SPECIAL CUBULATION OF STRICT HYPERBOLIZATION

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Abstract. We prove that the Gromov hyperbolic groups obtained by the
strict hyperbolization procedure of Charney and Davis are virtually compact
special, hence linear and residually finite, if the initial complex satisfies some
minor conditions. Our strategy consists in constructing an action of a hyper-
bolized group on a certain dual CAT(0) cubical complex. As a result, all the
common applications of strict hyperbolization are shown to provide manifolds
with virtually compact special fundamental group. In particular, we obtain ex-
amples of closed negatively curved Riemannian manifolds whose fundamental
groups are linear and virtually algebraically fiber.

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1. Introduction

Closed aspherical manifolds occupy a central place in manifold topology. For this class of manifolds, the Borel Conjecture predicts that two such manifolds are homeomorphic if and only if they have isomorphic fundamental groups – in other words, that the topology is entirely encoded in the fundamental group. A challenging problem is the question of examples. The fundamental group will always satisfy Poincaré Duality over \( \mathbb{Z} \) (i.e. they are \( PD_n \) groups), and the Wall Conjecture predicts that conversely, any \( PD_n \) group is the fundamental group of an aspherical manifold. Classically, there were two sources of examples of aspherical manifolds: they either arose from Lie theory, as quotients of contractible Lie groups by discrete subgroups, or from differential geometry, as non-positively curved manifolds.

In the late 1970’s, Gromov introduced two metric versions of non-positive curvature, \( CAT(0) \) spaces and Gromov hyperbolic spaces. Simply connected, complete, locally \( CAT(0) \) spaces are automatically contractible. In dimensions \( \geq 4 \), manifolds that support locally \( CAT(0) \) metrics form a new source of aspherical manifolds. Moreover, it is easy to produce such manifolds, through a process known as hyperbolization. This was originally outlined by Gromov in [Gro87], and subsequently developed by Davis and Januszkiewicz in [DJ91]. Hyperbolization is a functorial procedure, which inputs a simplicial complex, and outputs a locally \( CAT(0) \) space. In a later refinement, Charney and Davis in [CD95] developed a strict hyperbolization procedure, where the output is locally \( CAT(-1) \), i.e. admits a metric of negative curvature (as opposed to just non-positive curvature). The hyperbolization procedures have been used to produce examples of aspherical manifolds with various unexpected properties. In this work we show that one can construct hyperbolizations that have some additional algebraic regularity.

**Theorem 1.1.** Given a dimension \( n > 0 \), there exists a strict hyperbolization procedure \( \mathcal{H} \) with the following property. Let \( K \) be any \( n \)-dimensional simplicial complex, which is compact, homogeneous, and without boundary. Then the resulting hyperbolized space \( \mathcal{H}(K) \) has fundamental group \( G = \pi_1(\mathcal{H}(K)) \) which acts cocompactly on a \( CAT(0) \) cubical complex by cubical isometries.

Most of the paper is concerned with the proof of this result (see §3 and §4). The cubical complex in this statement is \( n \)-dimensional, but it is not locally compact and the action is not proper. Nevertheless, since the hyperbolization procedure is strict (i.e. \( \mathcal{H}(K) \) is locally \( CAT(-1) \)), the fundamental group \( G \) is Gromov hyperbolic. Therefore the work of Agol, Haglund–Wise, and Groves–Manning about special cube complexes (see [Ago13; HW08; GM18]) can be used to extract information about \( G \). A cubical complex is special if it admits a local isometry to the Salvetti complex of a right-angled Artin group (RAAG) (see [HW08]). A group is virtually compact special if it has a finite index subgroup which is the fundamental group of a compact special cubical complex.
Theorem 1.2. Under the hypotheses of Theorem 1.1, the fundamental group $G = \pi_1(\mathcal{H}(K))$ is Gromov hyperbolic and virtually compact special. In particular, $G$ enjoys the following properties.

1. $G$ virtually embeds in a right–angled Artin group (RAAG) (see [HW08]).
2. $G$ is linear over $\mathbb{Z}$ (see [DJ00; HW08]), hence is residually finite.
3. $G$ has separable quasiconvex subgroups (see [HW08]).
4. $G$ is virtually residually finite rationally solvable (RFRS) (see [Ago08]).
5. $G$ has the Haagerup property, hence does not have property (T) (see [NR97; CMV04]).
6. $G$ satisfies the strong Atiyah conjecture (see [Sch14]).
7. $G$ is virtually bi-orderable (see [DK92]).
8. $G$ virtually embeds in the mapping class group of a closed surface, in a braid group, and in the group of diffeomorphisms of $\mathbb{R}$ (see [Kob12; KK15; BKK16]).
9. $G$ admits a proper affine action on $\mathbb{R}^n$ for some $n \geq 1$ (see [DGK20]).

This is achieved in §5 via a study of the cube stabilizers for the action from Theorem 1.1 and using a criterion for improper actions from [GM18]. The special cubical complex in Theorem 1.2 comes from a geometric action on a CAT(0) cubical complex different from the one in Theorem 1.1; its dimension is in general larger than $n$ and not easy to bound. We note that the fact that hyperbolized groups do not have property (T) was already observed by Belegradek in [Bel07] without using cubical methods.

The use of strict hyperbolization (as opposed to non–strict hyperbolization procedures) is crucial here. Indeed, there are closed aspherical manifolds whose fundamental group is not Gromov hyperbolic and not residually finite (see [Mes90; Bel06]). A well-known question by Gromov asks whether all Gromov hyperbolic groups are residually finite. Theorem 1.2 suggests that the strict hyperbolization procedure introduced by Charney and Davis in [CD95] is unlikely to provide counterexamples to this question.

1.1. The main arguments. The hyperbolization procedure in Theorems 1.1 and 1.2 is the composition of two hyperbolization procedures. The first one is Gromov’s cylinder construction, which turns the simplicial complex $K$ into a non–positively curved cubical complex $\mathcal{S}(K)$ (see §2.4, and [Gro87, §3.4.A]). The second one is the strict hyperbolization procedure of Charney and Davis, which turns a non–positively curved cubical complex $X$ into a locally $\text{CAT}(-1)$ piecewise hyperbolic polyhedron $X_\Gamma$ (see §3, and [CD95]). Here $\Gamma$ is a certain uniform arithmetic lattice of simple type in $\text{SO}_0(n, 1) = \text{Isom}^+(\mathbb{H}^n)$, which needs to be chosen to define the strict hyperbolization procedure. The hyperbolization procedure in Theorem 1.1 is then given by $\mathcal{H}(K) = (\mathcal{S}(K))_\Gamma$, i.e. by the composition

$$K \mapsto X = \mathcal{S}(K) \mapsto X_\Gamma = (\mathcal{S}(K))_\Gamma,$$

and in this paper we are mostly concerned with the study of the second part, i.e. the strict hyperbolization of a cubical complex. Sections §3 and §4 lead to the following argument.

Proof of Theorem 1.1. Let $K$ be an $n$-dimensional simplicial complex, which is compact, homogeneous, and without boundary. Then the cubical complex $X = \mathcal{S}(K)$ is an $n$-dimensional cubical complex, which is compact, homogeneous, and
without boundary. Moreover, up to a barycentric subdivision of \( K \), the cubical complex \( X = \mathcal{G}(K) \) can be assumed to be foldable, i.e. to admit a combinatorial map \( f : X \to \square^n \) to the standard cube which is injective on each cube (see Proposition 2.7).

Foldability provides a collection of subspaces of \( X \) that we call mirrors. A mirror is defined as a connected component of the full preimage in \( X \) of a codimension-1 face of the standard cube \( \square^n \) under the folding \( f : X \to \square^n \) (see §3.4). Mirrors of \( X \) give rise to nice locally convex codimension-1 subspaces of the hyperbolized complex \( X_\Gamma \), which we still call mirrors. Lifting the collection of mirrors of \( X_\Gamma \) to the universal cover \( \tilde{X}_\Gamma \) of \( X_\Gamma \) provides a stratification of \( \tilde{X}_\Gamma \): a point is in the \( k \)-stratum if it is contained in \( n - k \) mirrors (where \( n = \dim K = \dim X = \dim X_\Gamma \)).

We construct a dual cubical complex \( \mathcal{C}(\tilde{X}_\Gamma) \) in which vertices are given by cells in this stratification, and edges correspond to codimension-1 inclusion of cells (see §4). The complex \( \mathcal{C}(\tilde{X}_\Gamma) \) comes with a natural height function on its vertices, recording the dimension of the corresponding cell. In particular, the link of each vertex splits into an ascending sublink and a descending sublink. The former is flag because the cubical complex \( X = \mathcal{G}(K) \) is non–positively curved, and the latter is flag because of a Helly property satisfied by collections of pairwise orthogonal hyperplanes in \( \mathbb{H}^n \) (see Lemma 4.9). It follows that links of vertices in \( \mathcal{C}(\tilde{X}_\Gamma) \) are flag (see Proposition 4.10), hence \( \mathcal{C}(\tilde{X}_\Gamma) \) is a non–positively curved cubical complex. Moreover, the separation properties of the collection of mirrors (see §3.7) imply that \( \mathcal{C}(\tilde{X}_\Gamma) \) is simply–connected, hence \( \text{CAT}(0) \) (see Theorem 4.29).

Finally, note that the action of \( G = \pi_1(X_\Gamma) = \pi_1(\mathcal{H}(K)) \) on \( \tilde{X}_\Gamma \) by deck transformations induces an action of \( G \) on \( \mathcal{C}(\tilde{X}_\Gamma) \), as desired (see Lemma 4.30). \( \square \)

The action from Theorem 1.1 is further studied in §5.

**Proof of Theorem 1.2.** Under the hypothesis and in the setting of Theorem 1.1, the Gromov hyperbolic group \( G = \pi_1(X_\Gamma) = \pi_1(\mathcal{H}(K)) \) acts on the dual \( \text{CAT}(0) \) cubical complex \( \mathcal{C}(\tilde{X}_\Gamma) \) cocompactly and by cubical isometries.

The action is not proper, but the cube stabilizers can be identified with suitable cell stabilizers for the action of \( G \) by deck transformations on the universal cover \( \tilde{X}_\Gamma \) of \( \mathcal{H}(K) \) (see §5.1). These stabilizers are quasiconvex subgroups both of \( G \) and of \( \Gamma \) (see §5.2). Arithmetic lattices like \( \Gamma \) are known to be virtually compact special by [HW12]. In particular, we obtain that cell stabilizers for the action of \( G \) on \( \mathcal{C}(\tilde{X}_\Gamma) \) are virtually compact special. It then follows from [GM18, Theorem D] that \( G \) itself is virtually compact special (see Theorem 5.15). \( \square \)

For the sake of clarity: the cubical complex that witnesses the specialness of \( G \) is not the cubical complex from Theorem 1.1. It is obtained via the construction in [GM18], and its dimension is in general higher than \( n = \dim K \). One of the benefits of working with the dual \( \text{CAT}(0) \) cubical complex \( \mathcal{C}(\tilde{X}_\Gamma) \) (as opposed to other available \( \text{CAT}(0) \) cubical complexes, such as \( \tilde{X} \)) is that the stabilizers for the action of \( G \) on \( \mathcal{C}(\tilde{X}_\Gamma) \) can be related to the stabilizers for the action on \( \tilde{X}_\Gamma \), which are more geometric in nature and easier to understand.

### 1.2. Classical applications of hyperbolization procedures

The interest in hyperbolization procedures is that they can be used to construct closed aspherical manifolds with various interesting properties. As a result of our Theorem 1.2,
many applications of the strict hyperbolization procedure introduced by Charney and Davis in [CD95] can be obtained with additional algebraic features (e.g. the properties (1)-(9) listed in Theorem 1.2). We now collect some of these applications.

1.2.1. Riemannian hyperbolization. The strict hyperbolization procedure introduced by Charney and Davis in [CD95] outputs a space with a metric which is locally CAT(−1) and piecewise hyperbolic: the space is obtained by gluing together copies of the hyperbolizing cube □nΓ. When the cell complex X used in the hyperbolization procedure is homeomorphic to a smooth manifold, the hyperbolized complex XΓ is homeomorphic to a manifold too, but the locally CAT(−1) metric can a priori have singularities where the boundaries of different copies of the hyperbolizing cube □nΓ meet. It was recently shown by Ontaneda in [Ont20] that the construction can be tweaked in such a way that the manifold XΓ supports a smooth Riemannian metric with strictly negative sectional curvatures (possibly with respect to a different smooth structure).

This was used in [Ont20, Corollary 5] to construct examples in any dimension n ≥ 4 of closed Riemannian n-manifolds of pinched negative curvature which are “new” in the sense that they are not homeomorphic to any of the previously known examples of Riemannian manifold of negative curvature, such as closed real hyperbolic manifolds (or more generally locally symmetric spaces of rank 1), or the Gromov-Thurston branched covers in [GT87], or the examples of Mostow-Siu in [MS80] or Deraux in [Der05]. These manifolds are also distinct from the recent examples constructed by Stover–Toledo in [ST21b; ST21a], as the latter are Kähler, while the result of strict hyperbolization cannot be Kähler by [Bel07, Theorem 1.8].

Our construction does not require the smoothness provided by Ontaneda’s work, but it is compatible with it, so we get the following.

Corollary 1.3. For any ε > 0 and n ≥ 4 there are closed Riemannian n-manifolds with the following properties:

- they have sectional curvatures in the interval [−1 − ε, −1];
- they are not homeomorphic to a locally symmetric space of rank 1, or one of the manifolds constructed by Gromov–Thurston, Mostow–Siu, Deraux, or Stover–Toledo;
- their fundamental groups are Gromov hyperbolic and virtually compact special (in particular, they satisfy properties (1)-(9) in Theorem 1.2).

Remark 1.4. Thanks to the solution of the Borel Conjecture for closed aspherical n-manifolds with Gromov hyperbolic fundamental group in dimension n ≥ 5 (see Bartels-Lück in [BL12]), the fundamental groups of these manifolds provide examples of Gromov hyperbolic groups that are not isomorphic to lattices in SO(n, 1) or the other real simple Lie groups of rank 1. While it is not a priori clear from their construction whether these groups are linear, they actually turn out to be virtually compact special, hence linear over ℤ and residually finite. We note that Giralt proved in [Gir17] that the fundamental groups of the Gromov-Thurston manifolds are also virtually compact special.

Similarly, other applications obtained by Ontaneda in [Ont20] can be taken to have additional algebraic features. For example, we have the following version of Corollary 2 of [Ont20].
Corollary 1.5. Let $\varepsilon > 0$. The cohomology ring of any finite CW–complex embeds in the cohomology ring of a closed Riemannian manifold which has sectional curvatures in $[-1 - \varepsilon, -1]$ and whose fundamental group is Gromov hyperbolic and virtually compact special (hence satisfies properties (1)-(9) in Theorem 1.2). In particular, it can be embedded into the cohomology ring of a Poincaré Duality subgroup of $\text{SL}_N(\mathbb{Z})$ (for $N$ large).

1.2.2. Pathological aspherical manifolds. Davis and Januszkiewicz used the hyperbolization procedures to construct aspherical manifolds exhibiting a variety of pathological behavior (see [DJ91]). As a consequence of our Theorem 1.2, these examples can now be constructed to have the added property that their fundamental groups are virtually compact special, hence satisfy properties (1)-(9) from Theorem 1.2. For the convenience of the reader, we collate some of their examples.

Corollary 1.6. It is possible to construct (topological) manifolds of the following types which are piecewise hyperbolic and locally $\text{CAT}(-1)$.

- A closed 4-manifold which is not homotopy equivalent to any PL 4-manifold (see [DJ91, §5a]).
- For $n = 4k, k \geq 2$, a closed $n$-manifold which is not homotopy equivalent to any smooth manifold (see [BLW10, Example 5.2]).
- For $n \geq 5$, a closed $n$-manifold whose universal cover is not homeomorphic to $\mathbb{R}^n$ (see [DJ91, §5b]).
- For $n \geq 5$, a closed $n$-manifold whose universal cover is homeomorphic to $\mathbb{R}^n$, but whose boundary at infinity is not homeomorphic to $S^{n-1}$ (see [DJ91, §5c]).

Moreover, in all these examples, the fundamental groups of the manifolds are Gromov hyperbolic and virtually compact special (in particular, they satisfy properties (1)-(9) from Theorem 1.2).

Remark 1.7. Concerning the first example in Corollary 1.6, taking products with tori yields examples in all dimensions $n \geq 4$ of closed aspherical $n$-manifolds not homotopy equivalent to any PL $n$-manifold. These manifolds will have fundamental group which is linear over $\mathbb{Z}$, but when $n \geq 5$ will only support a locally $\text{CAT}(0)$ metric due to the product structure. It would be interesting to produce examples in dimensions $n \geq 5$ which support locally $\text{CAT}(-1)$ metrics.

1.2.3. Representing cobordism classes. As another application, we can obtain representatives for cobordism classes that are both topologically and algebraically nice.

Corollary 1.8. Let $M$ be an arbitrary closed smooth manifold. Then $M$ is cobordant to an aspherical manifold $M'$, where $\pi_1(M')$ is a Gromov hyperbolic and virtually compact special (in particular, it satisfies properties (1)-(9) from Theorem 1.2).

Following an idea of Gromov (see [Gro87; Pau91]), one lets $K$ be the cone over a smooth triangulation $\tau$ of $M$. Then we apply the strict hyperbolization $\mathcal{H}(K)$, and note that since hyperbolization preserves links, the point $p \in \mathcal{H}(K)$ corresponding to the cone point will have link a copy of $\tau$. Thus, removing a small neighborhood of $p$ leaves us with a cobordism $W$ between $M$ and $M' := \mathcal{H}(\tau)$. Our Theorem 1.2 then applies to $M'$. Note that $\pi_1(W)$ itself contains $\pi_1(M)$, hence might not be linear (for instance, if $\pi_1(M)$ is a non-linear group).
Remark 1.9. Thom’s work showed that oriented cobordism classes are rationally represented by products of even dimensional complex projective spaces (see [MS74, Section 17]). So every smooth oriented closed manifold has a multiple which is cobordant to a non-negatively curved Riemannian manifold. In analogy, combining Davis–Januszkiewicz–Weinberger [DJW01], Charney–Davis [CD95], and Ontaneda [Ont20], one obtains that every smooth oriented closed manifold is cobordant to a strictly negatively curved Riemannian manifold.

In dimensions $\geq 5$, the Borel Conjecture is known to hold for aspherical manifolds with Gromov hyperbolic groups (see Bartels–Lück [BL12]). As such, the topological manifold $M' := \mathcal{H}(\tau)$ is completely determined, up to homeomorphism, by its fundamental group. So the discussion above in principle reduces the study of cobordism classes of manifolds of dimension $n \geq 5$, to the study of the corresponding $\pi_1(M')$. Our corollary further reduces it to the linear case.

