mathcal{S}}\right) \leqslant \mathscr{H}^{2 k+1}(Q) \asymp \sigma(Q)^{2 k+1},
\end{aligned}
$$

so $M_{2}$ satisfies the desired Carleson packing condition as well.
It remains to show that $M_{1}$ is Carleson. Recall that for every integer $j \leqslant 0$ and $Q \in \Delta_{j}$, we defined the nonmonotonicity of $Q$ by $w(Q)=\widehat{w}_{j}(E)(Q)$. If $Q \in \mathcal{S}$ and $\angle\left(V_{Q}, V_{Q(\mathcal{S})}\right)>\lambda$, then there are horizontal lines that cross $V_{Q}$ in the positive direction (i.e., from $V_{Q}^{-}$to $V_{Q}^{+}$) but cross $V_{Q(\delta)}$ in the negative direction. We will show that these lines contribute to the total nonmonotonicity of the stopping time region $\mathcal{S}$, which is defined to be the quantity $w(\mathcal{S})=$ $\sum_{Q \in s} w(Q)$. We will then bound the number of cubes in $M_{1}$ in terms of the nonmonotonicity.

Our first goal will be to prove that most lines that cross $V_{Q}$ also cross $\partial E$. As in Section 9.4.1, we identify $\mathbb{R}^{2 k}$ with the subspace $\mathrm{H}=\left\{v \in \mathbb{H}^{2 k+1}\right.$ : $z(v)=0\}$ of horizontal vectors. We will use $\angle_{\mathrm{H}}$ to denote Euclidean angles between vectors and planes in H .

Lemma 78. For any $\lambda_{0}>0$ and any $0<a<b$, if $\varepsilon>0$ is sufficiently small, then the following property holds. Let $\mathcal{G}=\mathcal{L}(\varepsilon)$ be a set satisfying Proposition 75. Fix $Q \in \mathcal{L}_{\mathcal{L}}$, and let $x \in N_{1}(Q)$ be a point such that $d\left(x, V_{Q}\right)<$ $3 \varepsilon \sigma(Q)$. Also, let $u \in \mathrm{H}$ be a horizontal vector such that $\angle_{\mathrm{H}}\left(u, \pi\left(V_{Q}\right)\right)>\lambda_{0}$ and $u$ points toward $V_{Q}^{+}$. Let $L=x u^{\mathbb{R}}$ be the horizontal line through $x$ in the direction of $u$, and let $\gamma(t)$ be its unit-speed parametrization, with $\gamma(0)=x$. Then we have

$$
\gamma((-b \sigma(Q),-a \sigma(Q))) \subseteq E^{c} \quad \text { and } \quad \gamma((a \sigma(Q), b \sigma(Q))) \subseteq E
$$

Proof. suppose that $\varepsilon \in(0, \infty)$ satisfies

$$
\begin{equation*}
0<\varepsilon<\min \left\{\frac{1}{4}, \frac{a \sin \lambda_{0}}{4}, \frac{1}{b+1}\right\} . \tag{135}
\end{equation*}
$$

Let $f: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ be the signed distance to $V_{Q}$, i.e.,

$$
f(p)= \begin{cases}d\left(p, V_{Q}\right) & p \in V_{Q}^{+}, \\ -d\left(p, V_{Q}\right) & p \in V_{Q}^{-}\end{cases}
$$

This is a linear function on $\mathbb{H}^{2 k+1}$. If $f(p)>\varepsilon \sigma(Q)$, then $p \in V_{Q}^{+}$and $d\left(p, V_{Q}\right)>\varepsilon \sigma(Q)$. Hence, if also $p \in N_{1 / \varepsilon}(Q)$, then $p \in E$. Similarly, if $f(p)<-\varepsilon \sigma(Q)$ and $p \in N_{1 / \varepsilon}(Q)$, then $p \in E^{c}$.

For all $t \in \mathbb{R}$, we have $\pi(\gamma(t))=\pi(x)+t u$, so

$$
f(\gamma(t))=f(x)+t \sin \left(\angle_{\mathrm{H}}\left(u, \pi\left(V_{Q}\right)\right)\right)
$$

By hypothesis, $|f(x)|<3 \varepsilon \sigma(Q)$. Consequently, if $t \in(a \sigma(Q), b \sigma(Q))$, then

$$
f(\gamma(t))>-3 \varepsilon \sigma(Q)+a \sigma(Q) \sin \left(厶_{\mathrm{H}}\left(u, \pi\left(V_{Q}\right)\right)\right) \stackrel{(135)}{>} \varepsilon \sigma(Q)
$$

Since $\gamma(t) \in N_{1+t / \sigma(Q)}(Q) \subseteq N_{1 / \varepsilon}(Q)$, it follows that $\gamma(t) \in E$. Likewise, if $t \in(-b \sigma(Q),-a \sigma(Q))$, then $f(\gamma(t))<-\varepsilon \sigma(Q)$ and $\gamma(t) \in N_{1 / \varepsilon}(Q)$, so $\gamma(t) \in E^{c}$.

Consequently, each stopping-time region in $\mathscr{F}_{1}$ contributes to the nonmonotonicity of $E$.

Lemma 79. For any $0<\lambda<\frac{\pi}{2}$, there exists $\varepsilon>0$ such that if $\mathcal{L}_{\mathcal{L}}=\mathcal{L}_{\mathcal{L}}(\varepsilon)$ satisfies Proposition 75 and $\mathcal{S} \in \mathscr{F}_{1}$, then $w(\mathcal{S}) \gtrsim_{\lambda} \sigma(\mathrm{Q}(\mathcal{S}))^{2 k+1}$.

Proof. Let $0<c<1$ be as in Lemma 47, so that for all $Q \in \Delta$, there is a point $x_{Q} \in Q$ such that $B_{c \sigma(Q)}\left(x_{Q}\right) \cap \partial E \subseteq Q$. Let $0<\varepsilon<c / 4$ be a number satisfying Lemma 78 for $\lambda_{0}=\lambda / 4, a=c / 4$ and $b=2$. Note that $\varepsilon$ can be taken to be a function of $\lambda$ and $C$.

For each $Q \in \Delta$, let $\widehat{w}_{Q}$ be the restriction of $\widehat{w}_{\log _{2} \sigma(Q)}(E)$ to $Q$, i.e., the measure defined by $\widehat{w}_{Q}(A)=\widehat{w}_{\log _{2} \sigma(Q)}(E)(A \cap Q)$ for every measurable $A \subseteq \mathbb{H}^{2 k+1}$. For every $\mathcal{S} \in \mathcal{F}$, denote $\widehat{w}_{\mathcal{S}}=\sum_{Q \in \mathcal{S}} \widehat{w}_{Q}$. Then

$$
\begin{equation*}
\widehat{w}_{\mathcal{S}}\left(\mathbb{H}^{2 k+1}\right)=\sum_{Q \in \mathcal{S}} \widehat{w}_{\log _{2} \sigma(Q)}(E)(Q)=\sum_{Q \in \mathcal{S}} w(Q)=w(\mathcal{S}) \tag{136}
\end{equation*}
$$

Suppose that $A \in m_{1}(S)$, and let $Q$ be a child of $A$ such that $Q \in \mathcal{G}$ and $\theta \stackrel{\text { def }}{=} \angle\left(V_{Q}, V_{Q(S)}\right)>\lambda$. We claim that $\widehat{w}_{\mathcal{S}}(A) \gtrsim_{\lambda} \sigma(A)^{2 k+1}$. To see this, write $\sigma=\sigma(Q)$ and $\sigma_{\mathcal{S}}=\sigma(\mathrm{Q}(\mathcal{S}))$. Let $S^{2 k-1} \subseteq \mathrm{H}$ be the unit sphere in the space of horizontal vectors. Let $u_{0} \in S^{2 k-1}$ bisect $V_{Q}$ and $V_{Q(S)}$; i.e., let $u_{0} \in \mathrm{H}$ be the horizontal unit vector that points into $\pi\left(V_{Q}^{+}\right) \cap \pi\left(V_{\mathrm{Q}(S)}^{-}\right)$and satisfies

$$
\angle_{\mathrm{H}}\left(\pi\left(V_{Q}\right), \pi\left(u_{0}\right)\right)=\angle_{\mathrm{H}}\left(\pi\left(V_{\mathrm{Q}(S)}\right), \pi\left(u_{0}\right)\right)=\frac{\theta}{2}
$$

Denote $U=\left\{u \in S^{2 k-1}: \angle_{\mathrm{H}}\left(u, u_{0}\right)<\lambda / 4\right\}$ and $W=B_{\varepsilon \sigma}\left(x_{Q}\right)$. For $x \in W$ and $u \in U$, let $L(x, u)$ be the horizontal line through $x$ in direction $u$ and let $\gamma(t)=$ $\gamma_{x, u}(t)=x u^{t}$ be its unit-speed parametrization. Then $\angle_{\mathrm{H}}\left(u, \pi\left(V_{Q}\right)\right)>\frac{\lambda}{4}$ and $d\left(x, V_{Q}\right)<\varepsilon \sigma+d\left(x_{Q}, V_{Q}\right)<2 \varepsilon \sigma$, so $\gamma((-b \sigma,-a \sigma)) \subseteq E^{c}$ and $\gamma((a \sigma, b \sigma)) \subseteq E$ by Lemma 78. Similarly, since

$$
d\left(x, V_{\mathrm{Q}(\mathcal{S})}\right) \leqslant d\left(x, x_{Q}\right)+d\left(x_{Q}, V_{\mathrm{Q}(\mathcal{S})}\right)<\varepsilon \sigma+\varepsilon \sigma_{\mathcal{S}}<2 \varepsilon \sigma_{\mathcal{S}}
$$

Lemma 78 implies that $\gamma\left(\left(-b \sigma_{\mathcal{S}},-a \sigma_{\mathcal{S}}\right)\right) \subseteq E$ and $\gamma\left(\left(a \sigma_{\mathcal{S}}, b \sigma_{\mathcal{S}}\right)\right) \subseteq E^{c}$. Let $I=[p, q] \subseteq \mathbb{R}$ be the maximal interval such that $(a \sigma, b \sigma) \subseteq I$ and $\gamma(I) \subseteq E$.

Then $I \subseteq\left[-a \sigma, a \sigma_{S}\right]$, and therefore $\sigma<(b-a) \sigma \leqslant \ell(I) \leqslant 2 a \sigma_{S}<\sigma_{S}$. Furthermore, $p \in[-a \sigma, a \sigma]$, so $\gamma(p) \in \partial E$ and

$$
d\left(\gamma(p), x_{Q}\right) \leqslant a \sigma+d\left(x, x_{Q}\right) \leqslant \frac{c \sigma}{4}+\varepsilon \sigma \leqslant \frac{c \sigma}{2} .
$$

By our choice of $x_{Q}$, this implies $\gamma(p) \in Q$.
Denote $L_{W, U}=\{L(x, u): x \in W, u \in U\}$. For each $L \in L_{W, U}$, the intersection $E \cap L(x, u)$ contains a maximal interval $I$ of length $\ell(I) \in\left(\sigma, \sigma_{\mathcal{S}}\right)$ with one endpoint in $A$. It follows that

$$
\begin{equation*}
\widehat{w}_{\mathcal{S}}(A) \geqslant \sum_{i=\log _{2} \sigma}^{\log _{2} \sigma_{\mathcal{S}}} \widehat{w}_{i}(A) \geqslant \frac{\mathcal{L}\left(L_{W, U}\right)}{2} \asymp \varepsilon^{2 k+1} \sigma^{2 k+1} \lambda^{2 k-1} \asymp_{\lambda} \sigma(A)^{2 k+1} . \tag{137}
\end{equation*}
$$

Since $\mathcal{S} \in \mathscr{F}_{1}$ and all the elements of $m_{1}(\mathcal{S})$ are disjoint, we have

$$
\begin{aligned}
w(\mathcal{S}) \stackrel{(136)}{\geqslant} \widehat{w}_{\mathcal{S}}\left(\bigcup_{A \in m_{1}(\mathcal{S})} A\right) & \stackrel{(137)}{\gtrsim} \sum_{A \in m_{1}(\mathcal{S})} \sigma(A)^{2 k+1} \gtrsim \mathscr{H}^{2 k+1}\left(\bigcup_{A \in m_{1}(\mathcal{S})} A\right) \\
& \stackrel{134}{\gtrsim} \sigma(\mathrm{Q}(\mathcal{S}))^{2 k+1} .
\end{aligned}
$$

Finally, for all $R \in \Delta$, by Lemma 79 we have

$$
\sum_{\mathcal{S} \in \mathscr{H}_{1}(R)} \sigma(\mathrm{Q}(\mathcal{S}))^{2 k+1} \lesssim \lambda \sum_{\mathcal{S} \in \mathscr{\mathscr { F }}_{1}(R)} w(\mathcal{S}) \stackrel{(113)}{\leqslant} \operatorname{Per}(E)(R) \asymp \sigma(R)^{2 k+1},
$$

so $\mathscr{F}_{1}$ satisfies a Carleson packing condition. This concludes the proof of Proposition 77.
9.6. Stopping-time regions are close to Lipschitz graphs. Finally, we prove that the coronization constructed in Proposition 77 is in fact a corona decomposition $\partial E$ by constructing intrinsic Lipschitz graphs that approximate the stopping-time regions.