Remark 1.10. More generally, Corollary 1.8 works for a PL manifold, or even for a triangulable topological manifold. Note that in all dimension $n \geq 4$ there exist compact topological manifolds that are not triangulable (see [Man16]).

1.2.4. Prescribing the Gromov boundary. The groups obtained by strict hyperbolization are Gromov hyperbolic groups, so it is natural to ask what their Gromov boundary looks like. For example, the groups obtained by Riemannian hyperbolization in [Ont20] (see §1.2.1) are fundamental groups of smooth Riemannian manifolds of negative curvature, hence their Gromov boundaries are spheres of the appropriate dimensions.

Corollary 1.11. Let $n \geq 1$, and let $M$ be a closed connected orientable PL $n$-manifold that bounds a compact orientable PL $(n + 1)$-manifold. Then there exists a Gromov hyperbolic group $G$ such that

- the Gromov boundary of $G$ is homeomorphic to the tree of manifolds $\mathcal{X}(M)$;
- $G$ is virtually compact special (hence satisfies (1)-(9) in Theorem 1.2).

The groups in this statement are the ones obtained by Świątkowski in [Ś20] via strict hyperbolization of certain pseudomanifolds in which the link of a point is either a sphere or a copy of the manifold $M$. The tree of manifolds $\mathcal{X}(M)$ is a compact metrizable space which is obtained, roughly speaking, as a certain limit of connected sums of copies of $M$.

1.2.5. Manifolds with exotic symmetries. The hyperbolization procedures satisfy a certain functorial property: automorphisms of the simplicial complex $K$ induce isometries of the hyperbolized complex $\mathcal{H}(K)$. This has been used by various authors to produce closed manifolds with interesting symmetries.

For example, if $G$ is a Gromov hyperbolic group which is a Poincaré Duality group over $\mathbb{Z}$, an easy application of Smith theory shows that the fixed subgroup $G^\sigma$ of an involution $\sigma \in \text{Aut}(G)$ is still a Poincaré Duality group, but over $\mathbb{Z}_2$. Farrell–Lafont in [FL04] used an exotic symmetry produced via strict hyperbolization, to give examples whose fixed subgroups are not Poincaré Duality over $\mathbb{Z}$. Our results now show that these examples can also be chosen to satisfy properties (1)-(9) in Theorem 1.2.

For another application, recall that in their seminal paper [BC00], Baum–Connes defined a trace map $tr : K_0(C^*_r G) \rightarrow \mathbb{R}$, where $C^*_r G$ is the reduced $C^*$-algebra of the discrete group $G$. They also formulated the trace conjecture, which predicted
that when $G$ is a group with torsion, the image of the trace map is contained in the additive subgroup of $\mathbb{Q}$ generated by $1/n$, where $n$ ranges over the order of finite subgroups of $G$. A counterexample to this conjecture was constructed by Roy [Roy99], using the Davis–Januszkiewicz (non-strict) hyperbolization procedure. She constructed a group $G$ whose only finite subgroups are isomorphic to $\mathbb{Z}_3$, and an element in $K_0(C^*_r G)$ whose trace equals $-1105/9$. Nevertheless, there is always the possibility that the original Baum–Connes trace conjecture might hold for certain restricted classes of groups. The computations carried out by Roy ([Roy99], pgs. 210-213) apply verbatim if one instead uses the Charney–Davis strict hyperbolization, so our results have the following consequence.

**Corollary 1.12.** There exists a Gromov hyperbolic group $G$ whose only finite subgroups are isomorphic to $\mathbb{Z}_3$, but where the image of the trace map contains $-1105/9$. Moreover, this group satisfies properties (1)-(9) in Theorem 1.2. In particular, the original Baum–Connes trace conjecture does not hold for the classes of groups (1)-(9) in Theorem 1.2.

**Remark 1.13.** Lück formulated a refinement of the original Baum–Connes trace conjecture: the image of the trace map is contained in the subring $\mathbb{Z}[1/|\text{Fin}(G)|]$, obtained from $\mathbb{Z}$ by inverting all the orders of finite subgroups of $G$. Lück showed that this refined Trace Conjecture holds for any group that satisfies the Baum–Connes Conjecture (see [Lü02]). In the subsequent literature, this refined version is what is commonly referred to as the Trace Conjecture. For Gromov hyperbolic groups, the Baum–Connes Conjecture was established by Lafforgue (see [Laf02]). Thus the group appearing in our Corollary 1.12 satisfies the refined trace conjecture.

### 1.3. Virtual algebraic fibering.

In this section we present new applications of a more algebraic flavor. We say a group $G$ *algebraically fibers* if it admits a surjective homomorphism to $\mathbb{Z}$ with finitely generated kernel. We say it *virtually algebraically fibers* if it has a finite index subgroup that algebraically fibers. Agol introduced the notion of *residually finite rationally solvable* (or RFRS) group in [Ago08] as a major ingredient in the solution of the Virtual Haken Conjecture and Virtual Fibering Conjecture. Kielak proved in [Kie20, Theorem 5.3] that a finitely generated virtually RFRS group virtually algebraically fibers if and only if its first $L^2$–Betti number vanishes. Fisher has extended this result in [Fis21, Theorem 6.9] to relate the vanishing of higher $L^2$–Betti numbers of $G$ to higher finiteness properties of the kernel of a virtual algebraic fibration.

All of the groups constructed in this paper via strict hyperbolization are virtually compact special, hence virtually RFRS, see [Ago08, Corollary 2.3]. In some cases, it is possible to prove vanishing of many $L^2$–Betti numbers (for instance for all the examples obtained by Ontaneda in [Ont20], provided the curvatures are sufficiently pinched; see below for details). Hence, we get several new examples of virtually compact special Gromov hyperbolic groups that admit a virtual algebraic fibration, whose kernel has good algebraic finiteness properties. On the other hand, these groups can be seen to be incoherent, and in some cases it is possible to see that the kernel of a virtual algebraic fibration is itself a witness to incoherence (i.e. is finitely generated but not finitely presented).

Before providing the details for our case, we note that similar arguments also work for arithmetic hyperbolic manifolds of simple type and for Gromov–Thurston
manifolds. These are known to be virtually specially cubulated (hence RFRS) by [HW12] and [Gir17] respectively.

1.3.1. Kernels with good algebraic finiteness properties. We start by constructing Gromov hyperbolic groups that virtually algebraically fiber, and are not isomorphic to groups that were previously known to have this property.

Corollary 1.14. For all $n \geq 4$ there is a closed Riemannian $n$-manifold $M$ with negative sectional curvatures and such that

- $\pi_1(M)$ virtually algebraically fibers;
- $\pi_1(M)$ is Gromov hyperbolic and virtually compact special (hence satisfies (1)-(9) in Theorem 1.2);
- $\pi_1(M)$ is not isomorphic to a uniform lattice in $\text{SO}(n,1)$ (or other real simple Lie group of rank 1), or to the fundamental group of a Gromov-Thurston, Mostow-Siu, Deraux, or Stover-Toledo manifold.

The manifolds in this statement are the ones constructed by Ontaneda in [Ont20] (see Corollary 1.3 above). As a result of our Theorem 1.2 the fundamental group of such a manifold $M$ is virtually compact special, and in particular it is virtually RFRS. Moreover, $M$ can be chosen to have sectional curvatures pinched in the interval $[-1 - \varepsilon, -1]$ for an arbitrarily small $\varepsilon > 0$. By a result of Donnelly and Xavier (see [DX84, §4], and also [JX00, Theorem 2.3]), if the curvatures of $M$ are sufficiently pinched (i.e. $\varepsilon$ is sufficiently small with respect to the dimension $n$), then $M$ does not have any non-zero $L^2$–harmonic $p$-forms, for $p$ in a certain range. In particular, $b_2^0(\pi_1(M)) = 0$. By [Kie20] we see that if $\varepsilon$ is small enough then $\pi_1(M)$ virtually algebraically fibers.

Furthermore, one can pinch the curvatures even more to force the vanishing of the $L^2$–Betti numbers for $p = 0, 1, \ldots, \lfloor \frac{n}{2} \rfloor - 1$. In particular, using results from [Fis21] one can obtain examples in which $\pi_1(M)$ virtually algebraically fibers with kernel of type $FP_{\lfloor \frac{n}{2} \rfloor - 1}(\mathbb{Q})$. Also note that in the even dimensional case $M$ can be chosen to satisfy the (weak) Hopf conjecture, i.e. $(-1)^{\frac{n}{2}}\chi(M) \geq 0$.

1.3.2. Kernels witness incoherence in dimension 4. We have discussed how to use strict hyperbolization to obtain Gromov hyperbolic groups that virtually algebraically fiber with kernel of type $FP_{\lfloor \frac{n}{2} \rfloor - 1}(\mathbb{Q})$. On the other hand, these kernels should not be expected to have better finiteness properties. Indeed, in the context of the previous paragraph, we can show that in dimension $n = 4$ these kernels are not finitely presented (i.e. not of type $F_2$).

To see this, notice that Chern–Weil theory implies that the Euler characteristic of a closed negatively curved 4-manifold is strictly positive (see [Che55]). This prevents the kernel of an algebraic fibration of $\pi_1(M)$ from being finitely presented, as we now describe. We thank Genevieve Walsh for sharing the following argument with us. (This appears in [KVW21].)

Lemma 1.15. Let $M$ be a closed aspherical 4-manifold such that $\chi(M) \neq 0$. If $\pi_1(M)$ virtually algebraically fibers, then $\pi_1(M)$ is incoherent (the kernel is not finitely presented).

Proof. Suppose $\pi_1(M)$ virtually algebraically fibers, and let $G$ be the finite index subgroup of $\pi_1(M)$ which surjects to $\mathbb{Z}$ with finitely generated kernel $K$. Notice that $G$ is a $PD_4$ group with $\chi(G) \neq 0$ (since Euler characteristic is multiplicative
by index), and that $Z$ is a $PD_1$ group. Assume by contradiction that $K$ is finitely presented (i.e. type $F_2$). Then $K$ is in particular of type $FP_2$, and it follows from [HK07, Corollary 1.1] that $K$ is a $PD_3$ group. In particular $K$ has finite homological type (and the same is true for $Z$). So, by the properties of Euler characteristics on short exact sequences (see [Bro82, Chapter IX, 7.3(d)]) we can conclude that $\chi(G) = \chi(K)\chi(Z) = 0$. This contradicts the fact that $\chi(G) \neq 0$. □

An alternative argument for this Lemma, under the additional assumption that $\pi_1(M)$ is virtually RFRS, was shared with us by Kevin Schreve.

Proof. In the same set up, if by contradiction $K$ is finitely presented, then it is in particular of type $FP_2(\mathbb{Q})$. So, since $G$ is also virtually RFRS, by [Kie20; Fis21] we get that $b_1^{(2)}(G) = b_2^{(2)}(G) = 0$. But $G$ is a $PD_4$ group, so by duality this implies that all $L^2$-Betti numbers vanish. This gives $\chi(G) = 0$, which is again absurd. □

As a result we obtain the following statement. The manifold in it is once again one of the manifolds obtained by Ontaneda, with curvatures sufficiently pinched.

Corollary 1.16. There exists a closed 4-dimensional Riemannian manifold $M$ with negative sectional curvatures and such that

- $\pi_1(M)$ is incoherent (it virtually algebraically fibers with non finitely presented kernel);
- $\pi_1(M)$ is Gromov hyperbolic and virtually compact special (hence satisfies $(1)$-$(9)$ in Theorem 1.2);
- $\pi_1(M)$ is not isomorphic to a uniform lattice in $SO(n,1)$ (or other real simple Lie group of rank 1), or to the fundamental group of a Gromov-Thurston, Mostow-Siu, or Stover- Toledo manifold.

Remark 1.17. The situation in dimension 4 is quite different from that in dimension 5. Indeed, Italiano, Martelli, Migliorini in [IMM21] obtained a 5-dimensional cusped hyperbolic manifold that fibers over the circle. Its fundamental group algebraically fibers, with kernel of finite type, and in particular finitely presented. The hyperbolic groups obtained by suitable Dehn filling on these examples were shown to fiber with kernel of finite type. Moreover, a recent preprint of Groves and Manning shows that some of these groups are virtually compact special (see [GF22]).

Remark 1.18. When $n \geq 5$, the groups obtained by strict hyperbolization done with a sufficiently large piece (as in Ontaneda) contain subgroups isomorphic to uniform arithmetic lattices in $SO(4,1)$. The incoherence of these subgroups (see [KPV08; Ago08; Kap13]) gives incoherence of the hyperbolized groups, but these subgroups are not fibers themselves (for instance because they are quasiconvex). Notice that this approach does not work in dimension $n = 4$, as uniform lattices in $SO(3,1)$ are coherent.

Structure of the paper. This paper is structured as follows. In §1 we present the motivation, the context, the statements, and the major applications of our results. In §2, we provide combinatorial and metric background about cell complexes and hyperbolization procedures. §3 is devoted to a description of Charney–Davis strict hyperbolization procedure for a cubical complex $X$. In particular, we study the geometry of the universal cover $\tilde{X}_\Gamma$ of the hyperbolized complex $X_\Gamma$ in terms of
a certain collection of convex subspaces called mirrors. This provides a graph
of spaces decomposition of $X_{\Gamma}$. In §4 we construct and study a dual CAT(0)
cubical complex $C(\tilde{X}_{\Gamma})$. Finally in §5 we study the action of the hyperbolized
group $\Gamma_{X} = \pi_{1}(X_{\Gamma})$ on this dual cubical complex $C(\tilde{X}_{\Gamma})$, and prove that $\Gamma_{X}$ is
virtually compact special.

Common terminology and notation. The numbers in parenthesis refer to the
section(s) in which each item is introduced or discussed.

- The hyperbolizing lattice $\Gamma$ (§3.1) and the cubical complex $X_{\Gamma}$.
- The hyperbolized complex $X_{\Gamma}$ (§3.2), and its universal cover $\tilde{X}_{\Gamma}$ (§3.3).
- The hyperbolized cube $\square^{n}_{\Gamma}$ (§3.1), and its universal cover $\tilde{\square}^{n}_{\Gamma}$ (§3.3).
- The hyperbolized groups $\Gamma_{\square^{n}} = \pi_{1}(\square^{n}_{\Gamma})$ and $\Gamma_{X} = \pi_{1}(X_{\Gamma})$ (§3.2).
- The folding map $f : X \to \square^{n}$ of a foldable complex (§2.2), and the induced
  map $f_{\Gamma} : X_{\Gamma} \to \square^{n}_{\Gamma}$ on the hyperbolized complex (§3.2).
- The Charney-Davis map $g : \square^{n}_{\Gamma} \to \square^{n}$, and the induced map $g_{X} : X_{\Gamma} \to X$
on the hyperbolized complex (§3.2).
- The dual cubical complex $C(\tilde{X}_{\Gamma})$ (§4).
- A face $F \subseteq \square^{n}$.
- A cube $C \subseteq X$ (if $X$ is a cubical complex) or $Q \subseteq C(\tilde{X}_{\Gamma})$.
- A cell $\sigma \subseteq X_{\Gamma}$, $\tilde{X}_{\Gamma}$ (§3.5), $\boxtimes^{n}$, $\square^{n}_{\Gamma}$, $\tilde{\square}^{n}_{\Gamma}$ (§3.3).
- A tile $\tau$ in $\tilde{X}_{\Gamma}$ (§3.3), and the dual tile $C(\tau)$ in $C(\tilde{X}_{\Gamma})$ (§4.2).
- A mirror $M$ in $\tilde{X}_{\Gamma}$ (§3.4), and the dual mirror $C(M)$ in $C(\tilde{X}_{\Gamma})$ (§4.3).
- An edge–path $p$ in $C(\tilde{X}_{\Gamma})$, its length $\ell(p)$, its height $h(p)$ (§4.1, §4.2), the
  number of $(p, M)$–crossings $m(p, M)$ with respect to a mirror $M$, and its
  total mirror complexity $m(p)$ (§4.3).

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2. Cell complexes and hyperbolization procedures

We collect in this section some background material used in our constructions.
In §2.1 we review the basics about cell complexes, and in §2.2 we focus on foldable
complexes, i.e. complexes that can be folded down to a standard simplex or cube.
In §2.3 we introduce a general template for the study of hyperbolization procedures
for foldable complexes. In §2.4 we review a specific hyperbolization procedure due
to Gromov.

2.1. Combinatorial and metric geometry of cell complexes. In this section
we collect background material about cell complexes, mainly to fix notation and
terminology; for a detailed treatment see [BH99, §L7, §L15]. Let us denote by $M_{k}^{n}$
the simply connected Riemannian manifold of dimension $n$ and constant sectional
curvature $k$: for instance $M_{1}^{n} = S^{n}$ is the round sphere, $M_{0}^{n} = E^{n}$ is the Euclidean
space, and $M^n_{n-1} = \mathbb{H}^n$ is the hyperbolic space. An isometrically embedded copy of $M_k^n$ inside $M_k^n$ will be called a $d$–plane, or a hyperplane if $d = n - 1$.

2.1.1. Cells. A cell in $M_k^n$ is defined to be the convex hull of a finite set of points; if $k > 0$ we are going to also require that it is contained in an open ball of radius $\frac{\pi}{2\sqrt{k}}$. The dimension of a cell $C$ is the smallest $d$ such that $C$ is contained in a $d$–plane. A cell of dimension $d$ will also be called a $d$–cell. The interior of $C$ is its interior inside this $d$–plane. A face $F$ of $C$ is a subspace of the form $F = H \cap C$ where $H$ is a hyperplane such that $C$ lives in one of the two closed half–spaces bounded by $H$, and $H \cap C \neq \emptyset$. A face is itself a cell, and we call vertices and edge of $C$ the faces of dimension 0 and 1 respectively.

2.1.2. Cell complexes. An $M_k^n$–cell complex is a topological space $X$ obtained by gluing together cells from $M_k^n$ by isometries of their faces, in such a way that each cell embeds in $X$ and the intersection of any two cells is either empty or a cell. Notice that this definition is more strictly restrictive than the one in [BH99, Definition 1.7.37] (which allows one to glue two faces of the same cell), and the one in [CD95] (in which cells are allowed to intersect in a proper union of faces). Both conditions can be satisfied by performing a cellular subdivision. On the other hand, we do not require cell complexes to be locally compact at this stage, i.e. a vertex can be contained in infinitely many cells.

We call an $M_k^n$–cell complex simply a cell complex when we do not need to keep track of $M_k^n$. For instance we will denote by $\Delta^n$ the standard $n$–simplex and by $\Box^n = [0, 1]^n$ the standard $n$–cube; these are cells in $M_0^n$. A simplicial complex is a cell complex obtained by gluing simplices, and a cubical complex is a cell complex obtained by gluing cubes.

The dimension of a cell complex is the maximum dimension of its cells. We say that an $n$–dimensional cell complex is homogeneous if every cell is contained in a cell of dimension $n$, and that it is without boundary if every $(n-1)$–cell is contained in at least two different $n$–cells. For all $k = 0, \ldots, n$, the $k$–skeleton of $X$ is the subspace consisting of all the cells of dimension at most $k$, and will be denoted by $X^{(k)}$. A subcomplex of $X$ is a closed subspace $Y \subseteq X$ which is a union of cells of $X$. If $X$ and $Y$ are cell complexes, a continuous function $f : X \to Y$ is a combinatorial map if for every cell $C$ of $X$ we have that $f$ is a homeomorphism from $C$ to a cell $f(C)$ of $Y$.

Given a cell complex $X$, its barycentric subdivision $\mathcal{B}(X)$ is the simplicial complex whose $k$–simplices correspond to strictly ascending sequences of faces $F_0 \subset \cdots \subset F_k$ of $X$. There exists a natural (non–combinatorial) homeomorphism $X \to \mathcal{B}(X)$. We refer to [BH99, §I.7.44-48] for more details. Similarly, if $X$ is a cubical complex, then its cubical subdivision is the cubical complex obtained by subdividing each $n$–cube along midcubes into $2^n$ cubes.

Remark 2.1. By definition, a cell is compact, it has finitely many faces, and it can be realized as the intersection of finitely many closed half–spaces (see [BH99, §I.7]). We want to warn the reader that one of the main object under investigation in this paper (see §3.5) is obtained by gluing together certain “generalized cells”, i.e. subsets of $\mathbb{H}^n$ which are given by the intersection of an infinite but locally finite collection of closed half–spaces. These subsets are convex but not compact, so the resulting space is not strictly speaking a cell complex. However, some of the usual
tools for the study of cell complexes can be applied in this context (e.g. links). We will highlight the subtleties in the construction whenever relevant.

2.1.3. Links. Let $X$ be a cell complex. We define the link of points and cells as follows (see [BH99, §I.7] for more details). Let $p$ be a point of an $n$–cell $C \subseteq M^n_k$. Then we define the link $\text{lk}(p, C)$ to be the space of unit vectors in the tangent space at $p$ inside $C$. Measuring the angle between vectors endows $\text{lk}(p, C)$ with a natural length metric, which makes it isometric to an $(n – 1)$–cell in $S^{n-1}$. The link $\text{lk}(p, X)$ of $p$ in $X$ is then defined by gluing together the links $\text{lk}(p, C_i)$, where $C_i$ ranges over the cells of $X$ containing $p$. This is naturally an $M^{n-1}_k$–cell complex. If $Y$ is a sufficiently regular subspace of $X$ containing $p$ (e.g. a subcomplex), then the link $\text{lk}(p, Y)$ is defined analogously, restricting to vectors along $Y$.

Let $F$ be a $d$–face of an $n$–cell $C \subseteq M^n_k$. Then we define the link $\text{lk}(F, C)$ to be the subspace of unit vectors in the tangent space at an interior point of $F$, which are pointing into $C$ and are orthogonal to $F$. As before, this is naturally an $(n – k – 1)$–cell in $S^{n-1}$. The link $\text{lk}(C, X)$ of a $k$–cell $C \subseteq X$ is then defined by gluing together the links $\text{lk}(C, C_i)$, where $C_i$ ranges over the cells of $X$ containing $C$. It is naturally an $M^{n-k-1}_k$–cell complex. Finally, if $Y \subseteq X$ is a subcomplex of $X$ containing $C$, the link $\text{lk}(C, Y)$ of $C$ in $Y$ is defined analogously, by restricting to the cells of $Y$ that contain $C$. Observe that if $X$ is a simplicial or cubical complex, then the link of a $d$–cell $C$ is a simplicial complex in which vertices are given by the $(d + 1)$–cells containing $C$, and in which $m + 1$ vertices span an $m$–simplex if and only if the corresponding cells are contained in a $(d + m + 1)$–cell.

2.1.4. Spaces and complexes of bounded curvature. We will consider the usual notions of curvature for metric spaces, such as being locally $\text{CAT}(k)$ or Gromov hyperbolic (see [BH99, §II.1, §III.H.1] for more details). In particular, we will say a space is non-positively curved if it is locally $\text{CAT}(0)$, and negatively curved if it is locally $\text{CAT}(k)$ for some $k < 0$. Note that if $k < 0$ then a $\text{CAT}(k)$ space is in particular Gromov hyperbolic (see [BH99, Proposition III.H.1.2]), and that if $k \leq 0$ then a $\text{CAT}(k)$ space is uniquely geodesic (see [BH99, Proposition II.1.4]). Whenever $x, y$ are points in a uniquely geodesic space, we denote by $[x, y]$ the unique geodesic between them.

Let $X$ be a $M^n_k$–cell complex. Each cell can be naturally endowed with a metric from $M^n_k$, and these can be glued together to make $X$ into a complete geodesic metric space, as soon as there are only finitely many isometry classes of cells in $X$ (see [BH99, Theorem I.7.50]). When equipped with this metric, $X$ is said to be a cell complex of piecewise constant curvature $k$; we say it is piecewise spherical, Euclidean, or hyperbolic when $k = 1, 0, -1$ respectively. If not otherwise specified, a simplicial or cubical complex is always endowed with its standard piecewise Euclidean metric.

It is natural to ask for conditions under which a complex of piecewise constant curvature is globally a space of bounded curvature, namely a $\text{CAT}(k)$ space. For cubical complexes this is completely controlled by the links of vertices. In a cubical complex, cells are isometric to the standard Euclidean cube $\Box^n = [0, 1]^n$, so the link of a vertex is a piecewise spherical simplicial complex, in which all edges have length $\frac{\pi}{2}$. The following is known as Gromov’s link condition (see [BH99, Theorems II.5.18, II.5.20]). A simplicial complex is flag if any $k + 1$ pairwise adjacent vertices span a $k$–simplex.
Lemma 2.2. Let $X$ be a cubical complex. Then the following are equivalent.

1. $X$ is non–positively curved (i.e. locally CAT(0)).
2. The link of each vertex is a flag simplicial complex.
3. The link of each vertex is a CAT(1) simplicial complex.

2.2. Foldable complexes. Here we consider the notion of foldability for simplicial and cubical complexes that we will require later. The first definition is essentially from [BSa99, §1] (but see also [CD95, Definition 7.2], and [Xie04] for a more recent discussion).

A simplicial (respectively cubical) $n$-dimensional complex $X$ is foldable if it admits a combinatorial map $f : X \to \triangle^n$ (respectively $f : X \to \square^n$) such that its restriction to each cell of $X$ is injective. Such a map will be called a folding for $X$. Notice that in a foldable complex the cells are necessarily embedded. This is the main reason why we have incorporated this condition in the definition of cell complex in §2.1.

Foldability has some immediate consequences. If $X$ is foldable, and $p : Y \to X$ is a combinatorial map which is injective on each cell, then $Y$ is foldable too. In particular any covering of a foldable complex is foldable. Moreover if $X$ is foldable, then the links of cells of codimension 2 are bipartite graphs. We collect below some examples in the cubical case; analogous ones can be constructed for the simplicial case.

Example 2.3 (Foldable and not foldable cubical complexes).

1. A graph is foldable if and only if it is bipartite (Figure 1, left).
2. The rose $R_m$ consisting of $m$ squares with a vertex in common is foldable if and only if $m$ is even (Figure 1, right).
3. Foldability of $X$ implies that links of codimension 2 cells are bipartite. However, foldability is not completely determined by this property. For example, let $X$ be the cubical complex obtained by taking the product $\partial \Delta^2 \times \mathbb{R}$, where $\mathbb{R}$ is endowed with the standard cell structure induced by $\mathbb{Z}$. Then the links of vertices are cycles of length 4 (hence they are bipartite), but $X$ is not foldable; notice that the universal cover of $X$ identifies the square complex defined by $\mathbb{Z}^2$ in $\mathbb{R}^2$, which is foldable (compare [BSa99, Lemma 1.2]).

A main source of foldability comes from barycentric subdivisions; the following is well-known (see [BSa99, Lemma 2.1]), we include a proof for completeness (see left of Figure 2 for an example).
Lemma 2.4. If $X$ is a cell complex, then $B(X)$ is a foldable simplicial complex.

Proof. Let $X$ have dimension $n$, and consider the simplex spanned by $\{0, \ldots, n\}$; this is just the standard simplex $\Delta^n$. Then we can define a map $f : B(X) \to \Delta^n$ by sending a vertex of $B(X)$ to the number which is equal to the dimension of the corresponding cell in $X$. \hfill \Box

Figure 2. The barycentric subdivision of the rose of 3 squares is a foldable simplicial complex, but its cubical subdivision is not a foldable cubical complex.

On the other hand, if $X$ is a non–foldable cubical complex of dimension at least 2, then its cubical subdivision is still non–foldable (see Figure 2, right). In §2.4 we will review Gromov’s construction and show that it can be used to turn any cubical complex into a foldable one (mildly changing the topology).

2.3. Hyperbolization procedures. In this section we set a framework for the study of certain constructions, which take a cell complex as input and return a non-positively curved space as output. The resulting space is in particular always aspherical, so the topology of the original complex is altered. What is interesting is that this can happen in a controlled way that allows to preserve some features of the original complex. Constructions of this type are generally known as hyperbolization procedures (or asphericalization procedures). They were first introduced by Gromov (see [Gro87, §3.4.A]), and then popularized by several authors (see [DJ91; Pan91; CD95; DFL14; Ont20]).

All the hyperbolization procedures we will consider in this paper are obtained by different incarnations of the same abstract construction, which we now review briefly, referring the reader to [Wil63] or [DJ91, §1] for more details. The naive idea is to fix some topological space $S$ and then replace every top-dimensional cell of a complex $X$ by a copy of $S$. For this gluing to be well–defined, it is common to assume that both $X$ and $S$ are equipped with chosen maps to a reference space.

For concreteness let us consider the following set up. Let us denote by $\sigma^n$ the standard simplex $\Delta^n$ or the standard cube $\square^n$, and let us fix a topological space $S$, equipped with a continuous map $g : S \to \sigma^n$, and a foldable simplicial or cubical complex $X$, equipped with a fixed folding $f : X \to \sigma^n$. One then considers the fibered product $\mathcal{H}_S(X) = \{(x, s) \in X \times S \mid f(x) = g(s)\}$, i.e. the space obtained via the pullback square in Figure 3.

Note that the construction endows $\mathcal{H}_S(X)$ with natural continuous maps $g_X : \mathcal{H}_S(X) \to X$ and $f_S : \mathcal{H}_S(X) \to S$, which are just the restrictions of the projections onto the factors of $X \times S$, and which make the diagram commute. Properties
Figure 3. A template for hyperbolization procedures

of the pair \((S, g)\) will result in properties of the space \(\mathcal{H}_S(X)\), and the art of hyperbolization consists in crafting a pair \((S, g)\) which yields some interesting properties on \(\mathcal{H}_S(X)\). For a trivial example, consider the case \(S\) consists of a single point. Then \(\mathcal{H}_S(X)\) is just the discrete set \(f^{-1}(g(S))\).

The following lemma identifies a mild condition under which the space \(\mathcal{H}_S(X)\) looks like a collection of copies of \(S\) (compare the remark on page 321 of [Wil63]). We explicitly remark that we do not assume \(S\) to be compact until part (3) of this lemma. This will be relevant in §3.3 for the study of a certain combinatorial decomposition of a space into non-compact pieces.

**Lemma 2.5.** Let \(g : S \to \sigma^n\) be surjective, and let \(C \subseteq X\) be an \(n\)–cell. Then the following hold.

1. The map \(f_S\) restricts to a homeomorphism \(\varphi : g_X^{-1}(C) \to S\).
2. The map \(\varphi^{-1} \circ f_S : \mathcal{H}_S(X) \to g_X^{-1}(C)\) is a retraction.
3. If \(X\) and \(S\) are compact, then \(\mathcal{H}_S(X)\) is compact too.

**Proof.** Let us denote by \(f_C\) the restriction of the folding map \(f\) to \(C\). Note that \(f_C : C \to \sigma^n\) is a homeomorphism. To prove (1), let \(\varphi : g_X^{-1}(C) \to S\) be the restriction of \(f_S\) to \(g_X^{-1}(C)\). Then \(\varphi\) is continuous, because it is just the restriction of the projection \(X \times S \to S\). Injectivity and surjectivity of \(\varphi\) follow respectively from those of \(f_C\). To conclude, we construct an explicit continuous inverse. Consider the map \(\lambda : S \to C, \lambda(s) = f_C^{-1}(g(s))\). Notice it is well-defined (because \(g\) is surjective), and continuous. Then the map \(\psi : S \to g_X^{-1}(C) \subseteq X \times S, \psi(s) = (\lambda(s), s)\) provides a continuous inverse to \(\varphi\).

Now (2) follows from (1), as every element of \(g_X^{-1}(C)\) is of the form \((\lambda(s), s)\). Finally, to prove (3), observe that if \(X\) is compact, then it consists of finitely many \(n\)–cells. As a result of the previous argument, \(\mathcal{H}_S(X)\) is covered by finitely many copies of the compact space \(S\), hence it is compact. \(\square\)

Depending on the applications in which they are interested, authors differ on what additional geometric conditions they require on the association \(X \to \mathcal{H}(S, X)\), hence they start with different spaces \((S, g)\). We refer the reader to [DJ91] for a very general treatment of how properties of \((S, g)\) imply properties of \(\mathcal{H}(S, X)\). Some commonly required conditions are the following

1. (Hyperbolicity): \(\mathcal{H}(S, X)\) admits a non-positively curved metric.
2. (Functoriality): if \(Z \subseteq X\) is a locally convex subcomplex, then \(\mathcal{H}(S, Z) \subseteq \mathcal{H}(S, X)\) is a locally convex subspace.
(3) (Local structure): if $C \subseteq X$ is an $n$–cell, then $\mathcal{H}_S(C)$ is an $n$-manifold with boundary and corners, and $\text{lk}(\mathcal{H}_S(C), \mathcal{H}_S(X)) \cong \text{lk}(C, X)$. In particular, if $X$ is a manifold, then $\mathcal{H}_S(X)$ is a manifold too.

(4) (Homology surjectivity): the map $g_X : \mathcal{H}_S(X) \rightarrow X$ induces a surjection on homology.

The association $X \rightarrow \mathcal{H}_S(X)$ is then called the hyperbolization procedure defined by $(S, g)$. We call $S$ the hyperbolizing cell, and $\mathcal{H}_S(X)$ the hyperbolized complex. Despite the name (established in the literature), the output $\mathcal{H}_S(X)$ of a hyperbolization procedure is a metric space which a priori is just non-positively curved. A strict hyperbolization is one for which $\mathcal{H}_S(X)$ is negatively curved. In this paper we will consider a (non–strict) hyperbolization for simplicial complexes due to Gromov (see §2.4), and a strict hyperbolization for cubical complexes due to Charney and Davis (see §3).

Remark 2.6. If $(S, g)$ is a given hyperbolizing cell, $g : S \rightarrow \sigma^n$ is surjective, and $F \subseteq \sigma^n$ is a closed face of the $n$–cell $\sigma^n$, then the subspace $g^{-1}(F)$ will be called a face of $S$. The dimension of a face of $S$ is defined to be simply the dimension of the corresponding face of $\sigma^n$. Note that a face of $S$ does not need to be connected. When this happens, $\mathcal{H}_S(X)$ may fail to be simply connected, even if both $X$ and $S$ are. For some interesting examples, see [DJ91, 1b.1], or consider the elementary one in Figure 4. Despite their non trivial role in the construction, most of the times the maps $f$ and $g$ are omitted from the notation.

**Figure 4.** Example for Remark 2.6: a face of the hyperbolizing cell $S$ is not connected and $\mathcal{H}_S(X)$ is not simply connected. The maps $f$ and $g$ here are defined by the vertex coloring.

2.4. **Gromov’s cylinder construction.** In this section we review a construction, due to Gromov, which turns a simplicial complex $K$ into a foldable cubical complex $\mathcal{G}(K)$ having non-positive curvature (see [Gro87, §3.4.A] for the original source, or [DJ91, §4c], [Pan91, §4], and references therein, for expository accounts).

The construction uses induction on dimension and pullback simultaneously, following this scheme. For each dimension $n \geq 1$ we will first define $\mathcal{G}(\Delta^n)$ and a map $g : \mathcal{G}(\Delta^n) \rightarrow \Delta^n$, then for any foldable $n$–dimensional simplicial complex $K$, with a folding $f : K \rightarrow \Delta^n$, we will define $\mathcal{G}(K)$ via the pullback square (compare §2.3)
Finally, for a general $K$ (not necessarily foldable), we will define $\mathcal{G}(K) = \mathcal{G}(B(K))$ (recall that the barycentric subdivision is always foldable by Lemma 2.4). Note that in any case the construction equips $\mathcal{G}(K)$ with a natural map to $\triangle^n$.

For $n = 1$ we set $\mathcal{G}(\triangle^1) = \triangle^1$, and we define $g : \mathcal{G}(\triangle^1) \to \triangle^1$ to be just the identity. By the pullback construction this defines $\mathcal{G}(K)$ and a map $g : \mathcal{G}(K) \to \triangle^1$ for all simplicial graphs $K$. Concretely, when $K$ is a simplicial graph, then $\mathcal{G}(K) = K$ if $K$ is bipartite, and $\mathcal{G}(K) = B(K)$ otherwise; the folding to $\triangle^1$ is induced by the bipartition.

Let us now assume by induction that for any simplicial complex $K$ of dimension at most $n - 1$ the space $\mathcal{G}(K)$ is defined, and is endowed with a map to $\triangle^{n-1}$. In order to define $\mathcal{G}(\triangle^n)$, consider a reflection of $\partial \triangle^n$, and induce a reflection on $\mathcal{G}(\partial \triangle^n)$. Let $U, V$ be the two closed half-spaces exchanged by the reflection, and define

$$\mathcal{G}(\triangle^n) = \mathcal{G}(\partial \triangle^n) \times [-1, 1] / \sim$$

where $(u, t) \sim (u', t')$ if and only if $|t| = |t'| = 1$ and $u = u' \in U$. Notice that taking a further quotient which identifies also points on $V$, one would get a map $\mathcal{G}(\triangle^n) \to \mathcal{G}(\partial \triangle^n) \times S^1$, and we can think of $\mathcal{G}(\triangle^n)$ as being obtained from $\mathcal{G}(\partial \triangle^n) \times S^1$ by cutting a slit in it along a half–fiber (see Figure 5).