Proposition 80. For every $\eta, \theta \in(0,1)$, there are $\varepsilon>0$ and $\lambda \in(0,1)$ with the following property. If $\mathcal{S} \subseteq \Delta$ is a coherent set such that for all $Q \in \mathcal{S}$, there is a vertical half-space $V_{Q}^{+}$with bounding plane $V_{Q}$ such that

$$
\begin{equation*}
d_{N_{\frac{1}{\varepsilon}}(Q)}\left(V_{Q}^{+}, E\right) \leqslant \varepsilon \sigma(Q) \quad \text { and } \quad \angle\left(V_{Q}, V_{Q(s)}\right)<\lambda \tag{138}
\end{equation*}
$$

(e.g., a stopping-time region in the coronization constructed in Proposition 77), then there is an intrinsic Lipschitz graph $\Gamma \subseteq \mathbb{H}^{2 k+1}$ with Lipschitz constant at most $\eta$ such that for all $Q \in \mathcal{S}$,

$$
d_{N_{4}(Q)}\left(\Gamma^{+}, E\right) \leqslant \theta \sigma(Q)
$$

We choose $\lambda=\frac{\eta}{10}$. The constant $\varepsilon$ will depend on $\eta$ and $\theta$, but it will satisfy $\varepsilon<10^{-3}$. Let $M=\mathrm{Q}(\mathcal{S})$. We can apply an isometry of $\mathbb{H}^{2 k+1}$ so that $V_{M}=\left\{h: \mathbb{H}^{2 k+1}: x_{k}(h)=0\right\}$ and $V_{M}^{+}=\left\{h \in \mathbb{H}^{2 k+1}: x_{k}(h)>0\right\}$. Then
the one-parameter subgroup $\left\langle X_{k}\right\rangle$ is a horizontal line orthogonal to $V_{M}$. For $y \in \mathbb{H}^{2 k+1}$ we denote $L_{y}=y\left\langle X_{k}\right\rangle$. Let $\mathcal{L}^{\perp}=\left\{L_{x}: x \in \mathbb{H}^{2 k+1}\right\}$ be the set of cosets of $\left\langle X_{k}\right\rangle$ and, for any $U \subseteq \mathbb{H}^{2 k+1}$, let $\mathcal{L}^{\perp}(U)=\left\{L_{u}: u \in U\right\}$ be the set of cosets that pass through $U$. If $g: \mathcal{L}^{\perp}(U) \rightarrow \mathbb{R}$ is a function, we let $\bar{g}: U \rightarrow \mathbb{R}$ denote the corresponding function $\bar{g}(y)=g\left(L_{y}\right)$ that is constant on cosets.

If $V$ is a vertical plane such that $\angle\left(V, V_{M}\right)<\frac{\pi}{2}$, then it intersects cosets of $\left\langle X_{k}\right\rangle$ transversely, and we define $\Pi_{V}: \mathbb{H}^{2 k+1} \rightarrow V$ to be the projection whose fibers are the $L_{y}$ 's. That is, we let $\Pi_{V}(y)$ be the intersection point of $L_{y}$ and $V$. Let $\Pi=\Pi_{V_{M}}$ be the projection to $V_{M}$, i.e., $\Pi(y)=y X_{k}^{-x_{k}(y)}$.

One computes that if $V \subseteq \mathbb{H}^{2 k+1}$ is a vertical plane, then for every $v \in V$, $t \in \mathbb{R}$ and $y \in \mathbb{H}^{2 k+1}$,

$$
d\left(v X_{k}^{t}, V\right)=|t| \cos \left(\angle\left(V, V_{M}\right)\right) \quad \text { and } \quad d\left(y, \Pi_{V}(y)\right)=\frac{d(y, V)}{\cos \left(\angle\left(V, V_{M}\right)\right)}
$$

In particular, if $V=V_{Q}$ for some $Q \in S$, then $\angle\left(V, V_{S}\right) \leqslant \lambda$, and therefore

$$
\begin{equation*}
d\left(y, \Pi_{V}(y)\right) \leqslant 2 d(y, V) \tag{139}
\end{equation*}
$$

Recall that for $\rho>0$, we define $\rho Q=\partial E \cap N_{\rho}(Q)$.
Lemma $81 . \mathcal{L}^{\perp}\left(N_{9 / 2}(M)\right) \subseteq \mathcal{L}^{\perp}(10 M)$.
Proof. For all $h \in M,\left|x_{k}(h)\right| \leqslant \varepsilon \sigma(M)$, so every $v \in N_{9 / 2}(M)$ satisfies $\left|x_{k}(v)\right| \leqslant(9 / 2+\varepsilon) \sigma(M)$. Therefore $d(v, \Pi(v)) \leqslant 5 \sigma(M)$ and $\Pi(v) \in N_{19 / 2}(M)$. Since $\Pi(v) \in V_{M}$, by the first inequality in (138) we have $\Pi(v) X_{k}^{2 \varepsilon \sigma(M)} \in E$ and $\Pi(v) X_{k}^{-2 \varepsilon \sigma(M)} \notin E$. It follows that there is $q \in \partial E \cap L_{\Pi(v)}$ such that $d(\Pi(v), q) \leqslant 2 \varepsilon \sigma(M)$. Then $q \in 10 M$ and $L_{v}=L_{q} \in \mathcal{L}^{\perp}(10 M)$, as desired.

As in Chapter 8 of [30], we will construct $\Gamma$ using the function $d_{\mathcal{S}}$ that we defined in Section 6. The quantity $d_{\mathcal{S}}(v)$ measures the amount of control we have over $E$ near $v$.

Lemma 82. If $0<\alpha<1$, and if $\varepsilon>0$ is sufficiently small (depending on $\alpha$ and $C)$, then for every point $v \in \mathbb{H}^{2 k+1}$, and every $r \in\left(\alpha d_{\mathcal{S}}(v), \sigma(M) / \alpha\right]$, there exists $R \in \mathcal{S}$ and a vertical half-space $V^{+}=V_{R}^{+}$such that $d_{B_{r}(v)}\left(E, V^{+}\right)$ $\leqslant \varepsilon r / \alpha$.

Proof. Suppose that $0<\varepsilon<\alpha^{2} / 4$. Note that $\alpha d_{\mathcal{S}}(v)<\sigma(M) / \alpha$ by our assumption on $r$, so $d(v, M) \leqslant d_{\mathcal{S}}(v)<\sigma(M) / \alpha^{2}$ and thus $B_{r}(v) \in$ $N_{\alpha^{-2}+\alpha^{-1}}(M) \subseteq N_{\varepsilon^{-1}}(M)$.

We first consider the simpler case that $r \in[\alpha \sigma(M), \sigma(M) / \alpha]$. Let $V^{+}=$ $V_{M}^{+}$. Then

$$
d_{B_{r}(v)}\left(E, V^{+}\right) \leqslant d_{N_{\frac{1}{\varepsilon}}(M)}\left(E, V^{+}\right) \stackrel{(138)}{\leqslant} \varepsilon \sigma(M) \leqslant \frac{\varepsilon r}{\alpha}
$$

Otherwise, if $r<\alpha \sigma(M)$, then by the definition of $d_{\mathcal{S}}(v)$, there is a cube $Q \in \mathcal{S}$ such that $d(v, Q)+\sigma(Q)<r / \alpha$. Since $\sigma(Q)<r / \alpha<\sigma(M), Q$ has an ancestor $A$ (possibly $A=Q$ ) such that $r /(2 \alpha) \leqslant \sigma(A)<r / \alpha$. Then $d(v, A) \leqslant$ $d(v, Q)<r / \alpha \leqslant 2 \sigma(A)$, so $v \in N_{2}(A)$ and $B_{r}(v) \subseteq N_{2+2 \alpha}(A) \subseteq N_{1 / \varepsilon}(A)$. If $V^{+}=V_{A}^{+}$, then

$$
d_{B_{r}(v)}\left(E, V^{+}\right) \leqslant d_{N_{\frac{1}{\varepsilon}}(A)}\left(E, V^{+}\right) \stackrel{(138)}{\leqslant} \varepsilon \sigma(A)<\frac{\varepsilon r}{\alpha} .
$$

Lemma 82 implies the following lemma, which is a Heisenberg-version of Lemma 8.4 of [30], showing that $M$ satisfies the intrinsic Lipschitz condition on scales above $d_{\delta}$.

Lemma 83. If $0<\beta<1$ and if $\varepsilon>0$ is sufficiently small (depending on $\beta, \lambda$, and $C)$, then for all $v, w \in 20 M$ such that $d(v, w) \geqslant \beta \min \left\{d_{\mathcal{S}}(v), d_{\mathcal{S}}(w)\right\}$, we have

$$
\begin{equation*}
\left|x_{k}(v)-x_{k}(w)\right| \leqslant 2 \lambda d(v, w) \tag{140}
\end{equation*}
$$

Proof. Let $\alpha=\min \left\{\frac{1}{40+C}, \beta\right\}$, and let $0<\varepsilon<\frac{\lambda \alpha}{4}$ be a number satisfying Lemma 82.

If necessary, switch $v$ and $w$ so that $d_{\mathcal{S}}(v) \leqslant d_{\mathcal{S}}(w)$. Let $r=d(v, w)$. Then $r \leqslant \operatorname{diam}\left(N_{20}(M)\right) \leqslant(40+C) \sigma(M)$ and $r \geqslant \beta d_{\mathcal{S}}(v)$, so $r \in\left(\alpha d_{\mathcal{S}}(v), \sigma(M) / \alpha\right]$ and by Lemma 82 , there is a vertical half-space $V^{+}$such that $d_{B_{r}(v)}\left(E, V^{+}\right) \leqslant$ $\varepsilon r / \alpha$. Consequently, there are points $v^{\prime}, w^{\prime} \in V$ such that $d\left(v, v^{\prime}\right) \leqslant \lambda r / 4$ and $d\left(w, w^{\prime}\right) \leqslant \lambda r / 4$. In particular, $d\left(v^{\prime}, w^{\prime}\right) \leqslant 3 r / 2$.

Since $V^{+}=V_{R}^{+}$for some $R \in \mathcal{S}$, we have $\angle\left(V, V_{S}\right)<\lambda$, so $\mid x_{k}\left(v^{\prime}\right)-$ $x_{k}\left(w^{\prime}\right) \mid \leqslant \lambda d\left(v^{\prime}, w^{\prime}\right) \leqslant 3 \lambda r / 2$ and $\left|x_{k}(v)-x_{k}(w)\right| \leqslant \frac{\lambda}{2} r+\left|x_{k}\left(v^{\prime}\right)-x_{k}\left(w^{\prime}\right)\right| \leqslant$ $2 \lambda d(v, w)$.

In particular, if $\Gamma_{0}=d_{\delta}^{-1}(0)$, then $\Gamma_{0}$ is an intrinsic Lipschitz graph over $\Pi\left(\Gamma_{0}\right)$.

Corollary 84. If $w \in \Gamma_{0}$ and $v \in 20 M$, then $\left|x_{k}(v)-x_{k}(w)\right| \leqslant 2 \lambda d(v, w)$ and $L_{w} \cap 20 M=\{w\}$.

Proof. If $v \in 20 M$, then Lemma 83 implies that $\left|x_{k}(v)-x_{k}(w)\right| \leqslant 2 \lambda d(v, w)$. If $v \in L_{w} \cap 20 M$, then $d(v, w)=\left|x_{k}(v)-x_{k}(w)\right| \leqslant 2 \lambda d(v, w)$, so $d(v, w)=0$ and $v=w$.

We will extend $\Gamma_{0}$ and construct the rest of $\Gamma$ by smoothing $M$ using a partition of unity. For the rest of this section, we will take $\alpha=\beta=10^{-3}$ and assume that $\varepsilon$ is small enough that Lemmas 82 and 83 hold.