By induction, $\mathcal{G}(\partial \triangle^n)$ is well–defined, and it comes with a map $\mathcal{G}(\partial \triangle^n) \to \triangle^{n-1}$. Notice that the boundary of $\mathcal{G}(\triangle^n)$ consists of two copies of $V$, glued along a subspace identifiable with $U \cap V$. In other words, $\partial \mathcal{G}(\triangle^n)$ can be naturally identified with $\mathcal{G}(\partial \triangle^n)$, hence $\partial \mathcal{G}(\triangle^n)$ comes with a map to $\partial \triangle^n$. This map can be extended
to a map $\mathcal{G}(\Delta^n) \to \Delta^n$ as follows: take a regular neighborhood $N \cong \partial \mathcal{G}(\Delta^n) \times [0,1]$ inside $\mathcal{G}(\Delta^n)$, and identify $\Delta^n$ with the cone over $\partial \Delta^n$. Then extend the map over $N$ along the cone direction, and collapse the complement of $N$ to the cone point. This completes the construction of $\mathcal{G}(\Delta^n)$ and a map $g : \mathcal{G}(\Delta^n) \to \Delta^n$. Arguing as above (i.e. with the template from §2.3), this also defines $\mathcal{G}(K)$ for any simplicial complex $K$.

**Proposition 2.7.** If $K$ is a simplicial complex, then $\mathcal{G}(K)$ is a foldable cubical complex of non–positive curvature. Moreover if $K$ is homogeneous (respectively without boundary, locally compact, or compact), then $\mathcal{G}(K)$ is homogeneous (respectively without boundary, locally finite, or compact) too.

**Proof.** First we show that $\mathcal{G}(K)$ admits the structure of a cubical complex, starting with the case $K = \Delta^n$. This is clear for $\mathcal{G}(\Delta^1) = [0,1] = \square^1$. Then, arguing by induction, $\mathcal{G}(\Delta^n)$ inherits a cubical structure from the one of $\mathcal{G}(\partial \Delta^n) \times [-1,1]$. Here we think of $[-1,1]$ as being given the standard cubical structures as a union of two unit intervals, and we give $\mathcal{G}(\partial \Delta^n) \times [-1,1]$ the standard cubical structure coming from the fact that $\square^{n-1} \times \square^1 = \square^n$. Since $\mathcal{G}(K)$ is in general defined via the pullback construction (see §2.3), it inherits a natural cubical structure from $\mathcal{G}(\Delta^n)$.

We now prove that the cubical complex $\mathcal{G}(K)$ has the desired properties. Foldability is proven in [CD95, Lemma 7.5]. Non–positive curvature is proven in [DJ91, Proposition 4c.2(3)]. For the other properties we argue as follows. Note that for each $n$ the hyperbolizing cell $\mathcal{G}(\Delta^n)$ is homogeneous, has a single boundary component, and satisfies $\partial \mathcal{G}(\Delta^n) = g^{-1}(\partial \Delta^n)$. So, if $K$ is homogeneous then $\mathcal{G}(K)$ is homogeneous, and if $K$ is without boundary, the same holds for $\mathcal{G}(K)$. It is proved in [DJ91, Lemma 1e.1 and §4c] that Gromov’s construction preserves the local structure (e.g. links). This implies that if $K$ is locally finite, then so is $\mathcal{G}(K)$. In particular, by (3) in Lemma 2.5, if $K$ is compact, then so is $\mathcal{G}(K)$, because $\mathcal{G}(\Delta^n)$ is compact.

We have defined Gromov’s construction for simplicial complexes. Given any cell complex $X$ we can first take its barycentric subdivision $\mathcal{B}(X)$ (which is a simplicial complex), and then apply Gromov’s construction to it.

**Corollary 2.8.** If $X$ is any cell complex, then $\mathcal{G}(\mathcal{B}(X))$ is a foldable cubical complex of non–positive curvature. Moreover if $X$ is homogeneous (respectively without boundary, locally compact, or compact), then $\mathcal{G}(X)$ is homogeneous (respectively without boundary, locally compact, or compact) too.

**Proof.** We know $K = \mathcal{B}(X)$ is a (foldable) simplicial complex (by Lemma 2.4), homeomorphic to $X$. Then the statements follow from Proposition 2.7. □

Gromov’s construction is known to satisfy even more properties, namely conditions (1)-(6) in [CD95] and (1), (2'), (3)-(5) in [DJ91]. Some of these are versions of conditions (1)–(4) from §2.3, while others deal with preservation of stable tangent bundles and rational Pontryagin classes, when they are defined. This is needed in the applications of the hyperbolization procedure to construct examples of closed aspherical manifolds with various prescribed features or pathologies (as in [DJ91; Out20]).
3. Strict hyperbolization

The hyperbolization procedure introduced by Charney and Davis in [CD95] is defined for cubical complexes, and fits in the framework outlined in §2.3, in the sense that it is determined by the choice of a hyperbolizing cell. Differently from Gromov’s cylinder construction (described in §2.4), this procedure is not defined by induction. Rather, for each dimension \( n > 0 \) a hyperbolizing cell is defined independently, and defines a hyperbolization procedure for \( n \)-dimensional cubical complexes.

While the original construction is a bit more general than the version we use here, we find it convenient to impose some mild restrictions on the cubical complex in order to simplify the presentation. From now on assume \( X \) is admissible, i.e. it satisfies the following conditions (see § 2 for definitions):

1. cubical;
2. homogeneous, without boundary;
3. foldable;
4. non-positively curved;
5. locally compact.

This setting, consistent with that of [Xie04], is more general than the one in [BSa99], as we do not require gallery–connectedness. In particular, we allow \( X \) to be a pseudomanifold. On the other hand, the first two conditions are a bit more restrictive than the corresponding ones in [CD95], while the other ones are the same. More precisely, if \( X \) is foldable, then necessarily cubes of \( X \) are embedded. In [CD95] they allow two cubes to meet in a proper union of faces; note that such faces have to be disjoint in each cube, because non–positive curvature guarantees that links of vertices are simplicial. In particular, up to performing cubical subdivision, one can always assume that \( X \) is cubical. Finally we remark that at this stage we are only assuming local finiteness instead of compactness of \( X \). While in our main theorems (Theorems 1.1 and 1.2) we assume that the complex is compact (in order to get a hyperbolic group), most of the geometric and combinatorial arguments do not need that, and in §5.8 we actually need to consider a certain hyperbolization of \( \mathbb{R}^n \).

The main contribution of this section is to define some subspaces of the space that results from strict hyperbolization on an admissible cubical complex \( X \). We call such subspaces mirrors, and prove that their lifts to the universal cover are convex and separating (see Proposition 3.14 and Proposition 3.37 respectively). Along the way, we also study a combinatorial decomposition of the universal cover (see §3.3 and §3.5) that will be the starting point for the construction of the dual cubical complex in §4.

3.1. The hyperbolizing cell. The hyperbolizing cell used in this hyperbolization procedure is a certain hyperbolic manifold with boundary and corners, obtained by cutting a closed hyperbolic manifold along a suitable collection of pairwise orthogonal totally geodesic codimension–1 submanifolds. While the existence of such an object is clear in dimension 2 (see Figure 6), the construction in higher dimension requires some arithmetic methods involving quadratic forms (see §3.1.1 below for more details). Specifically, the construction relies on the choice of a suitable congruence subgroup \( \Gamma \) of an arithmetic lattice in \( \text{SO}_0(n, 1) \), so we will denote the hyperbolizing cell by \( \Box^\Gamma \). Here and in the following we denote by \( B_n \) the group of Euclidean isometries of the standard cube \( \Box^n \). Also recall from Remark 2.6 that a
Lemma 3.1 (Corollary 6.2 in [CD95]). For every \( n \geq 2 \) there exists a compact, connected, orientable hyperbolic \( n \)-manifold with corners \( \square^n \), an isometric action of \( B_n \) on \( \square^n \), and a \( B_n \)-equivariant and face-preserving map \( g : \square^k \to \square^n \), such that the following hold.

1. The poset of faces of \( \square^n \) is \( B_n \)-equivariantly isomorphic to that of \( \square^n \).
2. Each face of \( \square^n \) is totally geodesic.
3. The faces of \( \square_{k}^{n} \) intersect orthogonally.
4. Each 0-dimensional face is a single point.
5. The map \( g : \square_{k}^{n} \to \square^n \) and its restriction to each face have degree one.

We call \( \square^n \) the hyperbolizing cube, and \( g \) the Charney–Davis map. We denote by \( \Gamma_{\square^n} = \pi_1(\square^n) \) the fundamental group of the hyperbolizing cube.

Remark 3.2. In this hyperbolization procedure, a \( k \)-face of \( \square^n \) is guaranteed to be connected when \( k = 0, n \), but may be disconnected otherwise (see Remark 2.6, and the Remark after Corollary 6.2 in [CD95]). Nevertheless, by abuse of notation, we will denote by \( \square^k_{\Gamma} = g^{-1}(\square^k) \) the \( k \)-face of \( \square^n \), even when \( 0 < k < n \). Notice that \( \square^k_{\Gamma} \) is a priori different from the \( k \)-dimensional hyperbolizing cube, i.e. the hyperbolizing cell that one obtains by hyperbolizing a \( k \)-dimensional cube with a hyperbolizing lattice \( \Lambda \subseteq SO_0(k,1) \) for \( 0 < k < n \). Namely, \( \square^k_{\Lambda} \) is always connected by construction. Finally, with respect to (5) in Lemma 3.1, when \( \square^k_{\Gamma} \) is disconnected, there is a preferred component of \( \square^k_{\Gamma} \) on which \( g \) has degree one, while it has degree zero on the other components (see [CD95, Lemma 5.9(b)] and §3.1.1 for details).

3.1.1. The construction of \( \square^n_{\Gamma} \). To construct the hyperbolizing cube \( \square^n_{\Gamma} \), Charney and Davis consider the hyperboloid model for \( \mathbb{H}^n \) inside Minkowski space \( \mathbb{R}^{n,1} \), i.e. the space \( \mathbb{R}^{n,1} \) equipped with a quadratic form of signature \((n,1)\). The isometry group of \( \mathbb{R}^{n,1} \) is naturally identified with the indefinite orthogonal group \( O(n,1) \), and its connected component \( SO_0(n,1) \) is naturally identified with the group of orientation preserving isometries of \( \mathbb{H}^n \). Then they show that \( SO_0(n,1) \) contains an arithmetic lattice \( \Gamma \) which enjoys some key properties, from which the properties of \( \square^n_{\Gamma} \) in Lemma 3.1 follow. In particular, \( \Gamma \) is a cocompact torsion-free lattice of \( SO_0(n,1) \), whose normalizer in \( O(n,1) \) contains all the permutations of the coordinates \( x_1, \ldots, x_n \), and all the reflections in the coordinate hyperplanes \( H_i = \{ (x_1, \ldots, x_n, x_{n+1}) \in \mathbb{R}^{n,1} \mid x_i = 0 \} \) for \( i = 1, \ldots, n \). Note that these generate a group of isometries isomorphic to \( B_n \). We will refer to the lattice constructed in [CD95, §6] as the hyperbolizing lattice.

If \( \Gamma \) is such a lattice, then it acts freely, properly discontinuously, and cocompactly by orientation–preserving isometries on \( \mathbb{H}^n \). We can consider the closed connected oriented hyperbolic \( n \)-manifold \( M_{\Gamma} = \mathbb{H}^n / \Gamma \). The hyperplanes \( H_i \) descend to codimension–1 submanifolds \( M_i = H_i / \text{Stab}_{\Gamma}(H_i) \) which are closed, oriented, totally geodesic and pairwise orthogonal (see Figure 6). Then the hyperbolizing cell \( \square^n_{\Gamma} \) is defined to be the metric completion of the space \( M_{\Gamma} \setminus \bigcup_{i=1}^{n} M_i \), with respect to the length metric induced on the complement of \( \bigcup_{i=1}^{n} M_i \). This is the manifold with boundary and corners obtained by cutting \( M_{\Gamma} \) open along the submanifolds \( M_1, \ldots, M_n \) (see [CD95, §5]). In particular, the map \( g : \square^n_{\Gamma} \to \square^n \) is
induced by the collapse map $g_0 : M_\Gamma \to (S^1)^n$ obtained by applying the Pontryagin-Thom construction to $M_\Gamma$ with respect to each of the codimension–1 submanifolds $M_1, \ldots, M_n$.

Remark 3.3. It is implicit in [CD95] that a hyperbolizing lattice $\Gamma$ contains infinitely many other hyperbolizing lattices as proper subgroups. They still enjoy the properties which are relevant for the construction, and provide corresponding hyperbolizing cubes. As observed by Ontaneda in [Ont17, Lemma 2.1], this can be used to produce hyperbolizing cubes for which the normal injectivity radius of the faces is arbitrarily large.

3.2. The hyperbolized complex. Following the template of §2.3, to define the strict hyperbolization procedure of [CD95] we proceed as follows. For each dimension $n > 0$, we choose the hyperbolizing cell to be the hyperbolizing cube $(\square^n_\Gamma, g)$ defined in §3.1. Then for any foldable cubical complex $X$ of dimension $n$, we define the hyperbolized complex to be the space $X_\Gamma$ obtained as the fiber product of the folding map $f : X \to \square^n$ and the Charney-Davis map $g : \square^n_\Gamma \to \square^n$, i.e. by the pullback square in Figure 7.

![Diagram](image.png)

**Figure 6.** A hyperbolizing cube

**Figure 7.** The hyperbolized complex $X_\Gamma$ as a fibered product.

Remark 3.4. By (5) in Lemma 3.1 we know that $g$ is surjective. So, Lemma 2.5 allows us to think of $X_\Gamma$ as being obtained by replacing every $n$–cube of $X$ by a hyperbolizing cube $\square^n_\Gamma$, in the following sense (see Figure 8). If $C$ is a top–dimensional cube of $X$, then its preimage $g_X^{-1}(C)$ in $X_\Gamma$ is homeomorphic to $\square^n_\Gamma$ (see (1) in Lemma 2.5). Then one can endow $X_\Gamma$ with a length metric by gluing
Figure 8. Strict hyperbolization of a square complex.

together these local metrics. In particular, \( f_\Gamma : X_\Gamma \to \square^n_\Gamma \) induces an isometry \( g_X^{-1}(C) \to \square^n_k \) for each top-dimensional cube \( C \subseteq X \). For a concrete example, if \( X \) is (a suitable cubical subdivision of) the standard cubical structure on the \( n \)-torus, then \( X_\Gamma \) is a closed hyperbolic manifold (see [Bel07, Lemma 3.2] for details).

Indeed, the piecewise hyperbolic metric obtained by gluing the hyperbolizing cubes together has no singularity and is in fact globally smooth and hyperbolic.

We collect here some of the main properties of this construction which are relevant for our work.

**Proposition 3.5** (Proposition 7.1 in [CD95]). For every \( n \geq 2 \) and every \( n \)-dimensional foldable cubical complex \( X \), the space \( X_\Gamma \) carries the structure of an \( n \)-dimensional piecewise hyperbolic cell complex, and is endowed with a map \( g_X : X_\Gamma \to X \), such that the following hold.

1. If \( C \subseteq X \) is a \( k \)-cube, then \( g_X^{-1}(C) \subseteq X_\Gamma \) is isometric to a \( k \)-face of \( \square^n_\Gamma \), and \( \text{lk}(g_X^{-1}(C), X_\Gamma) \) is a piecewise spherical cell complex, isomorphic to \( \text{lk}(C, X) \).
2. If \( Z \subseteq X \) is locally convex subcomplex of \( X \), then \( g_X^{-1}(Z) \) is a locally convex subspace of \( X_\Gamma \).
3. If \( X \) is locally \( \text{CAT}(0) \), then \( X_\Gamma \) is locally \( \text{CAT}(-1) \).
4. If \( X \) is compact and locally \( \text{CAT}(0) \), then \( \Gamma_X = \pi_1(X_\Gamma) \) is a Gromov hyperbolic group.

**Remark 3.6.** The statement says in particular that if \( C \) is a top-dimensional cube of \( X \) then \( g_X^{-1}(C) \) is isometric to \( \square^n_\Gamma \) (compare Remark 3.4). On the other hand, if \( C \) is a \( k \)-cube with \( k < n \), then \( g_X^{-1}(C) \) is isometric to \( \square^k_\Gamma = g^{-1}(\square^k) \), i.e. the hyperbolization of a lower dimensional cell, as introduced in Remark 3.2.

If \( Z \subseteq X \) is a \( k \)-dimensional subcomplex, the subspace \( g_X^{-1}(Z) \) can be identified with the fibered product of the maps \( f_{_Z} : Z \to \square^k \) and \( g^k : \square^n_\Gamma \to \square^k \), respectively obtained by restricting the folding map \( f : X \to \square^n \) to \( Z \) and the Charney–Davis map \( g : \square^n_\Gamma \to \square^k \) (see Figure 9). Loosely speaking, hyperbolization trickles down to the lower dimensional skeletons of the complex \( X \).
Remark 3.7. In this construction the choice of $X$ and $\Gamma$ are essentially independent. In particular for any fixed cubical complex $X$ one can consider deeper hyperbolizations by taking deeper hyperbolizing lattices (see Remark 3.3). While the combinatorial geometry of the hyperbolized complex, controlled by $X$, remains unchanged under different choices of the hyperbolizing lattice, its hyperbolic geometry can be quantitatively improved by an appropriate choice of the hyperbolizing lattice, as observed by Ontaneda in [Ont17, Lemma 2.1].

Remark 3.8 (Finding codimension-1 subspaces). The original approaches to cubulating a group $G$ relied on producing sufficiently many codimension one subgroups inside $G$ (see [Sag95; Sag97; HW14; BW12a]).

Since the copies of $\Box^k_\Gamma$ in the hyperbolized complex $X_\Gamma$ are obtained from an arithmetic hyperbolic manifold, they contain a large supply of compact totally geodesic codimension one submanifolds. It is tempting to try and use these to produce codimension one subgroups in the hyperbolized group $\Gamma_X = \pi_1(X_\Gamma)$. The difficulty with this approach is due to lack of control on the angles at which these totally geodesic codimension one hypersurfaces intersect the boundary of $\Box^k_\Gamma$. This makes it unclear how to extend the proposed subspace past the boundary. One could take a geodesic extension, but it would not be clear what the global behaviour of the subspace would be (see left of Figure 10). Or one could take a geodesic reflection, but that would give rise to a kink angle (see right of Figure 10). Given that $\Box^k_\Gamma$ has fixed finite diameter, kink angles too far from right angles might prevent the subspace from even being quasiconvex.

You can try to control the kink angle, for instance by requiring the codimension one submanifold to be orthogonal to all faces of $\Box^k_\Gamma$. In this case, the extension would be a locally convex subspace of $X_\Gamma$. Examples of orthogonal subspaces can be obtained by noting that the symmetry group of the cube $B_n$ acts on $\Box^k_\Gamma$ (see Lemma 3.1). Each reflection of $B_n$ has some fixed point set, which meets the boundary orthogonally and is totally geodesic.

However, one can only find finitely many such subspaces, both in the orthogonal case and in the case of kink angles bounded away from 0 (see [Sha91] and [Fis+21, §5]). This would make it quite delicate to ensure that one can find enough such subspaces to apply the standard criteria for properness of the induced cubulation (such as those in [BW12a; HW14]). To address these issues, we turn to a different type of subspaces, which we call mirrors. These are defined in §3.4 using the foldability of $X$, and enjoy properties reminiscent of those of hyperplanes in a CAT(0) cube complex. For the sake of clarity, the collection of mirrors also fails to provide a proper action of $\Gamma_X = \pi_1(X_\Gamma)$ on a CAT(0) cubical complex in the usual
way. Nevertheless, in §4 we will be able to use mirrors to construct an action of $\Gamma_X$ on a CAT(0) cubical complex for which the cube stabilizers are manageable, and are in a certain sense already detected by the action of $\Gamma_X$ by deck transformations on the universal cover $\hat{X}_\Gamma$ (see § 5.1). The reader interested in these remarks should also compare this discussion with that in Remark 4.2 below.

3.3. Tiling, folding, and developing the universal cover. Recall that we are assuming $\hat{X}$ is an admissible complex, as defined at the beginning of §3. It follows from Proposition 3.5 (see also Lemma 2.5) that the hyperbolized complex $X_\Gamma$ admits a decomposition into hyperbolized cubes, analogous to the decomposition of $X$ into cubes. In this section we show how to obtain an analogous decomposition of the universal cover $\hat{X}_\Gamma$ of $X_\Gamma$ into pieces which are isometric to the universal cover $\hat{□}_n^\Gamma$ of the hyperbolizing cube. Let us denote by $\pi : \hat{X}_\Gamma \to X_\Gamma$ and $\pi_□ : \hat{□}_n^\Gamma \to □_n^\Gamma$ the universal covering projections.

We start by realizing the space $\hat{□}_n^\Gamma$ as a convex subset of $\mathbb{H}^n$. Let us consider once again the coordinate hyperplanes $H_i = \{(x_1, \ldots, x_n, x_{n+1}) \in \mathbb{R}^{n+1} \mid x_i = 0\}$ for $i = 1, \ldots, n$ (introduced in §3.1.1). An open $\Gamma$–cell is a connected component of the complement in $\mathbb{H}^n$ of the collection of $\Gamma$–orbits of the hyperplanes $H_i$. A $\Gamma$–cell is the closure of an open $\Gamma$–cell. Notice that all $\Gamma$–cells are convex, isometric to each other, and that $\Gamma$ permutes them transitively. It follows from the construction of $\hat{□}_n^\Gamma$ in §3.1.1 that the universal cover $\hat{□}_n^\Gamma$ of $\hat{□}_n^\Gamma$ can be isometrically identified with any $\Gamma$–cell (see Figure 11).

While it might be tempting to think that $\hat{X}_\Gamma$ is obtained via some simple fibered product construction involving $\hat{X}$ and $\hat{□}_n^\Gamma$, that is not the case.

Remark 3.9 (What $\hat{X}_\Gamma$ is not). Note that $\hat{X}_\Gamma \neq (\hat{X})_\Gamma$, i.e. the universal cover of the hyperbolization of $X$ is not the hyperbolization of the universal cover of $X$. Indeed, $(\hat{X})_\Gamma$ is not simply connected, because it retracts to $\hat{□}_n^\Gamma$ by (2) in Lemma 2.5. Analogously, $\hat{X}_\Gamma$ is not the fiber product of $\hat{X}$ and $\hat{□}_n^\Gamma$ either. Indeed, note that the faces of $\hat{□}_n^\Gamma$ (i.e. the preimages of faces of $\hat{□}_n^\Gamma$ via the map $g \circ \pi_\Gamma$) are disconnected (see Figure 11). This prevents the fiber product of $\hat{X}$ and $\hat{□}_n^\Gamma$ from being simply connected, as observed in Remark 2.6.
In order to address this, and get a working understanding of $\widetilde{X}_\Gamma$, we consider the intermediate space $X'_\Gamma$ obtained as a fibered product of $X$ and $\widetilde{\Box}_\Gamma$ along the maps $f : X \to \Box^n$ and $g \circ \pi_{\Box} : \widetilde{\Box}_\Gamma \to \Box^n$ (see Figure 12). By the universal property of pullbacks we have an induced map $\pi' : X'_\Gamma \to X_\Gamma$.

**Lemma 3.10.** The map $\pi' : X'_\Gamma \to X_\Gamma$ is a covering map.

**Proof.** By the composition law for pullbacks the space $X'_\Gamma$ is actually the same as the pullback of $f_\Gamma : X_\Gamma \to \Box^n$ and $\pi_{\Box} : \widetilde{\Box}_\Gamma \to \Box^n$. In particular, the map $\pi'$ is the...
pullback of the universal covering projection $\pi_\square$ along the map $f$, hence is itself a covering map.

In particular, $X'_\Gamma$ can be endowed with a length metric that makes $\pi'$ a local isometry (see [BH99, Proposition I.3.25]), and the universal cover $\tilde{X}_\Gamma$ can be realized as the universal cover of this space $X'_\Gamma$, even in a metric sense. Let $\pi'' : \tilde{X}_\Gamma \to X'_\Gamma$ denote the universal covering projection.

![Figure 13. Tiles in $\tilde{X}_\Gamma, X_\Gamma$ and $X$.](image)

We define a tile of $X_\Gamma$ to be a subspace of the form $g^{-1}_X(C)$, for $C$ a top-dimensional cube of $X$. Recall from Remark 3.4 that each tile of $X_\Gamma$ is isometric to $\square^n_\Gamma$. In complete analogy, we define a tile in $X'_\Gamma$ and in $\tilde{X}_\Gamma$ to be a connected component of the lift of a tile from $X_\Gamma$ via the covering maps $\pi'$ and $\pi = \pi' \circ \pi''$ respectively. We refer to this decomposition into tiles as the tiling of each of these spaces (see Figure 13). Note that, since the complex $X$ is assumed to be admissible, each point of $X$ is either contained in the interior of a tile, or in the intersection of at least two tiles. Moreover the folding map $f$ of $X$ induces an analogous map on $X_\Gamma$ and its covering spaces, as established in the next lemma.

**Lemma 3.11.** The map $\tilde{f}_\Gamma = f'_\Gamma \circ \pi'' : \tilde{X}_\Gamma \to X'_\Gamma \to \square^n_\Gamma$ restricts to an isometry between each tile of $\tilde{X}_\Gamma$ and $\square^n_\Gamma$.

**Proof.** Recall that $X'_\Gamma$ is defined via a pullback construction, in the sense of §2.3. Therefore, by (1) in Lemma 2.5, the map $f'_\Gamma : X'_\Gamma \to \square^n_\Gamma$ restricts to a homeomorphism between each tile of $X'_\Gamma$ and $\square^n_\Gamma$. Since the metric on $X'_\Gamma$ is lifted from $X_\Gamma$ via $\pi'$, and $f_\Gamma$ restricts to an isometry between each tile of $X_\Gamma$ and $\square^n_\Gamma$ (see Remark 3.4), the map $f'_\Gamma$ actually gives an isometry between a tile of $X'_\Gamma$ and $\square^n_\Gamma$. Since the tiles of $X'_\Gamma$ are simply connected, they lift isometrically to tiles of $\tilde{X}_\Gamma$ via $\pi''$. In particular, $\pi''$ maps a tile of $\tilde{X}_\Gamma$ isometrically onto a tile of $X'_\Gamma$. Therefore the map $\tilde{f}_\Gamma = f'_\Gamma \circ \pi''$ maps a tile of $\tilde{X}_\Gamma$ isometrically onto $\square^n_\Gamma$, just by composition. \quare

The map $\tilde{f}_\Gamma$ from Lemma 3.11 will be called the folding map of $\tilde{X}_\Gamma$. The composition of the folding map $\tilde{f}_\Gamma$ with any isometric embedding $\varphi : \square^n_\Gamma \to C$ onto a $\Gamma$–cell $C \subseteq \mathbb{H}^n$ will be called a developing map for $\tilde{X}_\Gamma$.

**Remark 3.12.** The restriction of a developing map to a tile is an isometric embedding of a tile into $\mathbb{H}^n$ as a $\Gamma$–cell. Moreover if $T_1, T_2$ are two tiles of $\tilde{X}_\Gamma$ meeting
along a codimension–1 subspace \( Z \), and \( \phi_1 : T_1 \to C_1 \subseteq \mathbb{H}^n \) is an isometric embedding onto a \( \Gamma \)–cell that maps \( Z \) into some hyperplane \( H \), then post-composing \( \phi_1 \) with the reflection across \( H \) provides an isometric embedding \( \phi_2 \) of \( T_2 \) as a \( \Gamma \)–cell \( C_2 \) adjacent to \( C_1 \) along \( H \). The two embeddings can be glued together to give an isometric embedding of \( T_1 \cup T_2 \) onto the union of two \( \Gamma \)–cells \( C_1 \cup C_2 \) adjacent along \( H \). This can be “analytically continued” by sequentially extending over adjacent tiles, to obtain a globally defined map \( \tilde{X}_\Gamma \to \mathbb{H}^n \). However, this does not result in a global isometric embedding \( \tilde{X}_\Gamma \to \mathbb{H}^n \) in general. This is due to the fact that links in \( X \) can be very large, which gives rise to overlaps and singularities.

### 3.4. Mirrors: convexity

In this section we exploit foldability to define a collection of convex subcomplexes of \( X \), and induce corresponding subspaces in \( X_\Gamma \) and \( \tilde{X}_\Gamma \). Let \( Y \) be a foldable cubical complex of dimension \( n \) (in the following we will consider \( Y = X \) and \( Y = \tilde{X} \) depending on the situation). If \( f : Y \to \square^n \) is a fixed folding and \( F \subseteq \square^n \) is a codimension-1 face, then we define a mirror in \( Y \) to be a connected component of \( f^{-1}(F) \).

**Proposition 3.13.** Let \( Y \) be an admissible cubical complex. Then each mirror is a locally convex and geodesically complete subcomplex of \( Y \). In particular, if \( Y \) is CAT(0), then each mirror is convex.

**Proof.** For the first statement see [Xie04, Proposition 2.3] (and references therein such as [BSa99, Lemma 3.2(4)]). In the CAT(0) case, local convexity implies global convexity (see for instance [BW12b, Theorem 1.6, 1.10], or [RC16, Theorem 1.1]). \( \square \)

![Figure 14. Mirrors in \( \tilde{X}_\Gamma \), \( X_\Gamma \) and \( X \).](image)

We now define a mirror in \( \tilde{X}_\Gamma \) to be a connected components of \( \tilde{f}^{-1}(F) \), where \( F \) is a codimension-1 face of \( \square^n \) and \( \tilde{f} \) is the map given by the composition \( \tilde{f} = f \circ g_X \circ \pi : \tilde{X}_\Gamma \to X_\Gamma \to X \to \square^n \) (see Figure 15). Equivalently, we could define it as a connected component of the full preimage of a mirror of \( X \) via the map \( g_{\tilde{X}} = g_X \circ \pi \), but we find it convenient to use this definition. We will say that \( M \) folds to \( F \), and we will denote by \( M \) the collection of all mirrors in \( \tilde{X}_\Gamma \). Mirrors in \( X_\Gamma \) are defined in the analogous way using the map \( f \circ g_X \).

**Proposition 3.14.** Let \( X \) be an admissible cubical complex. Then each mirror of \( \tilde{X}_\Gamma \) is a closed connected convex subspace of \( \tilde{X}_\Gamma \).
Figure 15. The hyperbolized complex $X_\Gamma$ and the maps used to define mirrors.

**Proof.** Let $M$ be a mirror of $\tilde{X}_\Gamma$, and let $F \subseteq \Box^n$ be the codimension-1 face to which it folds. By definition $M$ is connected and closed. To prove convexity we argue as follows. Let $Z = g_X(\pi(M)) \subseteq X$, and notice that $Z$ is a mirror of $X$ that folds to $F$. By Proposition 3.13 we know that $Z$ is locally convex in $X$. By (2) in Proposition 3.5, we also know that $g_X^{-1}(Z)$ is locally convex in $X_\Gamma$, and therefore $M \subseteq \tilde{X}_\Gamma$ is locally convex too. By (3) in Proposition 3.5 we also know that $X_\Gamma$ is locally $\text{CAT}(-1)$. In particular $M$ is a closed and locally convex subspace in the $\text{CAT}(0)$ space $\tilde{X}_\Gamma$. Therefore it is convex (again by [BW12b, Theorem 1.6, 1.10], or [RC16, Theorem 1.1]).

3.5. **Stratification of** $\tilde{X}_\Gamma$. In this section we use the collection $\mathcal{M}$ of mirrors, introduced in §3.4, to define a stratification of $\tilde{X}_\Gamma$. The open $k$–stratum $\Sigma^k$ of $\tilde{X}_\Gamma$ is the subspace consisting of points that fold into the interior of a $k$–face of $\Box^n$ via the map $f = f \circ g_X \circ \pi : \tilde{X}_\Gamma \to X_\Gamma \to X \to \Box^n$, or equivalently to the interior of a $k$–cube of $X$ via the map $g_X = \pi \circ g_X : \tilde{X}_\Gamma \to X_\Gamma \to X$ (see Figure 15). Notice that $\Sigma^k$ is a locally closed subspace. An open $k$–cell is a connected component of $\Sigma^k$. A $k$–cell is the closure of an open $k$–cell. We say that a cell $\sigma$ folds to the face $F = f(\sigma) \subseteq \Box^n$ and to the cube $C = g_X(\sigma) \subseteq X$. The integer $k$ will be referred to as the dimension of a $k$–cell. An $(n-k)$–cell is a proper subset of the intersection of $k$ mirrors. In particular 0–cells are points, and $n$–cells are tiles (as defined in §3.3). We call 0–cells vertices, and 1–cells edges of the stratification.

**Remark 3.15 (Cellular structure).** We explicitly observe that this choice of strata does not define a stratified space structure on $\tilde{X}_\Gamma$ in the sense of [BH99, Definition II.12.1]. Moreover, the decomposition of $\tilde{X}_\Gamma$ into cells does not turn it into a genuine cell complex, as defined in §2.1. Indeed, while an open $k$–cell is homeomorphic to an open disk of dimension $k$, a $k$–cell is not homeomorphic to a closed disk of dimension $k$ as soon as $k \geq 2$. Its boundary in $\tilde{X}_\Gamma$ consists of an infinite union of lower–dimensional cells, so it is neither connected nor compact. For instance, an $n$–cell (i.e. a tile) is isometric to a $\Gamma$–cell (see Figure 11).

Nevertheless, we can still recover a lot of the classical behavior and tools, by observing that cells are convex and that the link of cells and points can be defined in analogy to the classical case (see §2.1.3). We gather here preliminary results
Lemma 3.16. Let $\sigma \subseteq \widetilde{X}_\Gamma$ be a cell. Then $\sigma$ is convex.

Proof. Let $C = \tilde{g}_X(\sigma) \subseteq X$ be the cube of $X$ to which $\sigma$ folds. By (2) in Proposition 3.5, we know $g_X^{-1}(C)$ is locally convex in $X_\Gamma$. Since $\sigma$ is by definition a connected component of $\pi^{-1}(g_X^{-1}(C))$ and $\pi$ is a local isometry, we can conclude that it is a locally convex subspace of $\widetilde{X}_\Gamma$. Arguing similarly to previous proofs of convexity, we can conclude that $\sigma$ is convex, because it is closed and locally convex in the CAT(0) space $\widetilde{X}_\Gamma$ (see (3) in Proposition 3.5 and [BW12b, Theorem 1.6, 1.10], or [RC16, Theorem 1.1]).

We now proceed to the study of links. Consider the universal covering map $\pi : \widetilde{X}_\Gamma \to X_\Gamma$. By Proposition 3.5, $X_\Gamma$ is a piecewise hyperbolic cell complex, so the link of points and cells in $X_\Gamma$ is well–defined (see §2.1.3). Since $\pi$ is a local isometry, we can just identify the link of points and cells in $\widetilde{X}_\Gamma$ with the links of the corresponding points and cells in $X_\Gamma$.

Lemma 3.17. Let $\sigma \subseteq \widetilde{X}_\Gamma$ be a cell.

1. Let $C = \tilde{g}_X(\sigma) \subseteq X$ be the cube to which it folds. Then $\tilde{g}_X$ induces an isomorphism between $\text{lk}(\sigma, \widetilde{X}_\Gamma)$ and $\text{lk}(C, X)$.

2. Let $\sigma$ be contained in another cell $\tau$. Let $F = \tilde{f}(\sigma), E = \tilde{f}(\tau) \subseteq \boxtimes^n$ be the faces to which they fold. Then $\tilde{f}$ induces an isomorphism between $\text{lk}(\sigma, \tau)$ and $\text{lk}(F, E)$.

3. Let $\sigma$ be a $k$–cell. Then $\text{lk}(\sigma, \widetilde{X}_\Gamma)$ is a piecewise spherical simplicial complex with vertices given by the $(k+1)$–cells containing $\sigma$, and in which $m+1$ vertices span an $m$–simplex if and only if the corresponding $(k+1)$–cells are contained in a $(k+m+1)$–cell.

Proof. The map $\tilde{g}_X : \widetilde{X}_\Gamma \to X$ is the composition of the map $\pi : \widetilde{X}_\Gamma \to X_\Gamma$, which preserves links because it is a covering map, and the map $g_X : X_\Gamma \to X$, which preserves links thanks to (1) in Proposition 3.5. This proves (1).

To prove (2) we argue similarly. The map $\tilde{f} : \widetilde{X}_\Gamma \to \boxtimes^n$ is the composition of the map $\tilde{g}_X : \widetilde{X}_\Gamma \to X$, which preserves links by (1), and the folding map $f : X \to \boxtimes^n$. By definition of folding, $f$ is a combinatorial isomorphism on each cube of $X$. If $B$ is the cube to which $\tau$ folds, the folding induces an isomorphism between $\text{lk}(C, B)$ and $\text{lk}(F, E)$.

Finally, (3) follows from (1), the fact that $\tilde{g}_X : \widetilde{X}_\Gamma \to X$ maps cells of $\widetilde{X}_\Gamma$ to cubes of $X$ preserving inclusion relations, and the fact that the link of a cell in a cubical complex carries a piecewise spherical simplicial structure as described in the statement.

Lemma 3.18. Let $\sigma_1, \sigma_2 \subseteq \widetilde{X}_\Gamma$ be two cells. Then either $\sigma_1 \cap \sigma_2$ is empty or it is a cell.