Let $\psi=1 / 50$ and let $\left\{c_{i}\right\}_{i \in I} \subseteq 10 M$ be as in Lemma 54 , so that the balls $\left\{B_{i}=B_{\psi d_{s}\left(c_{i}\right)}\left(c_{i}\right)\right\}_{i \in I}$ are disjoint, and so that if $4 \leqslant \rho<50$, then the balls $\left\{\rho B_{i}=B_{\rho \psi d_{\mathcal{S}}\left(c_{i}\right)}\left(c_{i}\right)\right\}_{i \in I}$ cover $10 M \backslash \Gamma_{0}$ with bounded multiplicity. Since $d_{\mathcal{S}}$
is 1-Lipschitz, we have

$$
\forall i \in I, \quad q \in \rho B_{i} \Longrightarrow\left(1-\frac{\rho}{50}\right) d_{\delta}\left(c_{i}\right) \leqslant d_{\mathcal{S}}(q) \leqslant\left(1+\frac{\rho}{50}\right) d_{\mathcal{S}}\left(c_{i}\right)
$$

For each $i \in I$, we have $d_{\mathcal{S}}\left(c_{i}\right) \in\left(\alpha d_{\mathcal{S}}\left(c_{i}\right), \sigma(M) / \alpha\right]$, so by Lemma 82 there is a vertical plane $V_{i}$ such that

$$
\begin{equation*}
d_{B_{d_{\mathcal{S}}\left(c_{i}\right)}\left(c_{i}\right)}\left(V_{i}^{+}, E\right)=d_{50 B_{i}}\left(V_{i}^{+}, E\right) \leqslant \frac{\varepsilon d_{\mathcal{S}}\left(c_{i}\right)}{\alpha} . \tag{141}
\end{equation*}
$$

Lemma 85. For all $L \in \mathcal{L}^{\perp}(10 M)$, let $I_{L}=\left\{i \in I: 8 B_{i} \cap L \neq \varnothing\right\}$, and let $\delta_{L}=\inf _{y \in L \cap 20 M} d_{\delta}(y)$. If $\varepsilon>0$ is sufficiently small, then
(1) there is an interval $U_{L} \subseteq L$ such that $\ell\left(U_{L}\right) \lesssim \varepsilon \delta_{L}$ and such that $20 M \cap L \subseteq$ $U_{L}$ and $\bigcup_{i \in I_{L}} L \cap V_{i} \subseteq U_{L}$;
(2) if $L \in \mathcal{L}^{\perp}(10 M), q \in L \cap 20 M$ and $i \in I_{L}$, then $\delta_{L} \asymp d_{\mathcal{S}}(q) \asymp d_{\mathcal{S}}\left(c_{i}\right)$;
(3) $\left|I_{L}\right| \asymp 1$.

Proof. Let $q \in L \cap 20 M$. We claim that $L \cap 20 M$ is contained in an interval of radius $\beta d_{\mathcal{S}}(q)$ centered at $q$. If $p \in L \cap 20 M$, then Lemma 83 implies that either $d(p, q)<\beta \min \left\{d_{\mathcal{S}}(p), d_{\mathcal{S}}(q)\right\}$ or $\left|x_{k}(p)-x_{k}(q)\right| \leqslant 2 \lambda d(p, q)<d(p, q)$. But $p$ and $q$ both lie on $L$, so $d(p, q)=\left|x_{k}(p)-x_{k}(q)\right|$, and thus $p \in B_{\beta d_{s}(q)}(q)$. Consequently, $d_{\delta}(p) \geqslant d_{\delta}(q)-\beta d_{\mathcal{S}}(q)$; since this holds for all $p, q \in L \cap 20 M$, it implies that $\delta_{L} \asymp d_{\delta}(q)$ for all $q \in L \cap 20 M$. In fact, we can prove that $L \cap 20 M$ actually lies in a smaller interval. Suppose that $p, q \in L \cap 20 M$. Let $i \in I$ be such that $q \in 8 B_{i}$. Since $d_{\mathcal{S}}$ is 1 -Lipschitz, we have $d_{\mathcal{S}}(q)<2 d_{\mathcal{S}}\left(c_{i}\right)$ and thus $p \in 9 B_{i}$. Denoting $\{v\}=\Pi_{V_{i}}(L)=L \cap V_{i}$, we have

$$
d(p, v) \stackrel{(139)}{\lessgtr} 2 d\left(p, V_{i}\right) \stackrel{(138)}{\lesssim} \varepsilon d_{\delta}\left(c_{i}\right) \asymp \varepsilon d_{\delta}(q) \asymp \varepsilon \delta_{L} .
$$

Similarly, $d(q, v) \lesssim \varepsilon \delta_{L}$. Hence, $\operatorname{diam}(L \cap 20 M) \asymp \varepsilon \delta_{L}$. This proves the first part of condition (1).

To prove the second part of condition (1), suppose that $L \in \mathcal{L}^{\perp}(10 M)$ and $i \in I_{L}$. Let $y_{i} \in 8 B_{i} \cap L$, and let $\left\{v_{i}\right\}=\Pi_{V_{i}}(L)$ as before. Then $d\left(c_{i}, V_{i}\right) \lesssim$ $\varepsilon d_{\mathcal{S}}\left(c_{i}\right)$ and

$$
d\left(v_{i}, y_{i}\right) \leqslant 2 d\left(y_{i}, V_{i}\right) \leqslant 2 d\left(y_{i}, c_{i}\right)+2 d\left(c_{i}, V_{i}\right) \stackrel{(141)}{\leqslant} 2 d\left(y_{i}, c_{i}\right)+2 \varepsilon d_{\mathcal{S}}\left(c_{i}\right) .
$$

If $\varepsilon$ is sufficiently small, then it follows that $v_{i} \in 17 B_{i}$, and by (141), there is some $q_{i} \in L \cap \partial E \cap 20 B_{i}$ such that $d\left(q_{i}, v_{i}\right) \lesssim \varepsilon d_{\mathcal{S}}\left(c_{i}\right)$. It follows that $\delta_{L} \asymp$ $d_{\mathcal{S}}\left(q_{i}\right) \asymp d_{\mathcal{S}}\left(c_{i}\right)$ and $q_{i} \in 20 M$, so we have $d\left(v_{i}, L \cap 20 M\right) \lesssim \varepsilon d_{\mathcal{S}}\left(c_{i}\right)$ and the second part of condition (1) holds, and also condition (2) holds. Furthermore if also $j \in I_{L}$ then $d\left(q_{i}, q_{j}\right) \lesssim \varepsilon \delta_{L} \asymp \varepsilon d_{\mathcal{S}}\left(c_{i}\right)$ by conditions (1) and (2). So, $q_{i} \in$ $21 B_{j}$, and since the balls $\left\{20 B_{j}\right\}_{j \in I_{L}}$ form a cover with bounded multiplicity, this implies that $\left|I_{L}\right| \asymp 1$.

Next, we define some auxiliary bump functions. Let $\mathfrak{b}: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ be a bump function such that $\mathfrak{b}\left(B_{5 \psi}\right)=1, \mathfrak{b}=0$ outside $B_{6 \psi}$ and such that $\mathfrak{b}$ is smooth and Lipschitz with respect to the Euclidean metric. That is, $\mathfrak{b}$ is a smooth function considered as a function on $\mathbb{R}^{2 k+1}$ and $\|\mathcal{G}\|_{\operatorname{Lip}\left(\ell_{2}^{2 k+1}\right)} \asymp 1$. This condition is needed below because nonhorizontal derivatives of $b$ will appear in some bounds. We denote the Euclidean metric on $\mathbb{H}^{2 k+1}$ by $d_{\text {Euc }}$, that is, $d_{\text {Euc }}\left(v, v^{\prime}\right)=\left\|v-v^{\prime}\right\|$ for all $v, v^{\prime} \in \mathbb{R}^{2 k+1}$. Note that if $\rho>0$ and $\gamma$ is a horizontal path inside $B_{\rho}$, then $\ell(\gamma) \asymp_{\rho} \ell_{\mathrm{Euc}}(\gamma)$. Therefore, if $v, v^{\prime} \in B_{\rho}$, then $d_{\text {Euc }}\left(v, v^{\prime}\right) \lesssim_{\rho} d\left(v, v^{\prime}\right)$. Hence $b$ is also Lipschitz with respect to the CarnotCarathéodory metric and $\|6\|_{\operatorname{Lip}\left(\mathbb{H}^{2 k+1}\right)} \asymp 1$.

For each $i \in I$, define $\tau_{i}: \mathbb{H}^{2 k+1} \rightarrow \mathbb{R}$ by

$$
\forall y \in \mathbb{H}^{2 k+1}, \quad \tau_{i}(y) \stackrel{\text { def }}{=} d_{\delta}\left(c_{i}\right) \mathcal{G}\left(\mathfrak{s}_{\frac{1}{d_{\mathcal{L}}\left(c_{i}\right)}}\left(c_{i}^{-1} y\right)\right) .
$$

The scaling and translation in the above definition sends $B_{i}$ to $B_{\psi}$, so $\tau_{i}=$ $d_{\mathcal{S}}\left(c_{i}\right)$ on $5 B_{i}, \tau_{i}$ vanishes outside the ball $6 B_{i}$, and $\left\|\tau_{i}\right\|_{\operatorname{Lip}\left(\mathbb{H}^{2 k+1}\right)} \asymp 1$.

We will use the $\left\{\tau_{i}\right\}_{i \in I}$ to construct a partition of unity on $\mathcal{L}^{\perp}$. For each $i \in I$, we define $\tau_{i}^{\prime}: \mathcal{L}^{\perp} \rightarrow \mathbb{R}$ by $\tau_{i}^{\prime}(L)=\tau_{i}\left(L \cap V_{i}\right)$. Equivalently, $\bar{\tau}_{i}^{\prime}(p)=\tau_{i}\left(\Pi_{V_{i}}(p)\right)$. Then $\tau_{i}^{\prime}$ is a continuous function that is zero outside $\mathcal{L}^{\perp}\left(V_{i} \cap 6 B_{i}\right)$.

Lemma 86. Let $T \stackrel{\text { def }}{=} \sum_{i \in I} \tau_{i}^{\prime}$. Then the following assertions hold true:
(1) For all $q \in 10 M$, we have $\bar{T}(q) \asymp d_{\delta}(q)$.
(2) For all $i \in I$, we have $\left\|\bar{\tau}_{i}^{\prime}\right\|_{\operatorname{Lip}\left(8 B_{i}\right)} \lesssim 1$.

Hence, the functions $\left\{\phi_{i} \stackrel{\text { def }}{=} \frac{\tau_{i}^{\prime}}{T}\right\}_{i \in I}$ form a partition of unity on $\mathcal{L}^{\perp}(10 M)$.
Proof. First, we show that $\bar{T}(q) \gtrsim d_{\delta}(q)$ for all $q \in 10 M$. If $q \in \Gamma_{0}$, then $d_{\mathcal{S}}(q)=0$ and there is nothing to prove. So, suppose that $q \in 10 M \backslash \Gamma_{0}$. Let $i \in I$ be such that $q \in 4 B_{i}$. Using (141), if $\varepsilon<\alpha / 2$, then $d\left(q, V_{i}\right) \leqslant d_{\delta}\left(c_{i}\right) / 2$. By (139), $d\left(\Pi_{V_{i}}(q), c_{i}\right) \leqslant 2 d\left(q, V_{i}\right)+d\left(q, c_{i}\right)<5 \psi d_{\mathcal{S}}\left(c_{i}\right)$. Thus $\Pi_{V_{i}}(q) \in 5 B_{i}$ and $\bar{T}(q) \geqslant \tau_{i}\left(\Pi_{V_{i}}(q)\right)=d_{\mathcal{S}}\left(c_{i}\right) \asymp d_{\delta}(q)$. On the other hand, Lemma 85 implies that for any $L \in \mathcal{L}^{\perp}(10 M)$, the set $I_{L}=\left\{i: 8 B_{i} \cap L \neq \varnothing\right\}$ has boundedly many elements, and if $q \in L \cap 10 M$, then $d_{\mathcal{S}}(q) \asymp d_{\mathcal{S}}\left(c_{i}\right)$ for each $i \in I_{L}$. Thus $\bar{T}(q) \lesssim d_{\mathcal{S}}(q)\left|I_{L}\right| \lesssim d_{\delta}(q)$. This proves the first part of the Lemma 86.

It remains to show that if $w_{1}, w_{2} \in 8 B_{i}$, then $\left|\tau_{i}\left(\Pi_{V_{i}}\left(w_{1}\right)\right)-\tau_{i}\left(\Pi_{V_{i}}\left(w_{2}\right)\right)\right| \lesssim$ $d\left(w_{1}, w_{2}\right)$. The main difficulty is that the map $\Pi_{V_{i}}$ is not Lipschitz with respect to the Carnot-Carathéodory metric; it is, however, smooth with respect to the Euclidean structure on $\mathbb{R}^{2 k+1}$.

Let $\xi: \mathbb{H}^{2 k+1} \rightarrow \mathbb{H}^{2 k+1}$ be the map $\xi(y)=\mathfrak{s}_{d_{S}\left(c_{i}\right)^{-1}}\left(c_{i}^{-1} y\right)$ that rescales and translates $B_{i}$ to $B_{\psi}$. This map sends lines in $\mathcal{L}^{\perp}$ to lines in $\mathcal{L}^{\perp}$, so
$\xi \circ \Pi_{V_{i}}=\Pi_{\xi\left(V_{i}\right)} \circ \xi$ and

$$
\begin{equation*}
\forall w \in \mathbb{H}^{2 k+1}, \quad \bar{\tau}_{i}^{\prime}(w)=d_{\delta}\left(c_{i}\right) \mathcal{b}\left(\xi\left(\Pi_{V_{i}}(w)\right)\right)=d_{\delta}\left(c_{i}\right) \mathfrak{b}\left(\Pi_{\xi\left(V_{i}\right)}(\xi(w))\right) . \tag{142}
\end{equation*}
$$

The projection $\Pi_{\xi\left(V_{i}\right)}$ is a smooth map in Euclidean coordinates. In fact, since $\xi\left(V_{i}\right)$ is a vertical plane with $\angle\left(V_{M}, \xi\left(V_{i}\right)\right) \leqslant \lambda$ and

$$
d\left(\mathbf{0}, \xi\left(V_{i}\right)\right)=\frac{d\left(c_{i}, V_{i}\right)}{d_{\delta}\left(c_{i}\right)} \lesssim \varepsilon,
$$

the Euclidean derivatives of $\Pi_{\xi\left(V_{i}\right)}$ are uniformly bounded inside the ball $B_{1}$, and hence the Lipschitz constant with respect to the Euclidean metric of the restriction of $\Pi_{\tilde{\varepsilon}\left(V_{i}\right)}$ to $B_{1}$ is bounded. That is, for $v_{1}, v_{2} \in B_{1}$, we have $d_{\text {Euc }}\left(\Pi_{\xi\left(V_{i}\right)}\left(v_{1}\right), \Pi_{\xi\left(V_{i}\right)}\left(v_{2}\right)\right) \lesssim d_{\text {Euc }}\left(v_{1}, v_{2}\right)$. Because $w_{1}, w_{2} \in 8 B_{i}$, we may apply this with $v_{1}=\xi\left(w_{1}\right), v_{2}=\xi\left(w_{2}\right) \in B_{1}$ to deduce that

$$
\begin{align*}
d_{\operatorname{Euc}}\left(\Pi_{\xi\left(V_{i}\right)}\left(\xi\left(w_{1}\right)\right), \Pi_{\xi\left(V_{i}\right)}\left(\xi\left(w_{2}\right)\right)\right) & \lesssim d_{\operatorname{Euc}}\left(\xi\left(w_{1}\right), \xi\left(w_{2}\right)\right) \\
& \lesssim d\left(\xi\left(w_{1}\right), \xi\left(w_{2}\right)\right)=\frac{d\left(w_{1}, w_{2}\right)}{d_{\delta}\left(c_{i}\right)} . \tag{143}
\end{align*}
$$

Consequently,

$$
\left|\bar{\tau}_{i}^{\prime}\left(w_{1}\right)-\bar{\tau}_{i}^{\prime}\left(w_{2}\right)\right| \stackrel{(142) \wedge(143)}{\leqslant} d_{\delta}\left(c_{i}\right)\|6\|_{\operatorname{Lip}\left(\ell_{2}^{2 k+1}\right)} \cdot \frac{d\left(w_{1}, w_{2}\right)}{d_{\mathcal{S}}\left(c_{i}\right)} \lesssim d\left(w_{1}, w_{2}\right)
$$

For each $i \in I$, let $f_{i}: \mathcal{L}^{\perp} \rightarrow \mathbb{R}$ be the function $f_{i}(L)=x_{k}\left(L \cap V_{i}\right)$. Then $\bar{f}_{i}(p)=x_{k}\left(\Pi_{V_{i}}(p)\right)$, so $\bar{f}_{i}$ is constant on vertical lines and

$$
\begin{equation*}
\left\|\bar{f}_{i}\right\|_{\operatorname{Lip}\left(\mathbb{H}^{2 k+1}\right)} \leqslant 2 \lambda . \tag{144}
\end{equation*}
$$

We define $f$ on $\mathcal{L}^{\perp}(10 M)$ by

$$
f(L)= \begin{cases}x_{k}\left(L \cap \Gamma_{0}\right) & \text { if } L \in \mathcal{L}^{\perp}\left(\Gamma_{0}\right) \\ \sum_{i \in I} \phi_{i}(L) f_{i}(L) & \text { if } L \in \mathcal{L}^{\perp}\left(10 M \backslash \Gamma_{0}\right) .\end{cases}
$$

By Corollary 84, this function is well defined. Let

$$
\Gamma \stackrel{\text { def }}{=}\left\{w \in \mathbb{H}^{2 k+1}: x_{k}(w)=\bar{f}(w)\right\}=\left\{v X_{k}^{\bar{f}(v)}: v \in V_{\delta}\right\} .
$$

This is an intrinsic graph over $\Pi_{V_{S}}(10 M)$.
Lemma 87. For all sufficiently small $\varepsilon>0$, the following assertions hold true:
(1) If $i \in I, q \in 10 M \cap 8 B_{i}$ and $w \in L_{q} \cap \Gamma$, then

$$
\begin{align*}
\left|x_{k}(q)-\bar{f}_{i}(q)\right| & \lesssim \varepsilon d_{\delta}(q),  \tag{145}\\
\left|x_{k}(q)-\bar{f}(q)\right| & =d(q, w) \lesssim \varepsilon d_{\delta}(q),  \tag{146}\\
d_{\delta}(w) & \asymp d_{\delta}(q) . \tag{147}
\end{align*}
$$

(2) $\Gamma$ is an intrinsic Lipschitz graph over $10 M$ with Lipschitz constant at most 10入.

Proof. If $q \in \Gamma_{0}$, then $w=q$, so (145) holds vacuously and (146) and (147) are trivial. So, suppose that $q \in 10 M \backslash \Gamma_{0}$ and that $w$ and $i$ are as above. By part (1) of Lemma 85, we have

$$
\left|x_{k}(q)-\bar{f}_{i}(q)\right|=d\left(q, \Pi_{V_{i}}(q)\right) \lesssim \varepsilon d_{\mathcal{S}}(q) .
$$

Since $f$ is an average of these $f_{i}$ 's, this implies (146). If $\varepsilon$ is small enough, then $d(q, w) \leqslant d_{\delta}(q) / 2$, so $d_{\delta}(w) \asymp d_{\mathcal{S}}(q)$. This proves the first part of Lemma 87.

To prove that $\Gamma$ is a Lipschitz graph, let $q_{1}, q_{2} \in 10 M$ and let $\left\{w_{j}\right\}=$ $L_{q_{j}} \cap \Gamma$. By the first part of the lemma, for $j=1,2$, we have $d\left(q_{j}, w_{j}\right) \lesssim \varepsilon d_{\delta}\left(q_{j}\right)$. We suppose that $\varepsilon$ is small enough that $d\left(q_{j}, w_{j}\right) \leqslant 10^{-3} \lambda d_{\mathcal{S}}\left(q_{j}\right)$. Since $d_{\mathcal{S}}$ is 1 -Lipschitz, this implies $d_{\mathcal{S}}\left(q_{j}\right) / 2 \leqslant d_{\mathcal{S}}\left(w_{j}\right) \leqslant 3 d_{\mathcal{S}}\left(q_{j}\right) / 2$.

We claim that $w_{1}$ and $w_{2}$ satisfy $\left|\bar{f}\left(w_{1}\right)-\bar{f}\left(w_{2}\right)\right|=\left|x_{k}\left(w_{1}\right)-x_{k}\left(w_{2}\right)\right| \leqslant$ $10 \lambda d\left(w_{1}, w_{2}\right)$. We consider two cases depending on whether $d\left(w_{1}, w_{2}\right)$ is greater or less than $10^{-2} \min \left\{d_{\mathcal{S}}\left(w_{1}\right), d_{\mathcal{S}}\left(w_{2}\right)\right\}$.

Suppose that $d_{\mathcal{S}}\left(w_{1}\right) \leqslant d_{\mathcal{S}}\left(w_{2}\right)$ and that $10^{-2} d_{\mathcal{S}}\left(w_{1}\right)<d\left(w_{1}, w_{2}\right)$. Then

$$
d_{\mathcal{S}}\left(w_{1}\right) \leqslant d_{\mathcal{S}}\left(w_{2}\right) \leqslant d\left(w_{1}, w_{2}\right)+d_{\mathcal{S}}\left(w_{1}\right) \leqslant 101 d\left(w_{1}, w_{2}\right),
$$

and also $d\left(q_{j}, w_{j}\right) \leqslant 2 \lambda \cdot 10^{-3} d_{\mathcal{S}}\left(w_{j}\right) \leqslant \lambda d\left(w_{1}, w_{2}\right) / 4$. It follows that

$$
\begin{equation*}
\frac{d\left(w_{1}, w_{2}\right)}{2} \leqslant d\left(q_{1}, q_{2}\right) \leqslant 2 d\left(w_{1}, w_{2}\right) \tag{148}
\end{equation*}
$$

and therefore $d\left(q_{1}, q_{2}\right)>\frac{d_{\delta}\left(w_{1}\right)}{200} \geqslant \frac{d_{\delta}\left(q_{1}\right)}{400} \geqslant \beta d_{\delta}\left(q_{1}\right)$. By Lemma 83 applied to $q_{1}$ and $q_{2}$, we find

$$
\begin{aligned}
\left|x_{k}\left(w_{1}\right)-x_{k}\left(w_{2}\right)\right| \leqslant \mid x_{k}\left(q_{1}\right) & -x_{k}\left(q_{2}\right) \mid+d\left(q_{1}, w_{1}\right)+d\left(q_{2}, w_{2}\right) \\
& \stackrel{(140)}{\leqslant} 2 \lambda d\left(q_{1}, q_{2}\right)+\lambda d\left(w_{1}, w_{2}\right) \stackrel{(148)}{\leqslant} 5 \lambda d\left(w_{1}, w_{2}\right),
\end{aligned}
$$

as desired.
Now suppose $d\left(w_{1}, w_{2}\right)<10^{-2} \min \left\{d_{\mathcal{S}}\left(w_{1}\right), d_{\mathcal{S}}\left(w_{2}\right)\right\}$. Then $w_{1}, w_{2} \notin \Gamma_{0}$ and $d_{\delta}\left(w_{1}\right) \asymp d_{\delta}\left(w_{2}\right)$; indeed, $d_{\delta}\left(w_{1}\right) \in\left[0.99 d_{\delta}\left(w_{2}\right), 1.01 d_{\delta}\left(w_{2}\right)\right]$. Let $\delta=$ $d_{\delta}\left(w_{1}\right)$, so that $\delta \asymp d_{\delta}\left(w_{j}\right) \asymp d_{\delta}\left(q_{j}\right)$.

We bound $\left|\bar{f}\left(w_{1}\right)-\bar{f}\left(w_{2}\right)\right|$ by bounding the behavior of the partition of unity. Suppose $\bar{\tau}_{i}^{\prime}\left(w_{1}\right)>0$. We claim that $w_{1}, w_{2} \in 8 B_{i}$. Since $\bar{\tau}_{i}^{\prime}\left(w_{1}\right)=$ $\tau_{i}\left(\Pi_{V_{i}}\left(w_{1}\right)\right) \neq 0$, we have $\Pi_{V_{i}}\left(w_{1}\right) \in 6 B_{i}$. Then

$$
d\left(w_{1}, \Pi_{V_{i}}\left(w_{1}\right)\right)=\left|\bar{f}\left(w_{1}\right)-\bar{f}_{i}\left(w_{1}\right)\right| \stackrel{(145) \wedge(146)}{\lesssim} \varepsilon d_{\delta}\left(q_{1}\right) .
$$

By part (2) of Lemma 85, we have $d_{\delta}\left(q_{1}\right) \asymp d_{\delta}\left(c_{i}\right)$, so if $\varepsilon$ is small enough, then $w_{1} \in 7 B_{i}$. Consequently, $d_{\mathcal{S}}\left(w_{1}\right) \leqslant 2 d_{\mathcal{S}}\left(c_{i}\right)$ and

$$
d\left(w_{1}, w_{2}\right) \leqslant \frac{d_{\mathcal{\delta}}\left(w_{1}\right)}{100} \leqslant \frac{d_{\mathcal{\delta}}\left(c_{i}\right)}{50},
$$

so $w_{2} \in 8 B_{i}$. Likewise, if $\bar{\tau}_{i}^{\prime}\left(w_{2}\right)>0$, then $w_{1}, w_{2} \in 8 B_{i}$.
Therefore, there are only boundedly many $i \in I$ such that $\bar{\tau}_{i}^{\prime}\left(w_{1}\right) \neq \bar{\tau}_{i}^{\prime}\left(w_{2}\right)$, and by Lemma 86 ,

$$
\left|\bar{\tau}_{i}^{\prime}\left(w_{1}\right)-\bar{\tau}_{i}^{\prime}\left(w_{2}\right)\right| \lesssim d\left(w_{1}, w_{2}\right)
$$

for any such $i$. It follows that $\left|\bar{T}\left(w_{1}\right)-\bar{T}\left(w_{2}\right)\right| \lesssim d\left(w_{1}, w_{2}\right)$.
Let $\Delta \bar{\phi}_{i}=\bar{\phi}_{i}\left(w_{1}\right)-\bar{\phi}_{i}\left(w_{2}\right)$ and $\Delta \bar{f}_{i}=\bar{f}_{i}\left(w_{1}\right)-\bar{f}_{i}\left(w_{2}\right)$. Since $\bar{T}\left(w_{i}\right) \asymp$ $d_{\delta}\left(w_{i}\right) \asymp \delta$,

$$
\begin{aligned}
\left|\Delta \bar{\phi}_{i}\right| & =\frac{\left|\bar{\tau}_{i}^{\prime}\left(w_{1}\right) \bar{T}\left(w_{2}\right)-\bar{\tau}_{i}^{\prime}\left(w_{2}\right) \bar{T}\left(w_{1}\right)\right|}{\bar{T}\left(w_{1}\right) \bar{T}\left(w_{2}\right)} \\
& \lesssim \delta^{-2}\left(\left|\bar{\tau}_{i}^{\prime}\left(w_{1}\right) \bar{T}\left(w_{2}\right)-\bar{\tau}_{i}^{\prime}\left(w_{1}\right) \bar{T}\left(w_{1}\right)\right|+\left|\bar{\tau}_{i}^{\prime}\left(w_{1}\right) \bar{T}\left(w_{1}\right)-\bar{\tau}_{i}^{\prime}\left(w_{2}\right) \bar{T}\left(w_{1}\right)\right|\right) \\
& \lesssim \delta^{-2}\left(\bar{\tau}_{i}^{\prime}\left(w_{1}\right) d\left(w_{1}, w_{2}\right)+\bar{T}\left(w_{1}\right) d\left(w_{1}, w_{2}\right)\right) \\
& \lesssim \frac{d\left(w_{1}, w_{2}\right)}{\delta} .
\end{aligned}
$$