Proof. Let $\sigma_1 \cap \sigma_2$ be non empty. If it contains either a single vertex, or a single edge, then we are done. So let us assume that it contains at least two edges. Also note that, since cells are convex by Lemma 3.16, the intersection $\sigma_1 \cap \sigma_2$ is convex. Therefore, if there are several edges then they cannot all be disjoint.
Let \( v \in \sigma_1 \cap \sigma_2 \) be a vertex, and let \( e,e' \) be two edges of \( \sigma_1 \cap \sigma_2 \) meeting at \( v \). Note that by Lemma 3.17 links in \( \tilde{X}_\Gamma \) are isomorphic to the corresponding links in \( X \). In particular, \( e \) and \( e' \) are edges of the cell \( \sigma_1 \) meeting at a vertex, so there is a 2–face \( \tau_1 \subseteq \sigma_1 \) containing both \( e \) and \( e' \). Analogously, we get a 2–face \( \tau_2 \subseteq \sigma_2 \) with the same property. Since these links are simplicial, necessarily we have \( \tau_1 = \tau_2 \), otherwise we would see a bigon in the link of \( v \). In particular, \( \tau_1 = \tau_2 \subseteq \sigma_1 \cap \sigma_2 \). This shows that any two edges of \( \sigma_1 \cap \sigma_2 \) meeting at \( v \) are adjacent in \( \text{lk}(v,\sigma_1 \cap \sigma_2) \). By Lemma 3.17 we know that \( \text{lk}(v,\tilde{X}_\Gamma) \cong \text{lk}(\tilde{g}_X(v),X) \), and this is a flag simplicial complex because \( X \) is non–positively curved (see Lemma 2.2). The same holds for \( \text{lk}(v,\sigma_1 \cap \sigma_2) \) because \( \sigma_1 \cap \sigma_2 \) is convex. In particular, all the edges of \( \sigma_1 \cap \sigma_2 \) that contain \( v \) are actually contained in a single cell of \( \sigma_1 \cap \sigma_2 \), which we denote \( \sigma_v \). Now, if \( v \), \( w \) are adjacent vertices of \( \sigma_1 \cap \sigma_2 \), then by uniqueness we have \( \sigma_v = \sigma_w \). Finally, by connectedness of \( \sigma_1 \cap \sigma_2 \), it follows that all vertices of \( \sigma_1 \cap \sigma_2 \) are contained in a single cell.

**Lemma 3.19.** Let \( \{ \sigma_j \mid j \in J \} \) be a collection of cells of \( \tilde{X}_\Gamma \). Then the following statements hold.

1. If \( \sigma = \bigcap_{j \in J} \sigma_j \) is not empty, then \( \sigma \) is the unique cell of maximal dimension contained in \( \sigma_j \) for all \( j \in J \).
2. If \( \bigcup_{j \in J} \sigma_j \) is contained in a single cell, then there exists a unique cell \( \sigma \) of minimal dimension containing \( \sigma_j \) for all \( j \in J \).

We refer to the cell in (1) (respectively (2)) of Lemma 3.19 as the lower cell (respectively upper cell) of the collection \( \{ \sigma_j \mid j \in J \} \).

**Proof.** Since \( X \) is finite–dimensional and locally compact, if \( J \) is infinite, then \( \sigma = \bigcap_{j \in J} \sigma_j \) is empty. So let as assume that \( J \) is finite. By Lemma 3.18 the intersection of finitely many cells is either empty or made of a single cell. This proves (1). To prove (2), assume by contradiction that there are two different cells of minimal dimension \( \sigma, \sigma' \) containing each \( \sigma_j \). Then \( \sigma \cap \sigma' \) is a proper union of cells, against Lemma 3.18.

**Lemma 3.20.** Let \( \tau \subseteq \tilde{X}_\Gamma \) be a cell. Let \( \sigma_1, \sigma_2 \subseteq \tau \) be cells of lower dimension, and let \( F_1,F_2 \subseteq \square^n \) be the faces to which they fold. If \( F_1 = F_2 \), then \( \sigma_1, \sigma_2 \) are either disjoint or equal.

**Proof.** Assume that \( \sigma_1, \sigma_2 \) are not disjoint, and let \( v \in \sigma_1 \cap \sigma_2 \) be a vertex. Let \( E \subseteq \square^n \) be the face to which \( \tau \) folds. We have that \( \tilde{f}(\sigma_1) = F_1 = F_2 = \tilde{f}(\sigma_2) \), and by (2) in Lemma 3.17 (with \( \sigma = v \)) the map \( \tilde{f} \) induces an isomorphism between \( \text{lk}(v,\tau) \) and \( \text{lk}(\tilde{f}(v),E) \). Therefore, \( \sigma_1 = \sigma_2 \).

**Lemma 3.21.** Let \( M \) be a mirror and let \( \tau \) be a \( k \)-cell of \( \tilde{X}_\Gamma \) not entirely contained in \( M \). If \( M \cap \tau \neq \emptyset \), then \( M \cap \tau \) is a \( (k-1) \)-cell.

**Proof.** First we show that \( M \cap \tau \) is a union of \( (k-1) \)-cells. Then we show that the union actually consists of a single cell.

Let \( \tilde{f} = f \circ g_X : \tilde{X}_\Gamma \to X_\Gamma \to X \to \square^n \) be the map that folds \( \tilde{X}_\Gamma \) to \( \square^n \), and let \( F \subseteq \square^n \) be the codimension–1 face to which \( M \) folds (i.e. \( F = \tilde{f}(M) \)). Similarly,
let $E \subseteq \square^n$ be the $k$–face to which $\tau$ folds (i.e. $E = \tilde{f}(\tau)$). Since $\tau \not\subseteq M$ we have $E \not\subseteq F$, and therefore $E \cap F$ is a $(k–1)$–face of $\square^n$. Let $p \in M \cap \tau$ be an arbitrary point. By Lemma 3.20, among the $(k–1)$–cells of $\tau$ that contain $p$, there is exactly one that folds to $E \cap F$; denote it by $\sigma_p$. Clearly $\sigma_p \subseteq \tau$. Moreover, since $\sigma_p$ and $M$ are non–disjoint, both fold into $F$, and $M$ is a mirror, we also have $\sigma_p \subseteq M$.

Therefore we have $M \cap \tau = \bigcup_{p \in M \cap \tau} \sigma_p$, i.e. $M \cap \tau$ is a union of $(k–1)$–cells that fold to $E \cap F$.

To see that $M \cap \tau$ actually consist of only one cell, assume by contradiction that $M \cap \tau$ contains two distinct $(k–1)$–cells $\sigma_1, \sigma_2$. Let $p_i \in \sigma_i$ and let $\gamma = [p_1, p_2]$ be the unique geodesic between them. Since $M$ and $\tau$ are both convex (by Proposition 3.14 and Lemma 3.16 respectively), we have that $\gamma \subset \tau \cap M$. Hence we find a path of cells in the boundary of $\tau$ that all fold to $E \cap F$. But this is absurd because different boundary cells of $\tau$ folding to the same face of $\square^n$ are necessarily disjoint, again by Lemma 3.20. \hfill $\square$

### 3.6. Graph of spaces decomposition for $\tilde{X}_\Gamma$

Our goal in §3.7 will be to prove that mirrors in $\tilde{X}_\Gamma$ enjoy a strong separation property. Our strategy will be to exploit a certain graph of spaces decomposition for $\tilde{X}_\Gamma$ (in the sense of [SW79]), which we introduce in this section, using the foldability of $X$ (see [BSa99; Xie04] for analogous constructions).

Recall from §3.4 that $\tilde{X}_\Gamma$ is equipped with a collection $M$ of closed convex subspaces called mirrors. For each $i = 1, \ldots, n$, let $M_i$ be the collection of mirrors of $\tilde{X}_\Gamma$ that fold to one of the two parallel $i$th faces of $\square^n = [0, 1]^n$, i.e. $\{ x_i = 0 \}$ and $\{ x_i = 1 \}$. Notice that by construction any two elements of $M_i$ are disjoint, and even have disjoint $\varepsilon$–neighborhoods for $\varepsilon$ sufficiently small (because $\Gamma$ is cocompact).

Let $\mathcal{C}_i$ be the collection of connected components of $\tilde{X}_\Gamma \setminus \cup_{M \in M_i} M$. For each mirror $M \in M_i$ and for each component $C \in \mathcal{C}_i$, consider the following equidistant space, obtained by pushing the mirror $M$ into the component $C$ (see Figure 16).

$$E_{M,C}^\varepsilon = \{ x \in C \mid d(x, M) = \varepsilon \}.$$

Notice that while we know $M$ is convex by Proposition 3.14, it is not clear whether $C$ is convex. A priori, $C$ could meet $M$ on more than one side, i.e. the closure of $C$ could contain a piece of $M$ in its interior. We will see this is not the case by considering a suitable graph of spaces decomposition of $\tilde{X}_\Gamma$. Our first step is to show that $E_{M,C}^\varepsilon$ is simply connected; in the process, we actually show it is a CAT($k$)–space for some $k \in (-1, 0)$. The idea for this can be summarized as follows: inside each tile, $E_{M,C}^\varepsilon$ looks like an equidistant hypersurface from a hyperplane in $\mathbb{H}^n$, and this is a non–positively curved hypersurface in $\mathbb{H}^n$ (see Figure 17). Then contribution from different tiles come together in a way that does not introduce any positive curvature along mirrors. We start from a preliminary lemma from classical hyperbolic geometry. Recall that a hyperplane in $\mathbb{H}^n$ is a totally geodesic copy of $\mathbb{H}^{n-1}$.

**Lemma 3.22.** Let $V \subseteq \mathbb{H}^n$ be a hyperplane, and let $\pi_V : \mathbb{H}^n \to V$ be the nearest point projection to $V$. Let $\varepsilon > 0$ and $S_V^\varepsilon = \{ x \in \mathbb{H}^n \mid d(x, V) = \varepsilon \}$. Then the following hold.

1. $S_V^\varepsilon$ is a smooth $(n–1)$–dimensional submanifold of $\mathbb{H}^n$.
2. For each $p \in S_V^\varepsilon$, the geodesic $[p, \pi_V(p)]$ is orthogonal to $V$ and $S_V^\varepsilon$.
3. For every other hyperplane $W$, if $V \cap W \neq \emptyset$, then $S_V^\varepsilon \cap W \neq \emptyset$. 


Figure 16. Some examples of equidistant spaces in dimension 2 and 3. Left: an equidistant space relative to a mirror $M$ in the vicinity of the intersection with two other mirrors. Here $\dim X = 3$, and all mirrors are locally Euclidean. Right: three equidistant spaces relative to the same mirror $M$ but three different complementary components, in the vicinity of the intersection with another mirror. Here $\dim X = 2$, and the mirrors branch, i.e. are not locally Euclidean.

Figure 17. Equidistant surface from a hyperplane.

(4) For every hyperplane $W$, $W$ is orthogonal to $S^\gamma_V$ if and only if $W$ is orthogonal to $V$.

(5) $\pi_V : S^\gamma_V \rightarrow V$ is a $\cosh^2(\varepsilon)$–conformal diffeomorphism.

(6) The induced metric on $S^\gamma_V$ has constant sectional curvature $-\frac{1}{\cosh^2(\varepsilon)}$.

Proof. The first five statements can be proved by explicit computations in the upper half–space model of $\mathbb{H}^n$, normalizing so that $V$ is a vertical hyperplane (see Figure 17). The computation for dimension $n = 3$ is carried out in detail in [Fen89, IV.5, page 58], and readily generalizes to higher dimensions. Finally, (6) follows from (5) and the general formula for the behavior of the sectional curvatures under rescaling.
For the next lemma, recall from §3.3 that tiles are closed by definition, and that a developing map is an isometric embedding of a tile into $\mathbb{H}^n$ as a $\Gamma$–cell.

**Lemma 3.23.** Let $M \in \mathcal{M}_i$ and $C \in \mathcal{C}_i$. Then for $\varepsilon > 0$ small enough the following hold.

1. For every mirror $N \in \mathcal{M}$, if $E^\varepsilon_{M,C} \cap N \neq \emptyset$ then $M \cap N \neq \emptyset$ and $C \cap N \neq \emptyset$.
2. For every tile $\tau$, if $E^\varepsilon_{M,C} \cap \tau \neq \emptyset$ then $M \cap \tau \neq \emptyset$ and $C \cap \tau \neq \emptyset$.
3. For every tile $\tau$ such that $E^\varepsilon_{M,C} \cap \tau \neq \emptyset$, and any developing map $\varphi : \tau \to \mathbb{H}^n$, $\varphi$ induces an isometry between $E^\varepsilon_{M,C} \cap \tau$ and $S_V \cap \varphi(\tau)$, where $V$ is the hyperplane containing $\varphi(M \cap \tau)$.
4. For every mirror $N \in \mathcal{M}$, if $E^\varepsilon_{M,C} \cap N \neq \emptyset$ then $E^\varepsilon_{M,C}$ is orthogonal to $N$.

**Proof.** To prove (1) note that if $E^\varepsilon_{M,C} \cap N \neq \emptyset$, then in particular $C \cap N \neq \emptyset$. Since $\Gamma$ is cocompact, there is a uniform lower bound $D > 0$ on the distance between disjoint mirrors. But $E^\varepsilon_{M,C} \cap N \neq \emptyset$ means that $N$ comes $\varepsilon$ close to $M$. By choosing $\varepsilon < D$ we can force $N$ to actually intersect $M$. The proof of (2) is analogous to that of (1). Suppose $E^\varepsilon_{M,C} \cap \tau \neq \emptyset$. Then clearly $C \cap \tau \neq \emptyset$. Moreover, a point in $E^\varepsilon_{M,C} \cap \tau$ witnesses that $d(M, \tau) < \varepsilon$, and by choosing $\varepsilon$ small enough we can ensure that this forces an intersection, again by cocompactness of $\Gamma$.

Now we consider (3). Suppose that $E^\varepsilon_{M,C} \cap \tau \neq \emptyset$. Then by (2) we know that $M \cap \tau \neq \emptyset$ and $C \cap \tau \neq \emptyset$. In particular $M$ appears as an $(n-1)$–cell in the boundary of $\tau$ thanks to Lemma 3.21. If we pick a developing map $\varphi$ for $\tau$, then $\varphi(\tau)$ is a $\Gamma$–cell, and $\varphi(M)$ is some hyperplane $V$ on its boundary (see Lemma 3.11 and Remark 3.12). Then the statement follows from the fact that $\varphi$ is an isometric embedding of $\tau$ into $\mathbb{H}^n$.

Finally, to prove (4), suppose that $E^\varepsilon_{M,C} \cap N \neq \emptyset$. Then by (1) we know that $N \cap M \neq \emptyset$. In particular by construction $N$ is orthogonal to $M$. Then the statement follows from (3), together with (4) in Lemma 3.22.

Next, our goal is to prove that equidistant spaces are negatively curved. In order to do this, we will study the geometry of links of points in $\tilde{X}_\Gamma$, along various subspaces (we refer the reader to §3.5 for definitions). Recall that the link of a point in $\tilde{X}_\Gamma$ is identified to the link of its projection to $X_\Gamma$.

**Remark 3.24.** All the subspaces of $\tilde{X}_\Gamma$ considered here (such as a mirror $M$, and the induced space $E^\varepsilon_{M,C}$) carry a natural locally finite cellular structure induced by that of $\tilde{X}_\Gamma$. Even if they are not genuine cell complexes (as in Remark 3.15), their projections to $X_\Gamma$ are, and links can be defined in analogy to the classical case.

**Lemma 3.25.** Let $M \in \mathcal{M}_i$, $C \in \mathcal{C}_i$, $p \in E^\varepsilon_{M,C}$. Then for $\varepsilon > 0$ small enough the following holds. Let $\tau_1, \ldots, \tau_m$ be the collection of tiles containing $p$, and let $T = \tau_1 \cup \cdots \cup \tau_m$. Then the following hold.

1. $\text{lk}(\pi_M(p), M \cap T)$ is CAT(1).
2. $\pi_M : E^\varepsilon_{M,C} \to M$ induces an isometry $\lambda_p : \text{lk}(p, E^\varepsilon_{M,C}) \to \text{lk}(\pi_M(p), M \cap T)$.
3. $\text{lk}(p, E^\varepsilon_{M,C})$ is CAT(1).

**Proof.** Of course, (3) follows from (1) and (2). For convenience, let us denote $L = \text{lk}(\pi_M(p), \tilde{X}_\Gamma)$, $L_T = \text{lk}(\pi_M(p), T)$, and $L_{M \cap T} = \text{lk}(\pi_M(p), M \cap T)$. We
have $L_{MCT} \subseteq L_T \subseteq L$. Equip $L_{MCT}$ and $L_T$ with the induced length metric. Let $\overrightarrow{p} \in L_T$ be the direction at $\pi_M(p)$ pointing to $p$ (see Figure 18).

**Figure 18.** Links of points along various subspaces in the proof of Lemma 3.25. Here $p$ is contained in four tiles and sits on the intersection of two mirrors $M', M''$. The vertical projection is the nearest point projection $\pi_M : E_{M,C}^p \to M$.

We start by proving (1). Since $\hat{X}_\Gamma$ is negatively curved, $L$ is CAT(1). In particular, balls of radius at most $\pi/2$ are $\pi$-convex and CAT(1). Since $\hat{X}_\Gamma$ is piecewise hyperbolic, $L$ is piecewise spherical. Moreover, by (4) in Lemma 3.23 we know that all the mirrors containing $p$ intersect $M$ orthogonally. Therefore, $L$ has a natural structure of all-right spherical complex in which $\overrightarrow{p}$ is a vertex (possibly up to subdivision if $\pi_M(p)$ is not a vertex). In particular, we have natural identifications $L_T = B(\overrightarrow{p}, \frac{\pi}{2})$ and $L_{MCT} = \partial B(\overrightarrow{p}, \frac{\pi}{2})$.

Let $C_1(Y)$ denote the spherical cone over a space $Y$, and denote the cone point by 0. Since $L$ is an all-right spherical complex, we have a natural isometry

$$\varphi : C_1(\partial B(\overrightarrow{p}, \frac{\pi}{2})) \to B(\overrightarrow{p}, \frac{\pi}{2})$$

defined as follows: $\varphi(0) = \overrightarrow{p}$, and for each $\overrightarrow{q} \in \partial B(\overrightarrow{p}, \frac{\pi}{2})$ and $0 < t \leq \frac{\pi}{2}$ let $\varphi(t, \overrightarrow{q})$ be the point at distance $t$ from $\overrightarrow{p}$ along the geodesic $[\overrightarrow{p}, \overrightarrow{q}]$. As a result, $C_1(L_{MCT}) = C_1(\partial B(\overrightarrow{p}, \frac{\pi}{2}))$ is CAT(1). By Berestovskii’s Theorem (see [BH99, p. II.3.14]) we conclude that $L_{MCT}$ is CAT(1) as desired.