By (144), we have $\left|\Delta \bar{f}_{i}\right| \leqslant 2 \lambda d\left(w_{1}, w_{2}\right)$. Hence,

$$
\begin{align*}
\left|\bar{f}\left(w_{2}\right)-\bar{f}\left(w_{1}\right)\right|= & \left|\sum_{i \in I} \bar{\phi}_{i}\left(w_{1}\right) \bar{f}_{i}\left(w_{1}\right)-\sum_{i \in I} \bar{\phi}_{i}\left(w_{2}\right) \bar{f}_{i}\left(w_{2}\right)\right| \\
\leqslant & \left|\sum_{i \in I} \Delta \bar{\phi}_{i} \cdot \bar{f}_{i}\left(w_{1}\right)\right|+\left|\sum_{i \in I} \bar{\phi}_{i}\left(w_{2}\right) \cdot \Delta \bar{f}_{i}\right| \\
\leqslant & \bar{f}\left(w_{1}\right)\left|\sum_{i \in I} \Delta \bar{\phi}_{i}\right|  \tag{149}\\
& +\sum_{i \in I}\left|\Delta \bar{\phi}_{i}\right| \cdot\left|\bar{f}_{i}\left(w_{1}\right)-\bar{f}\left(w_{1}\right)\right|+\sum_{i \in I} \bar{\phi}_{i}\left(w_{2}\right)\left|\Delta \bar{f}_{i}\right| \\
\leqslant & 0+\sum_{i \in I} O\left(\frac{d\left(w_{1}, w_{2}\right)}{\delta} \cdot \varepsilon \delta\right)+2 \lambda d\left(w_{1}, w_{2}\right) \sum_{i \in I} \bar{\phi}_{i}\left(w_{2}\right) \\
\leqslant & O\left(\varepsilon d\left(w_{1}, w_{2}\right)\right)+\lambda d\left(w_{1}, w_{2}\right)
\end{align*}
$$

where in (149) we used (145) and (146) to bound $\left|\bar{f}_{i}-\bar{f}\right|$. If $\varepsilon$ is sufficiently small, then this implies the desired estimate $\left|\bar{f}\left(w_{1}\right)-\bar{f}\left(w_{2}\right)\right| \leqslant 2 \lambda d\left(w_{1}, w_{2}\right)$.

Since $\Gamma$ is an intrinsic Lipschitz graph over $\Pi(10 M)$, Theorem 27 implies that $\Gamma$ can be extended to an intrinsic Lipschitz graph over $V_{M}$ with no increase in the Lipschitz constant. We also denote this extension by $\Gamma$, and we let $\Gamma^{+}$ be the half-space $\Gamma=\left\{w \in \mathbb{H}^{2 k+1}: x_{k}(w) \geqslant \bar{f}(w)\right\}$.

Finally, we claim that this satisfies (87). Let $Q \in \mathcal{S}$. By (128), we have $d_{N_{1 / \varepsilon}(Q)}\left(V_{Q}^{+}, E\right) \leqslant \varepsilon \sigma(Q)$. We claim that

$$
\begin{equation*}
d_{N_{\frac{9}{2}}(Q)}\left(V_{Q}^{+}, \Gamma^{+}\right) \lesssim \varepsilon \sigma(Q) \tag{150}
\end{equation*}
$$

Let $L \in \mathcal{L}^{\perp}\left(N_{9 / 2}(Q)\right)$. Suppose that $v \in L \cap V_{Q}$ and $w \in L \cap \Gamma$. By Lemma 81, there exists $q \in L \cap 10 Q$. By (128) and (139), we have

$$
\left|x_{k}(q)-x_{k}(v)\right|=d(q, v) \leqslant 2 d\left(q, V_{Q}\right) \leqslant 2 \varepsilon \sigma(Q)
$$

By (146), we have

$$
\left|x_{k}(w)-x_{k}(q)\right| \lesssim \varepsilon d_{\mathcal{S}}(q) \leqslant \varepsilon \sigma(Q)
$$

so $\left|x_{k}(w)-x_{k}(v)\right| \lesssim \varepsilon \sigma(Q)$.
We have $V_{Q}^{+} \cap L=[v, \infty)$ and $\Gamma^{+} \cap L=[w, \infty)$, so if $a \in L \cap\left(V_{Q}^{+} \triangle \Gamma^{+}\right)$, then either $a \in[v, w]$ or $a \in[w, v]$, depending on the order of $v$ and $w$. In both cases, $d\left(a, V_{Q}\right) \leqslant\left|x_{k}(w)-x_{k}(v)\right| \lesssim \varepsilon \sigma(Q)$ and $d(a, \Gamma) \lesssim \varepsilon \sigma(Q)$, so (150) holds. If $\varepsilon$ is sufficiently small, then Lemma 52 implies

$$
d_{N_{4}(Q)}\left(\Gamma^{+}, E\right) \leqslant d_{N_{\frac{9}{2}}(Q)}\left(\Gamma^{+}, V_{Q}^{+}\right)+d_{N_{\frac{1}{\varepsilon}}(Q)}\left(V_{Q}^{+}, E\right) \lesssim \varepsilon \sigma(Q) .
$$

If $\varepsilon$ is sufficiently small, then

$$
d_{N_{4}(Q)}\left(\Gamma^{+}, E\right) \leqslant \theta \sigma(Q) .
$$

This completes the proof of Proposition 80.

## 10. Historical comments on Sparsest Cut and directions for further research

Among the well-established deep and multifaceted connections between theoretical computer science and pure mathematics, the Sparsest Cut Problem stands out for its profound and often unexpected impact on a variety of areas. Previous research on this question came hand-in-hand with the development of remarkable mathematical and algorithmic ideas that spurred many further works of importance in their own right. Because the present work belongs to this tradition, here we will put it into context by elaborating further on the history of these investigations. We will also describe interesting directions for further research and open problems.

The first polynomial-time algorithm for Sparsest Cut with $O(\log n)$ approximation ratio was obtained in the important work [72], which studied the notable special case of Sparsest Cut with Uniform Demands; see Section 10.1 below. This work introduced a linear programming relaxation and developed influential techniques for its analysis, and it has led to a myriad of algorithmic applications. The seminal contributions [77], [7] obtained the upper bound
$\rho_{\mathrm{GL}}(n) \lesssim \log n$ in full generality by incorporating a classical embedding theorem of Bourgain [15], thus heralding the transformative use of metric embeddings in algorithm design. The matching lower bound on the integrality gap of this linear program was proven in [72], [77]. This showed for the first time that Bourgain's embedding theorem is asymptotically sharp and was the first (and very influential) demonstration of the power of expander graphs in the study of metric embeddings.

An $O(\sqrt{\log n})$ upper bound for the approximation ratio of the GoemansLinial algorithm in the special case of uniform demands was obtained in the important work [5]. This work relied on a clever use of the concentration of measure phenomenon and introduced influential techniques such as a "chaining argument" for metrics of negative type and the use of expander flows. [5] also had direct impact on results in pure mathematics, including combinatorics and metric geometry; see e.g., the "edge replacement theorem" and the estimates on the observable diameter of doubling metric measure spaces in [94]. The best-known general upper bound $\rho_{\mathrm{GL}}(n) \lesssim(\log n)^{\frac{1}{2}+o(1)}$ of [4] built on the (then very recent) development of two techniques: the chaining argument of [5] (through its more careful analysis in [67]) and the measured descent embedding method of [63] (through its phrasing as a gluing technique for Lipschitz maps in [67]). Another important input to [4] was a re-weighting argument of [20] that allowed for the construction of an appropriate "random zero set" from the argument of [5], [67]; see [90], [92] for more on this notion and its significance.

The impossibility result [59] that refuted the Goemans-Linial conjecture relied on a striking link to complexity theory through the Unique Games Conjecture (UGC), as well as an interesting use of discrete harmonic analysis (through [16]) in this context; see also [64] for an incorporation of a different tool from discrete harmonic analysis (namely, [54], following [58]) for the same purpose, as well as [21], [28] for computational hardness. The best impossibility result currently known [55] for Sparsest Cut with Uniform Demands relies on the development of new pseudorandom generators.

The idea of using the geometry of the Heisenberg group to bound $\rho_{\mathrm{GL}}(n)$ from below originated in [69], where the relevant metric of negative type was constructed through a complex-analytic argument, and initial (qualitative) impossibility results were presented through the use of Pansu's differentiation theorem [101] and the Radon-Nikodým Property from functional analysis; see, e.g., [13]. In [23], it was shown that the Heisenberg group indeed provides a proof that $\lim _{n \rightarrow \infty} \rho_{\mathrm{GL}}(n)=\infty$. This proof introduced a remarkable new notion of differentiation for mappings into $L_{1}(\mathbb{R})$, which led to the use of tools from geometric measure theory [38], [39] to study the problem. A different proof that $\mathbb{H}^{3}$ fails to admit a bi-Lipschitz embedding into $L_{1}(\mathbb{R})$ was found
in [24], where a classical notion of metric differentiation [60] was used in conjunction with the novel idea to consider monotonicity of sets in this context, combined with a sub-Riemannian-geometric argument that yielded a classification of monotone subsets of $\mathbb{H}^{3}$. The main result of [26] finds a quantitative lower estimate for the scale at which this differentiation argument can be applied, leading to a lower bound of $(\log n)^{\Omega(1)}$ on $\rho_{\mathrm{GL}}(n)$. This result relies on a mixture of the methods of [23] and [24] and overcomes obstacles that are not present in the original qualitative investigations. In particular, [26] introduced quantitative measures of nonmonotonicity that we use in the present work. The quantitative differentiation bound of [26] remains the best bound currently known, and it would be very interesting to discover the sharp behavior in this more subtle question.

The desire to avoid the (often difficult) need to obtain sharp bounds for quantitative differentiation motivated the investigations [8], [66]. In particular, [8] devised a method to prove sharp (up to lower order factors) nonembeddability statements for the Heisenberg group based on a cohomological argument and a quantitative ergodic theorem. For Hilbert space-valued mappings, [8] used a cohomological argument in combination with representation theory to prove the following quadratic inequality for every finitely supported function $\phi: \mathbb{H}_{\mathbb{Z}}^{2 k+1} \rightarrow L_{2}(\mathbb{R})$ :

$$
\begin{align*}
& \left(\sum_{t=1}^{\infty} \frac{1}{t^{2}} \sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left\|\phi\left(h Z^{t}\right)-\phi(h)\right\|_{L_{2}(\mathbb{R})}^{2}\right)^{\frac{1}{2}} \\
& \quad \lesssim\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{2}}\|\phi(h \sigma)-\phi(h)\|_{L_{2}(\mathbb{R})}^{2}\right)^{\frac{1}{2}} . \tag{151}
\end{align*}
$$

In [66] a different approach based on Littlewood-Paley theory was devised, leading to the following generalization of (151) that holds for every $p \in(1,2]$ and every finitely supported $\phi: \mathbb{H}_{Z}^{2 k+1} \rightarrow L_{p}(\mathbb{R})$ :

$$
\begin{align*}
& \left(\sum_{t=1}^{\infty} \frac{1}{t^{2}}\left(\sum_{h \in \mathbb{H}_{\mathbb{Z}}^{2 k+1}}\left\|\phi\left(h Z^{t}\right)-\phi(h)\right\|_{L_{p}(\mathbb{R})}^{p}\right)^{\frac{2}{p}}\right)^{\frac{1}{2}}  \tag{152}\\
& \quad \leqslant C(p)\left(\sum_{h \in \mathbb{H}_{Z}^{2 k+1}} \sum_{\sigma \in \mathfrak{S}_{2}}\|\phi(h \sigma)-\phi(h)\|_{L_{p}(\mathbb{R})}^{p}\right)^{\frac{1}{p}}
\end{align*}
$$

for some $C(p) \in(0, \infty)$. See [66] for a strengthening of (152) that holds for general uniformly convex targets (using the recently established [79] vector-valued Littlewood-Paley-Stein theory for the Poisson semigroup). These functional inequalities yield sharp nonembeddability estimates for balls in $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$, but
the method of (152) inherently yields a constant $C(p)$ in (152) that satisfies $\lim _{p \rightarrow 1} C(p)=\infty$. The estimate (15) that we prove here for $L_{1}(\mathbb{R})$-valued mappings is an endpoint estimate corresponding to (152), showing that $C(p)$ actually remains bounded as $p \rightarrow 1$. This answers positively a conjecture of [66] and is crucial for the results that we obtain here.

As explained in Section 5, our proof of (151) uses the $\mathbb{H}^{3}$-analogue of (151). It should be mentioned at this juncture that the proofs of (151) and (152) in [8], [66] were oblivious to the dimension of the underlying Heisenberg group. An unexpected aspect of the present work is that the underlying dimension does play a role at the endpoint $p=1$, with the analogue of (15) (or Theorem 3) for $\mathbb{H}_{\mathbb{Z}}^{3}$ being in fact incorrect; recall Remark 4.

As we recalled above, past progress on the Sparsest Cut Problem came hand-in-hand with substantial mathematical developments. The present work is a culmination of a long-term project that is rooted in mathematical phenomena that are interesting not just for their relevance to approximation algorithms but also for their connections to the broader mathematical world. In the ensuing subsections we shall describe some further results and questions related to this general direction.
10.1. The Sparsest Cut Problem with Uniform Demands. An especially important special case of the Sparsest Cut Problem is when the demand matrix $D$ is the matrix all of whose entries equal 1 , i.e., $D=1_{\{1, \ldots, n\} \times\{1, \ldots, n\}} \in M_{n}(\mathbb{R})$, and the capacity matrix $C$ lies in $M_{n}(\{0,1\})$, i.e., all its entries are either 0 or 1 . This is known as the Sparsest Cut Problem with Uniform Demands. In this case $C$ can also be described as the adjacency matrix of a graph $G$ whose vertex set is $\{1, \ldots, n\}$ and whose edge set consists of those unordered pairs $\{i, j\} \subseteq$ $\{1, \ldots, n\}$ for which $C_{i j}=1$. With this interpretation, given $A \subseteq\{1, \ldots, n\}$ the numerator in (6) equals twice the number of edges that are incident to $A$ in $G$. And, since $D=\mathbf{1}_{\{1, \ldots, n\} \times\{1, \ldots, n\}}$, the denominator in (6) is equal to $2|A|(n-|A|) \asymp n \min \{|A|,|\{1, \ldots, n\} \backslash A|\}$. So, the Sparsest Cut Problem with Uniform Demands asks for an algorithm that takes as input a finite graph and outputs a quantity that is bounded above and below by universal constant multiples of its conductance [107] divided by $n$. The Goemans-Linial integrality gap corresponding to this special case is

$$
\rho_{\mathrm{GL}}^{\mathrm{unif}}(n) \stackrel{\text { def }}{=} \sup _{\substack{C \in M_{n}(\{0,1\}) \\ C \text { symmetric }}} \frac{\operatorname{OPT}\left(C, \mathbf{1}_{\{1, \ldots, n\} \times\{1, \ldots, n\}}\right)}{\operatorname{SDP}\left(C, \mathbf{1}_{\{1, \ldots, n\} \times\{1, \ldots, n\})}\right.}
$$

The Goemans-Linial algorithm furnishes the best-known approximation ratio also in the case of uniform demands. We have $\rho_{\mathrm{GL}}^{\mathrm{unif}}(n) \lesssim \sqrt{\log n}$ by the important work [5], improving over the previous bound $\rho_{\mathrm{GL}}^{\mathrm{unif}}(n) \lesssim \log n$ of [72]. As explained in [25], the present approach based on (fixed dimensional) Heisenberg groups cannot yield a lower bound on $\rho_{\mathrm{GL}}^{\text {unif }}(n)$ that tends to $\infty$ as $n \rightarrow \infty$.

The currently best-known lower bound [55] is $\rho_{\mathrm{GL}}^{\mathrm{unif}}(n) \geqslant \exp (c \sqrt{\log \log n})$ for some universal constant $c \in(0, \infty)$, improving over the previous bound $\rho_{\mathrm{GL}}^{\text {unif }}(n) \gtrsim \log \log n$ of [34]. Determining the asymptotic behavior of $\rho_{\mathrm{GL}}^{\text {unif }}(n)$ remains an intriguing open problem.
10.2. A sharp upper bound on $\rho_{\mathrm{GL}}(n)$. The $o(1)$ term in the bound $\rho_{\mathrm{GL}}(n)$ $\lesssim(\log n)^{\frac{1}{2}+o(1)}$ that is stated in [4] is $\frac{\log \log \log n}{\log \log n}$, i.e., $\rho_{\mathrm{GL}}(n) \lesssim \sqrt{\log n} \log \log n$. This specific lower-order factor can be reduced with more technical work; see the discussion in Section 6 of [4]. To the best of our knowledge, the question of removing the $o(1)$ term altogether was not previously investigated, since the pressing issue was to determine the correct power of $\log n$. Now that Theorem 5 is established, it is worthwhile to revisit this matter to obtain the sharp upper bound $\rho_{\mathrm{GL}}(n) \lesssim \sqrt{\log n}$. We conjecture that this estimate is indeed valid, and we expect that its proof should follow from a more careful application of available methods. We postpone this investigation to future work.
10.3. Sharp stability of monotone sets. We already asked above what is the sharp dependence in the quantitative differentiation statement (stability of monotone sets) of [26]. This is a challenging open question whose need was circumvented in the present context by adopting the strategy of aiming for the functional inequality (15) rather than for a sharp structure theorem for sets of finite perimeter in $\mathbb{H}^{2 k+1}$ that shows that they are approximately a half-space on a macroscopically large scale. The advantage of our approach is that it relies on a corona decomposition in which the approximation is by an intrinsic Lipschitz graph rather than a half-space, i.e., we avoid the need to "zoom in" all the way to the scale at which the boundary resembles a half-space. Intrinsic Lipschitz graphs can be much less structured than half-spaces, but we control them using cancellations (in a lower dimension) that can be captured through a decomposition of each $\mathbb{H}^{3}$-slice into irreducible representations. Regardless of any application to embedding theory, sharpening the bound of [26] on the stability of monotone sets is an interesting question in geometric measure theory.
10.4. Hardness. Proving NP-hardness of approximation for the Sparsest Cut Problem remains elusive if one wishes to prove inapproximability within any constant factor. (Such UGC-based hardness of approximation is available due to $[59,21]$.) This is a longstanding open question that seems to be very challenging. We speculate that the hardness of approximation threshold for Sparsest Cut is $\Theta(\sqrt{\log n})$, but we do not see a path towards such a statement even under the UGC.
10.5. Further questions related to quantitative rectifiability and singular integrals. It would be interesting to explore the impact of the methods of the present article on other questions related to quantitative rectifiability and
singular integral operators. In future work we will investigate how an application of these ideas to $\mathbb{R}^{n}$ rather than the Heisenberg group yields a new perspective on classical results such as the Jones Travelling Salesman Theorem [52], and we believe that there is scope for additional applications of the use of quantitative nonmonotonicity to govern hierarchical partitioning schemes as well as the type of "induction on dimension" that is used here to pass from $L_{2}$ bounds to their $L_{1}$ counterparts (in a higher dimension).

Added in proof. The questions that are raised in Remark 4 are resolved in our forthcoming manuscript [97]. Specifically, it is shown in [97] that when $k=1$, the estimate (5) holds for $q=4$ and does not hold for any $q<4$. Various applications of these statements are described in [97].