To prove (2) we argue as follows. By (3) in Lemma 3.23 and (5) in Lemma 3.22 we know that within each tile $\tau_k$ the projection $\pi_M$ is a conformal diffeomorphism, so it induces an isometry $\lambda_p^{\tau_k} : \text{lk}(p, E_{M,C}^p \cap \tau_k) \to L_k = \text{lk}(\pi_M(p), M \cap \tau_k)$. This
is enough in the case $m = 1$, i.e. when $p$ is contained in a single tile. When $m \ge 2$, by gluing together the maps $\lambda_p^\pi$, we obtain a map $\lambda_p : \text{lk} (p, E^\pi_{M,C}) \to L_1 \cup \cdots \cup L_m = L_{M \cap T}$. Notice that shooting geodesic rays from $\pi_M (p)$ into $T$ along directions in $L_T$ provides an isometry

$$\psi : \text{lk} (\overrightarrow{p}, L_T) \to \text{lk} (p, E^\pi_{M,C})$$

Combining this with the natural isometry

$$r : \text{lk} (\overrightarrow{p}, B \left(\overrightarrow{p}, \frac{\pi}{2} \right)) \to \partial B \left(\overrightarrow{p}, \frac{\pi}{2} \right)$$

and using the aforementioned identifications, we obtain the desired isometry

$$\text{lk} (p, E^\pi_{M,C}) \xrightarrow{\psi^{-1}} \text{lk} (\overrightarrow{p}, L_T) = \text{lk} \left(\overrightarrow{p}, B \left(\overrightarrow{p}, \frac{\pi}{2} \right) \right) \xrightarrow{r} \partial B \left(\overrightarrow{p}, \frac{\pi}{2} \right) = L_{M \cap T}.$$

Remark 3.26. Note that, in the notation of Lemma 3.25, $L_{M \cap T} = \text{lk} (\pi_M (p), M \cap T)$ is a closed subspace of $\text{lk} (\pi_M (p), M)$ which is possibly proper. Indeed, $\pi_M (p)$ might live on a lower dimensional cell, where $M$ might branch off away from $T$, as in Figure 19. However, all the branches make an angle of at least $\pi$ with each other, because $M$ is convex.

**Figure 19.** A mirror $M$ branching away from $T$, the union of tiles containing $p$ (other mirrors not displayed).

**Lemma 3.27.** Let $M \in \mathcal{M}_i$ and $C \in \mathcal{C}_i$. Then for $\varepsilon > 0$ small enough there is $k \in (-1,0)$ such that the following hold.

1. The metric induced on $E^\pi_{M,C}$ is locally CAT($k$).
(2) The nearest point projection $\pi_M : E^\varepsilon_{M,C} \to M$ maps non–constant local geodesics to non–constant local geodesics.

(3) The metric induced on $E^\varepsilon_{M,C}$ is CAT($k$).

Proof. To prove (1) we argue as follows. By (3) in Lemma 3.23, we know that, away from the intersection with mirrors, $E^\varepsilon_{M,C}$ is locally isometric (via a developing map) to an equidistant hypersurface in $\mathbb{H}^n$. Such a hypersurface is a manifold of negative curvature $k \in (-1, 0)$ by (5) in Lemma 3.22. By Remark 3.24, $E^\varepsilon_{M,C}$ is essentially a cell complex, so by [BH99, Theorem II.5.2] $E^\varepsilon_{M,C}$ is locally CAT($k$) if and only if the link of every vertex is a CAT(1) space. This condition is verified by (3) in Lemma 3.25.

Now we consider (2). By (3) in Lemma 3.23 and (5) in Lemma 3.22, we know that in the interior of each tile $\pi_M$ is a conformal diffeomorphism with constant conformal factor. Therefore it sends a local geodesic on $E^\varepsilon_{M,C}$ to a piecewise local geodesic on $M$, possibly broken at points where two or more tiles meet. To take care of those possibly singular point, we invoke (2) in Lemma 3.25, which guarantees that $\pi_M$ induces an isometric embedding of links also at those points. Indeed, if $p \in E^\varepsilon_{M,C}$ is such a break point, and $c$ is a geodesic on $E^\varepsilon_{M,C}$ through $p$, then the incoming and outgoing directions are at distance $D \geq \pi$ in $\text{lk} (p, E^\varepsilon_{M,C})$. Let $c' = \pi_M (c)$. Then $c'$ is a piecewise geodesic in $M$ through $\pi_M (p)$. With the notations of Lemma 3.25, the distance in $\text{lk} (\pi_M (p), M \cap T)$ between the incoming and outgoing directions is the same $D \geq \pi$. The distance in the full $\text{lk} (\pi_M (p), M)$ is not smaller, as $\text{lk} (\pi_M (p), M)$ does not contain geodesic loops shorter than $2\pi$ by convexity. So, $c'$ is a local geodesic in $M$ at $\pi_M (p)$. Moreover if $c$ is non-constant then $c'$ is non–constant because $\pi_M$ is locally injective.

To conclude, we prove (3). By (1) we know that $E^\varepsilon_{M,C}$ is locally CAT($k$), so we only need to prove that it is also simply connected. By contradiction, let $\gamma \in \pi_1 (E^\varepsilon_{M,C})$ be a non–trivial homotopy class. Since $E^\varepsilon_{M,C}$ is complete and non–positively curved, $\gamma$ is represented by a unique non–constant local geodesic $c_\gamma$. By (2) $\pi_M (c_\gamma)$ is a non–constant local geodesic on $M$. Since $M$ is complete and non–positively curved, $\pi_M (c_\gamma)$ is not nullhomotopic, which contradicts the fact that $M$ is contractible.

Remark 3.28. Note that if for a mirror $M$ and a tile $\tau$ the intersection $M \cap \tau$ was lower–dimensional, then the edge space would develop to an equidistant hypersurface from a lower–dimensional totally geodesic subspace of $\mathbb{H}^n$, which has some positive curvature. So, Lemma 3.21 (establishing that if a mirror intersects a tile then the intersection is a codimension-1 cell) is a key tool to prove that edge spaces are non–positively curved.

Proposition 3.29. $\widehat{X}_T$ admits the structure of a graph of spaces, with underlying graph a connected tree.

Proof. We define a graph $T_i$ as follows. Vertices come in two different families, namely a vertex $v_M$ for each mirror $M \in M_i$ and a vertex $v_C$ for each component $C \in \mathcal{C}_i$. Then we place one edge $e_{M,C}$ between $v_M$ and $v_C$ whenever $M$ intersects the closure of $C$. Vertex and edge spaces are defined as follows: we associate $M$ to $v_M$, $C$ to $v_C$, and $E^\varepsilon_{M,C}$ to the edge $e_{M,C}$ between them.

The edge maps to the two types of vertices are respectively given by the nearest point projection $\pi_M : E^\varepsilon_{M,C} \to M$ and the natural inclusion $i : E^\varepsilon_{M,C} \hookrightarrow C$. The
resulting space is homeomorphic to $\tilde{X}_\Gamma$. Notice that edge spaces are contractible by (3) in Lemma 3.27, so the gluing maps are automatically injective on fundamental groups.

We are left to show that $T_i$ is a connected tree. Connectedness of $T_i$ follows directly from that of $X$. There is a natural map $r_i : \tilde{X}_\Gamma \to T_i$ obtained by collapsing all the vertex spaces to points and all the cylinders over edge spaces to edges. Notice that $r_i$ is a retraction and $\tilde{X}_\Gamma$ is contractible, which forces $T_i$ to be simply connected. □

Remark 3.30. In this graph of spaces decomposition all the spaces involved are non–positively curved, but the edge maps are not local isometries. Moreover, further pathological behavior can arise depending on the structure of the mirrors, as we now discuss. Note that the following phenomena already arise in the setting of cubical complexes, i.e. are not introduced by the hyperbolization procedure.

On one hand, if the mirror $M$ branches (i.e. has non–locally Euclidean points) in such a way that different branches meet the closure of different complementary components, then the nearest point projections $\pi_M : E_{M,C} \to M$ from the individual edge spaces fail to be surjective.

On the other hand, if the mirror $M$ is such that a complementary component $C$ wraps around $M$ and meets it on different sides, then the map $\pi_M : E_{M,C} \to M$ fails to be injective. This would be the case for a mirror that separates locally but not globally, e.g. one that is contained in the closure of a single complementary component. In this case the corresponding vertex would be a boundary vertex for the tree $T_i$. We will see in § 3.7 that this failure of injectivity does not occur in our setting.

Remark 3.31 (A graph of groups decomposition for $\Gamma_X$). Note that $\Gamma_X = \pi_1(X_\Gamma)$ acts on $\tilde{X}_\Gamma$ sending mirrors to mirrors and preserving the coloring, i.e. each family $\mathcal{M}_i$. In particular it preserves this graph of spaces decomposition, hence it acts on the underlying graph, which has been seen to be a tree. The action is without global fix points and without inversions. This realizes $\Gamma_X = \pi_1(X_\Gamma)$ as a graph of groups. It is worth noticing that combination theorems are available in the literature, which provide a way to construct a cubulation of a group expressed as a graph of cubulated groups, when certain conditions are met (see for instance [HW12; HW19; Wis21]). In our context, the vertex groups are given by the fundamental groups of the mirrors from $\mathcal{M}_i$ and the components from $\mathcal{C}_i$. While it is reasonable to expect that the former are cubulated (e.g. arguing by induction on dimension), it is not at all clear that the latter should be. The guiding idea for the rest of the paper is that nevertheless those components can be further decomposed into tiles. The fundamental group of a tile can be shown to be cubulated (see Lemma 5.12), and the results of Groves and Manning from [GM18] then provide a way to combine the cubulation from each tile into a global cubulation.

3.7. Mirrors: separation. In this section we will prove a strong separation property for mirrors in $\tilde{X}_\Gamma$. In order to obtain convexity of the mirrors, in the proof of Proposition 3.14 we have used the fundamental fact that in a CAT(0) space local convexity implies global convexity. The same local–to–global property fails for separation, as shown by the following example.

Example 3.32. Consider the square complex $Y$ in the center of Figure 20. Consider the subcomplex $Z$ consisting of the central thick (red) edge. The subspace $Z$
is locally separating in $Y$, in the sense that for any $z \in Z$ and any arbitrarily small neighborhood $U_z$ of $z$ in $Y$, $U_z \setminus Z$ is disconnected. However, $Z$ is not separating, i.e. $Y \setminus Z$ is connected. Notice that $Y$ is a CAT(0) and foldable cubical complex, but $Z$ is not a full connected component of the preimage of a codimension-1 face, i.e. not a mirror.

In this example both $Y$ and $Z$ have boundary, but it can be modified to obtain an example without boundary. We start by attaching eight more squares following the pattern in Figure 20, and extending $Z$ with two more edges. In the resulting complex, no edge meeting $Z$ is a boundary edge, so we can keep adding squares (and extending $Z$) to get an admissible complex which displays the same pathology as the original one.

![Figure 20. A locally separating but not separating subcomplex in a CAT(0) square complex.](image)

When $Y$ is a homogeneous cubical complex of dimension $n$, every $k$–cube $F$ of $Y$ is contained in some $n$-cell. When $Y$ has no boundary, $F$ is contained in at least two distinct $n$-cubes. This motivates the following definition. Let $M$ be a mirror of $Y$ and let $F$ be a $k$–cube of $M$. A framing for $F$ is a choice of two $n$–cubes $\{C_1, C_2\}$ of $Y$ such that $F \subseteq C_1 \cap C_2 \subseteq M$. We note explicitly that this definition is relative to the fixed mirror $M$. For the next proof, we will make use of some properties of hyperplanes in CAT(0) cubical complexes. We refer the reader to [Sag95, Theorem 4.10] or [HW08, Example 3.3.(3), Lemma 13.3] for details and proofs.

![Figure 21. A framing for a cube $F$ on a mirror $M$.](image)

**Lemma 3.33.** Let $Y$ be a CAT(0) admissible cubical complex. Then each mirror separates $Y$. More precisely, let $M \subseteq Y$ be a mirror, let $F \subseteq M$ be a $k$–cube, and let $\{C_1, C_2\}$ be a framing for $F$. Then $C_1, C_2$ are contained in the closure of two distinct connected components of $Y \setminus M$. 

Proof. Let $v$ be a vertex on $F$, let $e_i$ be the edge of $C_i$ with starting point $v$ and endpoint in $C_i \setminus M$ (see Figure 21). Note that this edge exists because $C_i$ is $n$–dimensional, while $M$ is $(n-1)$–dimensional and convex, so that $M \cap C_i$ is some $(n-1)$–dimensional face $E_i$ of $C_i$. Also note that by definition of framing, $C_1 \cap C_2 \subseteq M$, and therefore $e_1 \neq e_2$. Let $H_i$ be the hyperplane of $Y$ dual to $e_i$. In particular this means that $H_i$ meets $C_i$ in the midcube orthogonal to $e_i$. Since $Y$ is CAT(0), we get that $H_1 \neq H_2$, $H_1 \cap H_2 = \emptyset$, $H_k \cap M = \emptyset$, and $Y \setminus H_k$ consists of exactly two components, one containing $M$ and one not containing $M$.

The carrier of a hyperplane $H$ in a CAT(0) is isomorphic to $H \times [0,1]$. By definition of mirror, if $M$ contains an $(n-1)$–cube of $H \times \{0\}$ then actually $H \times \{0\} \subseteq M$. Since $M$ contains the $(n-1)$–cell $E_i = C_i \cap M$ of $C_i$, and $E_i \subseteq H_i \times \{0\}$ by construction, we can conclude that $M$ contains $H_i \times \{0\}$ for $i = 1, 2$. It follows that any path from $H_1$ to $H_2$ must intersect $M$. In particular, $M$ separates $Y$ in at least two components, one containing $H_1$ and one containing $H_2$. The closures of such components contain $C_1$ and $C_2$ respectively.

We want to extend this result to mirrors in $\widetilde{X}_\Gamma$. To do this, we introduce the following terminology, in analogy with the cubical case. Let $M$ be a mirror of $\widetilde{X}_\Gamma$, and let $\sigma$ be a $k$–cell of $M$. A framing for $\sigma$ is the choice of two distinct $n$–cells $\tau_1, \tau_2$ such that $\sigma \subseteq \tau_1 \cap \tau_2 \subseteq M$. We begin by obtaining a weak separation property.

Lemma 3.34. Let $M \in \mathcal{M}_i$, let $\sigma \subseteq M$ be a $k$–cell, and let $\{\tau_1, \tau_2\}$ be a framing for $\sigma$. Then there exist two different components $C_1, C_2 \in \mathcal{C}_i$ whose closure contain $\tau_1, \tau_2$ respectively.

Proof. The map $g_X : X_\Gamma \to X$ lifts to a map $\alpha : \widetilde{X}_\Gamma \to \widetilde{X}$ between the universal covers. Note that it sends mirrors to mirrors. In particular we obtain a mirror $\alpha(M)$ and a $k$–cube $\alpha(\sigma) \subseteq \alpha(M)$ with a framing $\{\alpha(\tau_1), \alpha(\tau_2)\}$. By Lemma 3.33 we can conclude that $\alpha(\tau_1)$ and $\alpha(\tau_2)$ are separated by $\alpha(M)$ in $\widetilde{X}$. This implies that $\alpha^{-1}(\alpha(\tau_1))$ and $\alpha^{-1}(\alpha(\tau_2))$ are separated in $\widetilde{X}_\Gamma$ by $\alpha^{-1}(\alpha(M))$, i.e. the full preimage of the mirror $\alpha(M)$ in $\widetilde{X}_\Gamma$. Note that $\alpha^{-1}(\alpha(M))$ consist of infinitely many mirrors from $\mathcal{M}_i$; indeed, recall from Lemma 3.21 that disjoint $(n-1)$–cells of a tile belong to different mirrors. A fortiori, $\tau_1$ and $\tau_2$ are separated by the entire collection $\mathcal{M}_i$. In particular, there exists two different components $C_1, C_2 \in \mathcal{C}_i$ whose closure contain $\tau_1, \tau_2$ respectively, as desired.

Remark 3.35. Observe that in the proof of Lemma 3.34, it is not clear whether the framing is separated by $M$ itself. While the entire collection of mirrors $\mathcal{M}_i$ disconnects $\widetilde{X}_\Gamma$ into a collection of complementary components, it is not a priori clear that any single mirror separates $\widetilde{X}_\Gamma$.

Recall from Proposition 3.29 that $\widetilde{X}_\Gamma$ admits the structure of a graph of spaces over a connected tree $T_i$, and that there is a natural retraction $r_i : \widetilde{X}_\Gamma \to T_i$ obtained by collapsing all the vertex spaces to points and all the cylinders over edge spaces to edges.

Lemma 3.36. The tree $T_i$ has no boundary.

Proof. It is enough to show that each vertex has at least two neighboring vertices. Vertices of $T_i$ are either associated to mirrors from $\mathcal{M}_i$ or to components from $\mathcal{C}_i$. We analyze the two different cases separately. Let $v_C$ be the vertex associated to
a component $C \in \mathcal{C}_i$. Then $v_C$ has infinitely many edges coming into it, because $C$ has infinitely many mirrors from $M_i$ on its boundary (this is already true for a single tile: by Lemma 3.21, disjoint $(n-1)$–cells in the boundary of a tile belong to different mirrors).

Now let $v_M$ be the vertex associated to a mirror $M \in M_i$. Let $\sigma \subseteq M$ be an $(n-1)$-cell on it, and pick a framing $\{\tau_1, \tau_2\}$. By Lemma 3.34, there exist two different components $C_1, C_2 \in \mathcal{C}_i$ whose closure contain $\tau_1, \tau_2$ respectively. The corresponding vertices $v_{C_1}, v_{C_2}$ in $T_i$ are both adjacent to the vertex $v_M$ corresponding to $M$, as desired. \hfill \Box

The next result is the analogue of Lemma 3.33 from the cubical case.

**Proposition 3.37.** Each $M \in M_i$ separates $\tilde{X}_\Gamma$. More precisely, let $M \in M_i$ be a mirror, let $\sigma \subseteq M$ be a $k$–cell, and let $\{\tau_1, \tau_2\}$ be a framing for $\sigma$. Then $\tau_1, \tau_2$ are contained in the closure two distinct connected components of $\tilde{X}_\Gamma \setminus M$.

**Proof.** For the first statement, consider the natural retraction $r_i : \tilde{X}_\Gamma \to T_i$. Note that for each mirror $M \in M_i$ there is a corresponding vertex $v_M \in T_i$, and $M = r_i^{-1}(v_M)$. By Lemma 3.36 we know that $T_i$ is a tree with no boundary, hence any of its vertices disconnects it. Therefore $M = r_i^{-1}(v_M)$ disconnects $\tilde{X}_\Gamma$.