## References

[1] L. Ambrosio, Some fine properties of sets of finite perimeter in Ahlfors regular metric measure spaces, Adv. Math. 159 no. 1 (2001), 51-67. MR 1823840. Zbl 1002.28004. https://doi.org/10.1006/aima.2000.1963.
[2] L. Ambrosio, N. Fusco, and D. Pallara, Functions of Bounded Variation and Free Discontinuity Problems, Oxford Math. Monogr., The Clarendon Press, Oxford University Press, New York, 2000. MR 1857292. Zbl 0957.49001.
[3] A. Andoni, A. Naor, and O. Neiman, Snowflake universality of Wasserstein spaces, 2015. arXiv 1509.08677.
[4] S. Arora, J. R. Lee, and A. Naor, Euclidean distortion and the sparsest cut, J. Amer. Math. Soc. 21 no. 1 (2008), 1-21. MR 2350049. Zbl 1132.68070. https://doi.org/10.1090/S0894-0347-07-00573-5.
[5] S. Arora, S. Rao, and U. Vazirani, Expander flows, geometric embeddings and graph partitioning, J. ACM 56 no. 2 (2009), Art. 5, 37. MR 2535878. Zbl 1325.68255. https://doi.org/10.1145/1502793.1502794.
[6] P. Assouad, Plongements lipschitziens dans $\mathbf{R}^{n}$, Bull. Soc. Math. France 111 no. 4 (1983), 429-448. MR 0763553. Zbl 0597.54015. https://doi.org/10.24033/ bsmf. 1997.
[7] Y. Aumann and Y. Rabani, An $O(\log k)$ approximate min-cut max-flow theorem and approximation algorithm, SIAM J. Comput. 27 no. 1 (1998), 291-301. MR 1614825. Zbl 0910.05038. https://doi.org/10.1137/S0097539794285983.
[8] T. Austin, A. Naor, and R. Tessera, Sharp quantitative nonembeddability of the Heisenberg group into superreflexive Banach spaces, Groups Geom. Dyn. 7 no. 3 (2013), 497-522. MR 3095705. Zbl 1284.46019. https://doi.org/10.4171/ GGD/193.
[9] K. Ball, Markov chains, Riesz transforms and Lipschitz maps, Geom. Funct. Anal. 2 no. 2 (1992), 137-172. MR 1159828. Zbl 0788.46050. https://doi.org/ 10.1007/BF01896971.
[10] K. Ball, The Ribe programme, in Séminaire Bourbaki. Vol. 2011/2012. Exposés 1043-1058, Astérisque 352, Soc. Math. France, Paris, 2013, pp. Exp. No. 1047, viii, 147-159. MR 3087345. Zbl 1303.46019.
[11] H. Bass, The degree of polynomial growth of finitely generated nilpotent groups, Proc. London Math. Soc. (3) 25 (1972), 603-614. MR 0379672. Zbl 0259. 20045. https://doi.org/10.1112/plms/s3-25.4.603.
[12] F. Baudier and G. Lancien, Embeddings of locally finite metric spaces into Banach spaces, Proc. Amer. Math. Soc. 136 no. 3 (2008), 1029-1033. MR 2361878. Zbl 1142.46037. https://doi.org/10.1090/ S0002-9939-07-09109-5.
[13] Y. Benyamini and J. Lindenstrauss, Geometric Nonlinear Functional Analysis. Vol. 1, Amer. Math. Soc. Colloq. Publ. 48, Amer. Math. Soc., Providence, RI, 2000. MR 1727673. Zbl 0946. 46002.
[14] S. Blachère, Word distance on the discrete Heisenberg group, Colloq. Math. 95 no. 1 (2003), 21-36. MR 1967551. Zbl 1010.05034. https://doi.org/10.4064/ cm95-1-2.
[15] J. Bourgain, On Lipschitz embedding of finite metric spaces in Hilbert space, Israel J. Math. 52 no. 1-2 (1985), 46-52. MR 0815600. Zbl 0657.46013. https: //doi.org/10.1007/BF02776078.
[16] J. Bourgain, On the distributions of the Fourier spectrum of Boolean functions, Israel J. Math. 131 (2002), 269-276. MR 1942312. Zbl 1021.43004. https://doi.org/10.1007/BF02785861.
[17] A. Brunel and L. Sucheston, Sur quelques conditions équivalentes à la superréflexivité dans les espaces de Banach, C. R. Acad. Sci. Paris Sér. A-B 275 (1972), A993-A994. MR 0318843. Zbl 0245. 46016.
[18] D. Burago, Y. Burago, and S. Ivanov, A Course in Metric Geometry, Grad. Studies in Math. 33, Amer. Math. Soc., Providence, RI, 2001. MR 1835418. Zbl 0981.51016. https://doi.org/10.1090/gsm/033.
[19] S. Chawla, Sparsest cut, in Encyclopedia of Algorithms, Springer-Verlag, New York, 2008, pp. 868-870. https://doi.org/10.1007/978-0-387-30162-4_388.
[20] S. Chawla, A. Gupta, and H. Räcke, Embeddings of negative-type metrics and an improved approximation to generalized sparsest cut, ACM Trans. Algorithms 4 no. 2 (2008), Art. 22, 18. MR 2419119. Zbl 1297.68234. https: //doi.org/10.1145/1361192.1361199.
[21] S. Chawla, R. Krauthgamer, R. Kumar, Y. Rabani, and D. SivakumAR, On the hardness of approximating multicut and sparsest-cut, Comput. Complexity 15 no. 2 (2006), 94-114. MR 2243123. Zbl 1132.68418. https: //doi.org/10.1007/s00037-006-0210-9.
[22] J. Cheeger and B. Kleiner, On the differentiability of Lipschitz maps from metric measure spaces to Banach spaces, in Inspired by S. S. Chern, Nankai Tracts Math. 11, World Sci. Publ., Hackensack, NJ, 2006, pp. 129-152. MR 2313333. Zbl 1139.58004. https://doi.org/10.1142/9789812772688_0006.
[23] J. Cheeger and B. Kleiner, Differentiating maps into $L^{1}$, and the geometry of BV functions, Ann. of Math. (2) 171 no. 2 (2010), 1347-1385. MR 2630066. Zbl 1194.22009. https://doi.org/10.4007/annals.2010.171.1347.
[24] J. Cheeger and B. Kleiner, Metric differentiation, monotonicity and maps to $L^{1}$, Invent. Math. 182 no. 2 (2010), 335-370. MR 2729270. Zbl 1214.46013. https://doi.org/10.1007/s00222-010-0264-9.
[25] J. Cheeger, B. Kleiner, and A. Naor, A $(\log n)^{\Omega(1)}$ integrality gap for the sparsest cut SDP, in 2009 50th Annual IEEE Symposium on Foundations of Computer Science—FOCS 2009, IEEE Computer Soc., Los Alamitos, CA, 2009, pp. 555-564. MR 2648435. Zbl 1291.90318. https://doi.org/10.1109/FOCS. 2009.47.
[26] J. Cheeger, B. Kleiner, and A. Naor, Compression bounds for Lipschitz maps from the Heisenberg group to $L_{1}$, Acta Math. 207 no. 2 (2011), 291-373. MR 2892612. Zbl 1247.46020. https://doi.org/10.1007/s11511-012-0071-9.
[27] M. Christ, A $T(b)$ theorem with remarks on analytic capacity and the Cauchy integral, Colloq. Math. 60/61 no. 2 (1990), 601-628. MR 1096400. Zbl 0758. 42009. https://doi.org/10.4064/cm-60-61-2-601-628.
[28] J. Chuzhoy and S. Khanna, Polynomial flow-cut gaps and hardness of directed cut problems [extended abstract], in STOC'07-Proceedings of the 39th Annual ACM Symposium on Theory of Computing, ACM, New York, 2007, pp. 179-188. MR 2402441. Zbl 1232.68063. https://doi.org/10.1145/1250790.1250817.
[29] J. Chuzhoy and S. Khanna, Polynomial flow-cut gaps and hardness of directed cut problems, J. ACM 56 no. 2 (2009), Art. 6, 28. MR 2535879. Zbl 1325. 68096. https://doi.org/10.1145/1502793.1502795.
[30] G. David and S. Semmes, Singular integrals and rectifiable sets in $\mathbf{R}^{n}$ : Beyond Lipschitz graphs, Astérisque 193, Soc. Math. France, Paris, 1991. MR 1113517. Zbl 0743. 49018.
[31] G. David, Opérateurs intégraux singuliers sur certaines courbes du plan complexe, Ann. Sci. École Norm. Sup. (4) 17 no. 1 (1984), 157-189. MR 0744071. Zbl 0537.42016 . https://doi.org/10.24033/asens.1469.
[32] G. David, Wavelets and Singular Integrals on Curves and Surfaces, Lecture Notes in Math. 1465, Springer-Verlag, Berlin, 1991. MR 1123480. Zbl 0764. 42019. https://doi.org/10.1007/BFb0091544.
[33] G. David and S. Semmes, Analysis of and on Uniformly Rectifiable Sets, Math. Surveys Monogr. 38, Amer. Math. Soc., Providence, RI, 1993. MR 1251061. Zbl 0832.42008. https://doi.org/10.1090/surv/038.
[34] N. R. Devanur, S. A. Khot, R. Saket, and N. K. Vishnoi, Integrality gaps for sparsest cut and minimum linear arrangement problems, in STOC ${ }^{\prime} 06$ : Proceedings of the 38th Annual ACM Symposium on Theory of Computing, ACM, New York, 2006, pp. 537-546. MR 2277179. Zbl 1301.05332. https: //doi.org/10.1145/1132516.1132594.
[35] M. M. Deza and M. Laurent, Geometry of Cuts and Metrics, Algorithms Combin. 15, Springer-Verlag, Berlin, 1997. MR 1460488. Zbl 0885.52001. https: //doi.org/10.1007/978-3-642-04295-9.
[36] J. Ding, J. R. Lee, and Y. Peres, Markov type and threshold embeddings, Geom. Funct. Anal. 23 no. 4 (2013), 1207-1229. MR 3077911. Zbl 1279.46013. https://doi.org/10.1007/s00039-013-0234-7.
[37] B. Franchi, R. Serapioni, and F. Serra Cassano, Meyers-Serrin type theorems and relaxation of variational integrals depending on vector fields, Houston J. Math. 22 no. 4 (1996), 859-890. MR 1437714. Zbl 0876.49014. Available at https://www.math.uh.edu/~hjm/restricted/ archive/v022n4/0859FRANCHI.pdf.
[38] B. Franchi, R. Serapioni, and F. Serra Cassano, Rectifiability and perimeter in the Heisenberg group, Math. Ann. 321 no. 3 (2001), 479-531. MR 1871966. Zbl 1057.49032. https://doi.org/10.1007/s002080100228.
[39] B. Franchi, R. Serapioni, and F. Serra Cassano, On the structure of finite perimeter sets in step 2 Carnot groups, J. Geom. Anal. 13 no. 3 (2003), 421-466. MR 1984849. Zbl 1064.49033. https://doi.org/10.1007/BF02922053.
[40] B. Franchi, R. Serapioni, and F. Serra Cassano, Intrinsic Lipschitz graphs in Heisenberg groups, J. Nonlinear Convex Anal. 7 no. 3 (2006), 423-441. MR 2287539. Zbl 1151. 58005.
[41] B. Franchi, R. Serapioni, and F. Serra Cassano, Differentiability of intrinsic Lipschitz functions within Heisenberg groups, J. Geom. Anal. 21 no. 4 (2011), 1044-1084. MR 2836591. Zbl 1234.22002. https://doi.org/10.1007/ s12220-010-9178-4.
[42] N. Garofalo and D.-M. Nhieu, Isoperimetric and Sobolev inequalities for Carnot-Carathéodory spaces and the existence of minimal surfaces, Comm. Pure Appl. Math. 49 no. 10 (1996), 1081-1144. MR 1404326. Zbl 0880.35032. https: //doi.org/10/dzvwvb.
[43] M. X. Goemans, Semidefinite programming in combinatorial optimization, Math. Programming 79 no. 1-3, Ser. B (1997), 143-161, Lectures on mathematical programming (ismp97) (Lausanne, 1997). MR 1464765. Zbl 0887.90139. https://doi.org/10.1007/BF02614315.
[44] M. Gromov, Groups of polynomial growth and expanding maps, Inst. Hautes Études Sci. Publ. Math. 53 (1981), 53-73. MR 0623534. Zbl 0474.20018. https: //doi.org/10.1007/BF02698687.
[45] M. Gromov, Carnot-Carathéodory spaces seen from within, in Sub-Riemannian Geometry, Progr. Math. 144, Birkhäuser, Basel, 1996, pp. 79-323. MR 1421823. Zbl 0864.53025.
[46] M. Grötschel, L. Lovász, and A. Schrijver, Geometric Algorithms and Combinatorial Optimization, second ed., Algorithms Combin. 2, SpringerVerlag, Berlin, 1993. MR 1261419. Zbl 0837.05001. https://doi.org/10.1007/ 978-3-642-78240-4.
[47] A. Gupta, R. Krauthgamer, and J. R. Lee, Bounded geometries, fractals, and low-distortion embeddings, in FOCS '03 Proceedings of the 44th Symposium on Foundations of Computer Science (October 11-14, 2003, Cambridge, MA, USA), IEEE Computer Society, Washington, DC, 2003, pp. 534-543. Available at https://dl.acm.org/citation.cfm?id=946308.
[48] P. HajŁasz and P. Koskela, Sobolev met Poincaré, Mem. Amer. Math. Soc. 145 no. 688 (2000), x+101. MR 1683160. Zbl 0954.46022. https://doi.org/10. 1090/memo/0688.
[49] S. Heinrich, Ultraproducts in Banach space theory, J. Reine Angew. Math. 313 (1980), 72-104. MR 0552464. Zbl 0412.46017. https://doi.org/10.1515/ crll.1980.313.72.
[50] A. Jaffe, J. R. Lee, and M. Moharrami, On the optimality of gluing over scales, Discrete Comput. Geom. 46 no. 2 (2011), 270-282. MR 2812509. Zbl 1219.68159. https://doi.org/10.1007/s00454-011-9359-3.
[51] P. W. Jones, Square functions, Cauchy integrals, analytic capacity, and harmonic measure, in Harmonic Analysis and Partial Differential Equations (El Escorial, 1987), Lecture Notes in Math. 1384, Springer, Berlin, 1989, pp. 24-68. MR 1013815. Zbl 0675.30029. https://doi.org/10.1007/BFb0086793.
[52] P. W. Jones, Rectifiable sets and the traveling salesman problem, Invent. Math. 102 no. 1 (1990), 1-15. MR 1069238. Zbl 0731.30018. https://doi.org/10.1007/ BF01233418.
[53] M. I. Kadec, Linear dimension of the spaces $L_{p}$ and $l_{q}$, Uspehi Mat. Nauk 13 no. 6 (84) (1958), 95-98. MR 0101486. Zbl 0087.10802. Available at http: //mi.mathnet.ru/eng/umn/v13/i6/p95.
[54] J. Kahn, G. Kalai, and N. Linial, The influence of variables on Boolean functions (extended abstract), in 29th Annual Symposium on Foundations of Computer Science (White Plains, New York, USA, 24-26 October 1988), IEEE Computer Society, 1988, pp. 68-80. https://doi.org/10.1109/SFCS.1988.21923.
[55] D. Kane and R. Meka, A PRG for Lipschitz functions of polynomials with applications to sparsest cut, in STOC'13—Proceedings of the 2013 ACM Symposium on Theory of Computing, ACM, New York, 2013, pp. 1-10. MR 3210761. Zbl 1293.65007. https://doi.org/10.1145/2488608.2488610.
[56] S. Khot, On the power of unique 2-prover 1-round games, in Proceedings of the Thirty-Fourth Annual ACM Symposium on Theory of Computing, ACM, New York, 2002, pp. 767-775. MR 2121525. Zbl 1192.68367. https://doi.org/ 10.1145/509907.510017.
[57] S. Kнот, Inapproximability of NP-complete problems, discrete Fourier analysis, and geometry, in Proceedings of the International Congress of Mathematicians 2010 (ICM 2010). Volume IV, Hindustan Book Agency, New Delhi, 2012, pp. 2676-2697. MR 2827989. Zbl 1252.68143. https://doi.org/10.1142/ 9789814324359_0163.
[58] S. Khot and A. Naor, Nonembeddability theorems via Fourier analysis, Math. Ann. 334 no. 4 (2006), 821-852. MR 2209259. Zbl 1102.46051. https://doi.org/ 10.1007/s00208-005-0745-0.
[59] S. A. Khot and N. K. Vishnoi, The unique games conjecture, integrability gap for cut problems and embeddability of negative-type metrics into $\ell_{1}, J . A C M$ 62 no. 1 (2015), Art. 8, 39. MR 3323774. Zbl 1321.68316. https://doi.org/10. 1145/2629614.
[60] B. Kirchheim, Rectifiable metric spaces: local structure and regularity of the Hausdorff measure, Proc. Amer. Math. Soc. 121 no. 1 (1994), 113-123. MR 1189747. Zbl 0806.28004. https://doi.org/10.2307/2160371.
[61] P. Klein, A. Agrawal, R. Ravi, and S. Rao, Approximation through multicommodity flow, in 31 st Annual Symposium on Foundations of Computer Science, Vol. I, II (St. Louis, MO, 1990), IEEE Comput. Soc. Press, Los Alamitos, CA, 1990, pp. 726-737. MR 1150728. https://doi.org/10.1109/FSCS. 1990. 89595.
[62] B. Kleiner, A new proof of Gromov's theorem on groups of polynomial growth, J. Amer. Math. Soc. 23 no. 3 (2010), 815-829. MR 2629989. Zbl 1246. 20038. https://doi.org/10.1090/S0894-0347-09-00658-4.
[63] R. Krauthgamer, J. R. Lee, M. Mendel, and A. Naor, Measured descent: a new embedding method for finite metrics, Geom. Funct. Anal. 15 no. 4 (2005), 839-858. MR 2221152. Zbl 1108.46010. https://doi.org/10.1007/ s00039-005-0527-6.
[64] R. Krauthgamer and Y. Rabani, Improved lower bounds for embeddings into $L_{1}$, SIAM J. Comput. 38 no. 6 (2009), 2487-2498. MR 2506299. Zbl 1191. 68869. https://doi.org/10.1137/060660126.
[65] V. Lafforgue and A. Naor, A doubling subset of $L_{p}$ for $p>2$ that is inherently infinite dimensional, Geom. Dedicata 172 (2014), 387-398. MR 3253787. Zbl 1305.30029. https://doi.org/10.1007/s10711-013-9924-4.
[66] V. Lafforgue and A. Naor, Vertical versus horizontal Poincaré inequalities on the Heisenberg group, Israel J. Math. 203 no. 1 (2014), 309-339. MR 3273443. Zbl 1312.46032. https://doi.org/10.1007/s11856-014-1088-x.
[67] J. R. Lee, On distance scales, embeddings, and efficient relaxations of the cut cone, in Proceedings of the Sixteenth Annual ACM-SIAM Symposium on Discrete Algorithms, ACM, New York, 2005, pp. 92-101. MR 2298254. Zbl 1297.68244.
[68] J. R. Lee, M. Mendel, and A. Naor, Metric structures in $L_{1}$ : dimension, snowflakes, and average distortion, European J. Combin. 26 no. 8 (2005), 11801190. MR 2163751. Zbl 1106.68086. https://doi.org/10.1016/j.ejc.2004.07.002.
[69] J. R. Lee and A. Naor, $L_{p}$ metrics on the Heisenberg group and the GoemansLinial conjecture, in Proceedings of 47th Annual IEEE Symposium on Foundations of Computer Science (FOCS 2006), IEEE Computer Soc., Los Alamitos, CA, 2006, pp. 99-108. https://doi.org/10.1109/FOCS.2006.47.
[70] J. R. Lee, A. Naor, and Y. Peres, Trees and Markov convexity, Geom. Funct. Anal. 18 no. 5 (2009), 1609-1659. MR 2481738. Zbl 1171.05318. https: //doi.org/10.1007/s00039-008-0689-0.
[71] J. R. Lee and A. Sidiropoulos, Near-optimal distortion bounds for embedding doubling spaces into $L_{1}$ [extended abstract], in STOC'11—Proceedings of the 43 rd ACM Symposium on Theory of Computing, ACM, New York, 2011, pp. 765-772. MR 2932027. Zbl 1288.90079. https://doi.org/10.1145/1993636. 1993737.
[72] T. Leighton and S. Rao, Multicommodity max-flow min-cut theorems and their use in designing approximation algorithms, J. ACM 46 no. 6 (1999), 787832. MR 1753034. Zbl 1065.68666. https://doi.org/10.1145/331524.331526.
[73] S. Li, Coarse differentiation and quantitative nonembeddability for Carnot groups, J. Funct. Anal. 266 no. 7 (2014), 4616-4704. MR 3170215. Zbl 1311. 46021. https://doi.org/10.1016/j.jfa.2014.01.026.
[74] S. Li, Markov convexity and nonembeddability of the Heisenberg group, Ann. Inst. Fourier (Grenoble) 66 no. 4 (2016), 1615-1651. MR 3494180. Zbl 06680283. https://doi.org/10.5802/aif.3045.
[75] N. Linial, Finite metric-spaces-combinatorics, geometry and algorithms, in Proceedings of the International Congress of Mathematicians, Vol. III (Beijing, 2002), Higher Ed. Press, Beijing, 2002, pp. 573-586. MR 1957562. Zbl 0997. 05019.
[76] N. Linial, Squared $\ell_{2}$ metrics into $\ell_{1}$, in Open Problems on Embeddings of Finite Metric Spaces, 2002, edited by J. Matoušek, p. 5. Available at http: //students.mimuw.edu.pl/~mk249015/www10/metrop.pdf.
[77] N. Linial, E. London, and Y. Rabinovich, The geometry of graphs and some of its algorithmic applications, Combinatorica 15 no. 2 (1995), 215-245. MR 1337355. Zbl 0827.05021. https://doi.org/10.1007/BF01200757.
[78] K. Makarychev, Y. Makarychev, and A. Vijayaraghavan, Bilu-Linial stable instances of max cut and minimum multiway cut, in Proceedings of the Twenty-Fifth Annual ACM-SIAM Symposium on Discrete Algorithms, ACM, New York, 2014, pp. 890-906. MR 3376427. Zbl 1315.05114. https://doi.org/ 10.1137/1.9781611973402.67.
[79] T. Martínez, J. L. Torrea, and Q. Xu, Vector-valued Littlewood-PaleyStein theory for semigroups, Adv. Math. 203 no. 2 (2006), 430-475. MR 2227728. Zbl 1111.46008. https://doi.org/10.1016/j.aim.2005.04.010.
[80] J. Matoušek, Lectures on Discrete Geometry, Grad. Texts in Math. 212, Springer-Verlag, New York, 2002. MR 1899299. Zbl 0999.52006. https://doi. org/10.1007/978-1-4613-0039-7.
[81] B. Maurey, Type, cotype and K-convexity, in Handbook of the Geometry of Banach Sspaces, Vol. 2, North-Holland, Amsterdam, 2003, pp. 12991332. MR 1999197. Zbl 1074.46006. https://doi.org/10.1016/S1874-5849(03) 80037-2.
[82] B. Maurey and G. Pisier, Séries de variables aléatoires vectorielles indépendantes et propriétés géométriques des espaces de Banach, Studia Math. 58 no. 1 (1976), 45-90. MR 0443015. Zbl 0344.47014. https://doi.org/10.4064/ sm-58-1-45-90.
[83] E. J. McShane, Extension of range of functions, Bull. Amer. Math. Soc. 40 no. 12 (1934), 837-842. MR 1562984. Zbl 60.0973.02. https://doi.org/10.1090/ S0002-9904-1934-05978-0.
[84] M. Mendel, Metric dichotomies, in Limits of Graphs in Group Theory and Computer Science, EPFL Press, Lausanne, 2009, pp. 59-76. MR 2562140. Zbl 1293.05386.
[85] M. Mendel and A. Naor, A note on dichotomies for metric transforms, 2011. arXiv 1102.1800.
[86] M. Mendel and A. Naor, Markov convexity and local rigidity of distorted metrics, J. Eur. Math. Soc. (JEMS) 15 no. 1 (2013), 287-337. MR 2998836. Zbl 1266.46016. https://doi.org/10.4171/JEMS/362.
[87] V. D. Milman and H. Wolfson, Minkowski spaces with extremal distance from the Euclidean space, Israel J. Math. 29 no. 2-3 (1978), 113-131. MR 0467255. Zbl 0374.46013. https://doi.org/10.1007/BF02762002.
[88] F. Montefalcone, Some relations among volume, intrinsic perimeter and onedimensional restrictions of BV functions in Carnot groups, Ann. Sc. Norm. Super. Pisa Cl. Sci. (5) 4 no. 1 (2005), 79-128. MR 2165404. Zbl 1150.49022. https://doi.org/10.2422/2036-2145.2005.1.04.
[89] R. Montgomery, A Tour of Subriemannian Geometries, Their Geodesics and Applications, Math. Surveys Monogr. 91, Amer. Math. Soc., Providence, RI, 2002. MR 1867362. Zbl 1044.53022.
[90] A. NaOR, $L_{1}$ embeddings of the Heisenberg group and fast estimation of graph isoperimetry, in Proceedings of the International Congress of Mathematicians. Volume III, Hindustan Book Agency, New Delhi, 2010, pp. 1549-1575. MR 2827855. Zbl 1232.46021.
[91] A. NaOR, An introduction to the Ribe program, Jpn. J. Math. 7 no. 2 (2012), 167-233. MR 2995229. Zbl 1261.46013. https://doi.org/10.1007/ s11537-012-1222-7.
[92] A. NaOR, Comparison of metric spectral gaps, Anal. Geom. Metr. Spaces 2 (2014), 1-52. MR 3208067. Zbl 1316.46023. https://doi.org/10.2478/ agms-2014-0001.
[93] A. Naor, Y. Peres, O. Schramm, and S. Sheffield, Markov chains in smooth Banach spaces and Gromov-hyperbolic metric spaces, Duke Math. J. 134 no. 1 (2006), 165-197. MR 2239346. Zbl 1108.46012. https://doi.org/10. 1215/S0012-7094-06-13415-4.
[94] A. Naor, Y. Rabani, and A. Sinclair, Quasisymmetric embeddings, the observable diameter, and expansion properties of graphs, J. Funct. Anal. 227 no. 2 (2005), 273-303. MR 2168076. Zbl 1104.68087. https://doi.org/10.1016/ j.jfa.2005.04.003.
[95] A. Naor and L. Silberman, Poincaré inequalities, embeddings, and wild groups, Compos. Math. 147 no. 5 (2011), 1546-1572. MR 2834732. Zbl 1267. 20057. https://doi.org/10.1112/S0010437X11005343.
[96] A. Naor and R. Young, The integrality gap of the Goemans-Linial SDP relaxation for sparsest cut is at least a constant multiple of $\sqrt{\log n}$, in $S T O C^{\prime} 17-$ Proceedings of the 49th Annual ACM SIGACT Symposium on Theory of Computing, ACM, New York, 2017, pp. 564-575. MR 3678211. Zbl 1370.68235. https://doi.org/10.1145/3055399.3055413.
[97] A. Naor and R. Young, Foliated corona decompositions and the $L_{1}$ distortion of balls in the 3-dimensional Heisenberg group, 2018.
[98] M. I. Ostrovskir, Embeddability of locally finite metric spaces into Ba nach spaces is finitely determined, Proc. Amer. Math. Soc. 140 no. 8 (2012), 2721-2730. MR 2910760. Zbl 1276.46013. https://doi.org/10.1090/ S0002-9939-2011-11272-3.
[99] M. I. Ostrovskiı, Metric Embeddings, De Gruyter Stud. Math. 49, De Gruyter, Berlin, 2013, Bilipschitz and coarse embeddings into Banach spaces. MR 3114782. Zbl 1279.46001. https://doi.org/10.1515/9783110264012.
[100] P. Pansu, Une inégalité isopérimétrique sur le groupe de Heisenberg, C. $R$. Acad. Sci. Paris Sér. I Math. 295 no. 2 (1982), 127-130. MR 0676380. Zbl 0502. 53039.
[101] P. Pansu, Métriques de Carnot-Carathéodory et quasiisométries des espaces symétriques de rang un, Ann. of Math. (2) 129 no. 1 (1989), 1-60. MR 0979599. Zbl 0678.53042. https://doi.org/10.2307/1971484.
[102] G. Pisier, Martingales with values in uniformly convex spaces, Israel J. Math. 20 no. 3-4 (1975), 326-350. MR 0394135. Zbl 0344.46030. https://doi.org/10. 1007/BF02760337.
[103] G. Pisier and Q. Xu, Non-commutative $L^{p}$-Spaces, in Handbook of the Geometry of Banach spaces, Vol. 2, North-Holland, Amsterdam, 2003, pp. 14591517. MR 1999201. Zbl 1046.46048. https://doi.org/10.1016/S1874-5849(03) 80041-4.
[104] S. RAO, Small distortion and volume preserving embeddings for planar and Euclidean metrics, in Proceedings of the Fifteenth Annual Symposium on Computational Geometry (Miami Beach, FL, 1999), ACM, New York, 1999, pp. 300-306. MR 1802217. https://doi.org/10.1145/304893.304983.
[105] F. Shahrokhi and D. W. Matula, The maximum concurrent flow problem, J. Assoc. Comput. Mach. 37 no. 2 (1990), 318-334. MR 1072261. Zbl 0696.68071. https://doi.org/10.1145/77600.77620.
[106] D. B. Shmoys, Cut problems and their application to divide-and-conquer, in Approximation Algorithms for NP-hard Problems, PWS Publishing Co., Boston, MA, 1997, Chapter 5, D.S. Hochbaum, ed., pp. 192-235.
[107] A. Sinclair and M. Jerrum, Approximate counting, uniform generation and rapidly mixing Markov chains, Inform. and Comput. 82 no. 1 (1989), 93-133. MR 1003059. Zbl 0668.05060. https://doi.org/10.1016/0890-5401(89)90067-9.
[108] R. Tessera, Quantitative property A, Poincaré inequalities, $L^{p}$-compression and $L^{p}$-distortion for metric measure spaces, Geom. Dedicata 136 (2008), 203220. MR 2443353. Zbl 1162.46043. https://doi.org/10.1007/s10711-008-9286-5.
[109] L. Trevisan, On Khot's unique games conjecture, Bull. Amer. Math. Soc. (N.S.) 49 no. 1 (2012), 91-111. MR 2869009. Zbl 1280.68102. https://doi.org/ 10.1090/S0273-0979-2011-01361-1.
[110] J. H. Wells and L. R. Williams, Embeddings and Extensions in Analysis, 84, Springer-Verlag, New York, 1975, Ergeb. Math. Grenzgeb. MR 0461107. Zbl 0324.46034. https://doi.org/10.1007/978-3-642-66037-5.
[111] P. Wojtaszczyk, Banach Spaces for Analysts, Cambridge Stud. Adv. Math. 25, Cambridge University Press, Cambridge, 1991. MR 1144277. Zbl 0724.46012. https://doi.org/10.1017/CBO9780511608735.
[112] R. Young, Quantitative nonorientability of embedded cycles, Duke Math. J. 167 no. 1 (2018), 41-108. MR 3743699. Zbl 06847242. https://doi.org/10.1215/ 00127094-2017-0035.
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[^1]:    ${ }^{1}$ One could also consider classes of Banach spaces that are defined the same way but with $\mathbb{H}_{\mathbb{Z}}^{3}$ replaced with $\mathbb{H}_{\mathbb{Z}}^{2 k+1}$ for each $k \in \mathbb{N}$. The currently available evidence suggests that the resulting classes of Banach spaces do not actually depend on $k$, but this has not been proven. For simplicity and concreteness, we restrict the discussion here to $k=1$, but all of the questions could be studied for general $k$ as well.