For the second statement, we fix a $k$–cell $\sigma \subseteq M$ and a framing $\{\tau_1, \tau_2\}$. By Lemma 3.34 we get two components $C_1, C_2 \in \mathcal{C}_i$ containing $\tau_1, \tau_2$ in their closures. Note that these are complementary components of the entire collection of mirrors $M_i$, not complementary components of the mirror $M$. The corresponding vertices $v_{C_1}, v_{C_2}$ in $T_i$ are both adjacent to the vertex $v_M$ corresponding to $M$, and are separated by $v_M$ in $T_i$, since $T_i$ is a tree (see Proposition 3.29). Arguing as above via the natural retraction $r_i : \tilde{X}_\Gamma \to T_i$, we can conclude that $\tau_1, \tau_2$ are separated by $M$ in $\tilde{X}_\Gamma$.

We conclude this section with some remarks about the construction that we have described.

**Remark 3.38 (Foldability is key).** Foldability of $X$ has played the role of some sort of combinatorial completeness, as it guarantees that if a mirror $M$ intersects a tile $T$, then $M$ goes across $T$ along a top dimensional subcomplex of the boundary. This has provided both features of non–positive curvature (see Remark 3.28) and separation properties (as in the proof of Lemma 3.33). Example 3.32 shows that neither is available if foldability is not taken into account in the definition of mirrors (even on a foldable complex).

**Remark 3.39 (Complexes with boundary).** The construction from §3.6 can be generalized to cubical complexes that have enough good mirrors (i.e. mirrors that admit a cell which locally separates a framing), and keeping track only of such mirrors in the construction of the tree of spaces. For instance, one could drop the assumption that $X$ is without boundary, and ignore the mirrors that are entirely contained in the boundary. One still gets a decomposition as a graph of spaces over a tree without boundary. Indeed, vertices associated to good mirrors still have degree at least 2. One may worry about vertices associated to components. Even if there is a cube of $X$ with only one face $F$ contained in a good mirror, each of the components $C \in \mathcal{C}_i$ of $\tilde{X}_\Gamma$ arising from it still has infinitely many boundary cells corresponding to $F$. This guarantees that the vertices of the tree which are...
associated to components in $\mathcal{C}_i$ still have infinite degree. We will pursue this point of view in a following paper on cubulation of relatively hyperbolic groups obtained by relative hyperbolization procedures (see [CD95; Bel07]).

4. The dual cubical complex

We define a cubical complex associated to the stratification of $\tilde{X}_Γ$ introduced in §3.5, and prove that it is a CAT(0) cubical complex (see Theorem 4.29). Recall that $X$ is assumed to be an admissible cubical complex (as defined at the beginning of §3). Let $n = \dim(X)$ be its dimension. The dual cubical complex is denoted $\mathcal{C}(\tilde{X}_Γ)$ and defined as follows.

- Vertices are given by all the $k$-cells in $\tilde{X}_Γ$ for $k = 0, \ldots, n$.
- Two vertices corresponding to cells $σ$ and $τ$ are connected by an edge if and only if $|\dim(σ) − \dim(τ)| = 1$, and either $σ \subseteq τ$ or $τ \subseteq σ$.
- For $k > 1$, we attach one $k$-dimensional cube whenever we see its 1-skeleton.

The resulting cell complex $\mathcal{C}(\tilde{X}_Γ)$ is a cubical complex (see Figure 22). Moreover, we can label its 0-skeleton by integers $0 \leq k \leq n$: if $v$ is a vertex dual to a $k$-cell $σ$, then we define the height of $v$ to be $h(v) = \dim(σ) = k$.

In this section we study the combinatorial geometry of $\mathcal{C}(\tilde{X}_Γ)$, by analyzing cubes and links in §4.1, some notions of complexity for edge–paths in §4.2 and §4.3, and how to use them to prove that $\mathcal{C}(\tilde{X}_Γ)$ is simply connected in §4.4. Before starting, the following two remarks address the relation between $\mathcal{C}(\tilde{X}_Γ)$ and other natural combinatorial structures associated to $\tilde{X}_Γ$ and its collection of mirrors $M$.

Remark 4.1 (The associated graded poset). The set of cells in $\tilde{X}_Γ$ can be partially ordered by inclusion. The result is a graded poset, whose rank function is given by the dimension of the corresponding cell. The height we just defined is induced by this rank function. One could construct the order complex of such a poset, by taking a simplex for every chain. This would result in a simplicial complex, and is not what we are considering here.

Remark 4.2 (The associated wallspace). Since mirrors are separating subspaces (see Proposition 3.37), the collection of mirrors can be used to define a wallspace structure $(\tilde{X}_Γ, M)$ on $\tilde{X}_Γ$, and one could consider the dual CAT(0) cubical complex $\mathcal{C}(\tilde{X}_Γ, M)$ associated to this wallspace by Sageev’s construction. We refer the reader to [Sag95; HP98; HW14] for details about this construction, and we only review the main ingredients here. Given a mirror $M$, any partition of the complementary components into two classes is called a wall associated to $M$. An orientation of a wall is a choice of one of the two classes. A vertex of $\mathcal{C}(\tilde{X}_Γ, M)$ can then be described as a consistent choice of orientation for each mirror.

When $X$ and all mirrors are homeomorphic to manifolds, each mirror of $\tilde{X}_Γ$ has exactly two complementary components. In this quite restrictive case, an orientation of a wall is just a choice of one of the two complementary components. Therefore vertices of $\mathcal{C}(\tilde{X}_Γ, M)$ correspond to tiles (i.e. $n$–cells) in the stratification of $\tilde{X}_Γ$, and two vertices are connected by an edge when the corresponding tiles are adjacent along a mirror. In particular, $\mathcal{C}(\tilde{X}_Γ, M)$ is an $n$–dimensional cubical complex that can be subdivided to recover $\mathcal{C}(\tilde{X}_Γ)$. However, if there are mirrors which have more than two complementary components (such as in Figures 14 and
Figure 22. The dual cubical complex $\mathcal{C}(\tilde{X}_\Gamma)$ superimposed on the stratification of $\tilde{X}_\Gamma$. (In this picture the dimension is $n = 2$. Key: $\circ$, $\odot$, and $\bullet$ denote a vertex of height 0, 1, 2 respectively.)

then we find vertices in $\mathcal{C}(\tilde{X}_\Gamma, \mathcal{M})$ which do not correspond to tiles from the stratification of $\tilde{X}_\Gamma$ (they are not canonical vertices, in the terminology of [HW14]). As a result, the dimension of $\mathcal{C}(\tilde{X}_\Gamma, \mathcal{M})$ is usually higher than that of $X$, and it is more challenging to relate the actions of $\Gamma_X$ on $\tilde{X}_\Gamma$ and on $\mathcal{C}(\tilde{X}_\Gamma, \mathcal{M})$.

4.1. Cubes and links. In this section we explore basic facts about the cubical geometry of $\mathcal{C}(\tilde{X}_\Gamma)$. While this complex is not locally compact (see Remark 4.8), its dimension is the same as that of $X$ (see Lemma 4.5), and the links of vertices are flag complexes (see Proposition 4.10).

The first two lemmas show that squares and cubes in $\mathcal{C}(\tilde{X}_\Gamma)$ admit unique vertices of minimum and maximum height. Recall that the height of a vertex is the
dimension of its dual cell, and notice that, by definition of $C(\tilde{X}_\Gamma)$, if $u, v$ are adjacent vertices, then $|h(u) - h(v)| = 1$.

![Figure 23. A square in $C(\tilde{X}_\Gamma)$.](image)

**Lemma 4.3.** Let $S$ be a square of $C(\tilde{X}_\Gamma)$. Let $v_1, v_2, v_3, v_4$ be its vertices, with $v_2$ and $v_4$ non-adjacent in $S$. If $h(v_2) = h(v_4)$, then $|h(v_1) - h(v_3)| = 2$. In particular there is a unique vertex of maximal (respectively minimal) height, and the cell dual to it contains (respectively is contained in) each of the cells dual to the other vertices.

**Proof.** Let $h = h(v_2) = h(v_4)$ be the common value of the height of $v_2$ and $v_4$. Since $v_1$ is adjacent to $v_2$ and $v_4$, we have $h(v_1) = h \pm 1$, and similarly for $v_3$ (see Figure 23). In particular $|h(v_1) - h(v_3)|$ is either 0 or 2. By contradiction let us assume that $|h(v_1) - h(v_3)| = 0$, i.e. $h(v_1) = h(v_3) = h \pm 1$. Without loss of generality we can assume that $h(v_1) = h(v_3) = h + 1$. (The case $h(v_1) = h(v_3) = h - 1$ is completely analogous, via a dual argument). For $j = 1, 2, 3, 4$, let $\sigma_j$ be the cell of $\tilde{X}_\Gamma$ dual to the vertex $v_j$. Since $v_1$ is adjacent to $v_2$ and $v_4$, and has higher height, $\sigma_1$ contains $\sigma_2$ and $\sigma_4$; the same holds for $\sigma_3$. So $\sigma_1 \cap \sigma_3$ contains $\sigma_2 \cup \sigma_4$, contradicting Lemma 3.18.

To prove the final statement, let us assume without loss of generality that $v_1$ is the vertex of maximal height and $v_3$ is the one of minimal height, i.e. $h(v_1) - 1 = h = h(v_3) + 1$. Then we have that $\sigma_3 \subseteq \sigma_2, \sigma_4 \subseteq \sigma_1$.

In the next lemma we extend this result to higher dimensional cubes of $C(\tilde{X}_\Gamma)$. By an edge–path in $C(\tilde{X}_\Gamma)$ we will mean a continuous path which is entirely contained in the 1–skeleton (i.e. is a sequence of edges). If an edge–path $p$ goes through vertices $v_0, \ldots, v_s$ of $C(\tilde{X}_\Gamma)$, we will write $p = (v_0, \ldots, v_s)$; note that the sequence of vertices completely determines the sequence of edges, hence the path. We call $p$ an edge–loop if it is a closed loop, i.e. $v_0 = v_s$. For an edge–path $p = (v_0, \ldots, v_s)$ we define $\ell(p) = s$ to be the length of $p$, i.e. the number of edges in it. We also define the height of $p$ to be $h(p) = \max\{h(v_0), \ldots, h(v_s)\}$. Notice that along each edge of $p$ the height must increase or decrease exactly by 1.

**Lemma 4.4.** Let $Q$ be a cube of $C(\tilde{X}_\Gamma)$. Then the following hold.

1. There is a unique vertex $v \in Q$ of minimal height. The cell dual to it is contained in each of the cells dual to the vertices of $Q$.  

![Diagram of a cube in $C(\tilde{X}_\Gamma)$](image)
(2) There is a unique vertex \( w \in Q \) of maximal height. The cell dual to it contains each of the cells dual to the vertices of \( Q \).

Proof. We prove the first statement; the second is obtained by an analogous argument. Let \( k \) be the minimal height of vertices of \( Q \), and assume by contradiction that there is at least a pair of vertices of \( Q \) of height \( k \). Consider an edge–path \( p = (v_0, \ldots, v_s) \) in \( Q \) such that \( h(v_0) = h(v_s) = k \), \( v_0 \neq v_s \), and such that \( p \) is an edge–path of minimal height among all edge–paths in \( Q \) joining a pair of vertices of height \( k \). This is well–defined since the height of such a path can only be an integer between 0 and \( n \). Let \( h(p) = h \) be the height of \( p \).

![Diagram](image)

Figure 24. Lowering a vertex of maximal height on a path in \( C(\widetilde{X}_\Gamma) \).

Let \( v_j \) be a vertex of \( p \) of maximal height \( h(v_j) = h = h(p) \). Then \( h(v_{j \pm 1}) = h - 1 \) (notice that \( k \geq 0 \) and \( h \geq k + 1 \geq 1 \)). Since \( (v_{j-1}, v_j, v_{j+1}) \) is part of the cube \( Q \), it must be contained in a square, i.e. there exists a vertex \( u_j \) of \( Q \) (not necessarily on \( p \) ), such that \( \{v_{j-1}, v_j, v_{j+1}, u_j\} \) span a square in \( Q \). Lemma 4.3 implies that \( h(u_j) = h - 2 \). We can construct a new path in \( Q \) starting from the path \( p \) by lowering the vertex of maximal height, i.e. by replacing \( v_j \) with \( u_j \) (see Figure 24). We repeat the same operation on all vertices of height \( h \) along the path, and let \( p' \) be the resulting path in \( Q \). We have that \( h(p') = h - 1 < h = h(p) \), contradicting the minimality of the height of \( p \). This concludes the proof by contradiction, and proves the uniqueness of a vertex \( v \) of minimal height \( k \) in \( Q \).

We are left to show that the cell \( \sigma \) dual to \( v \) is contained in all the cells dual to the other vertices of \( Q \). By contradiction suppose there are vertices in \( Q \) whose dual cells do not contain \( \sigma \); call such vertices exceptional. Let \( p = (v_0, \ldots, v_s) \) be an edge–path in \( Q \) with \( v_0 = v \), \( v_s \) an exceptional vertex, and having minimal length among all edge–paths of \( Q \) between \( v \) and an exceptional vertex. We have \( h(v_{s-1}) = h(v_s) \pm 1 \). If \( h(v_{s-1}) = h(v_s) - 1 \), then the cell dual to \( v_{s-1} \) is contained in the cell dual to \( v_s \). By minimality of \( p \), we have that \( v_{s-1} \) is not exceptional, so the cell dual to \( v_{s-1} \) contains \( \sigma \), and hence \( v_s \) cannot be exceptional. Therefore \( h(v_{s-1}) = h(v_s) + 1 \) (as in Figure 25).

We keep walking backwards along \( p \) until we find a triple of vertices \( \{v_{j-1}, v_j, v_{j+1}\} \) such that \( h(v_j) = h(v_{j \pm 1}) + 1 \) (notice \( j \) is well defined and positive, since \( v \) is the unique vertex of minimal height in the whole \( Q \)). Arguing as before we complete to a square in \( Q \) with vertices \( \{v_{j-1}, v_j, v_{j+1}, u_j\} \); again by Lemma 4.3 we have \( h(u_j) = h(v_j) - 2 \) (see Figure 25). By minimality of \( p \), \( u_j \) must be non–exceptional,
and so we can change $p$ by replacing $v_j$ with $u_j$, without changing its length. Walking forward along $p$, we can keep changing the path without changing its length, until we are able to complete $\{v_{s-1}, v_s\}$ to a square $\{u_{s-2}, v_{s-1}, v_s, u_{s-1}\}$ in $Q$ with $u_{s-1}$ non-exceptional and with height $h(u_{s-1}) = h(v_{s-1}) - 2 = h(v_s) - 1$ (once again, see Figure 25). In particular, the cell dual to $u_{s-1}$ contains $\sigma$ and is contained in the cell dual to $v_s$, which contradicts the fact that the last vertex $v_s$ was chosen to be exceptional. □

As a consequence, we obtain the following statement.

**Lemma 4.5.** The complex $C(\tilde{X}_\Gamma)$ has dimension $\dim C(\tilde{X}_\Gamma) = \dim X = n$.

**Proof.** If $\tau$ is a tile of $\tilde{X}_\Gamma$ and $x$ one of its vertices, then the collection of cells containing $x$ and contained in $\tau$ provides a cube of dimension exactly $n$, so $\dim C(\tilde{X}_\Gamma) \geq n$, so we focus on the other inequality.

Let $Q$ be a cube of $C(\tilde{X}_\Gamma)$, and let $v_{\text{min}}$ be the vertex of minimal height in $Q$ (see Lemma 4.4). We claim that for each vertex $v \in Q$ we have

$$h(v) = h(v_{\text{min}}) + d_Q(v_{\text{min}}, v)$$

where $d_Q(v_{\text{min}}, v)$ is the distance in $Q$ of $v$ from $v_{\text{min}}$. Since the height can take values only between 0 and $n = \dim X$, this directly implies that

$$\dim Q = \max\{d_Q(v_{\text{min}}, v)\} = \max\{h(v) - h(v_{\text{min}})\} \leq n.$$ 

In order to prove the claim, pick a vertex $v \in Q$, and let $p = (v_0, \ldots, v_s)$ be an edge-path of minimal length $s = d_Q(v_{\text{min}}, v)$ in $Q$ from $v_0 = v_{\text{min}}$ to $v_s = v$. Since the height can at most increase by 1 along each edge of $p$, we have the inequality $h(v) \leq h(v_{\text{min}}) + d_Q(v_{\text{min}}, v)$. Assume by contradiction that the inequality is strict. Then the height is not monotonically increasing along $p$. Let $v_k$ be the first vertex of $p$ which is a local maximum for the height function. Arguing as above via Lemma 4.3, we look at the triple $v_{k-1}, v_k, v_{k+1}$, and complete it to a square with a fourth vertex $u_k$ such that $h(u_k) = h(v_k) - 2$. We can even assume without loss of generality that $k = 2$ (otherwise we proceed as in the proof of
Lemma 4.4 and change $p$ along squares walking backwards towards $v_{\min}$). But then $h(u_2) = h(v_2) - 2 = h(v_{\min})$. Minimality of $v_{\min}$ implies $u_2 = v_{\min}$, and therefore we get that $v_2$ was already adjacent to $v_{\min}$. This provides a path from $v_{\min}$ to $v$ of length at most $d_Q(v_{\min}, v) - 2$, which is a contradiction. \qed

We now turn to the study of links of vertices of $\mathrm{C}(\tilde{X}_\Gamma)$. Recall that $\mathrm{C}(\tilde{X}_\Gamma)$ is a cubical complex, hence its links are simplicial complexes. In particular, if $v \in \mathrm{C}(\tilde{X}_\Gamma)$ is a vertex, then vertices in $\link{v, \mathrm{C}(\tilde{X}_\Gamma)}$ correspond to vertices in $\mathrm{C}(\tilde{X}_\Gamma)$ which are adjacent to $v$. We begin with the following combinatorial characterization of simplices in the link of a vertex.

**Lemma 4.6.** Let $\sigma$ be a $k$-cell of $\tilde{X}_\Gamma$, and let $v$ be the dual vertex in $\mathrm{C}(\tilde{X}_\Gamma)$. Let $v_0, \ldots, v_m$ be a collection of vertices of $\mathrm{C}(\tilde{X}_\Gamma)$ adjacent to $v$, and let $\tau_0, \ldots, \tau_m$ be the dual cells in $\tilde{X}_\Gamma$. Then $v_0, \ldots, v_m$ induce a simplex in $\link{v, \mathrm{C}(\tilde{X}_\Gamma)}$ if and only if the following two conditions are satisfied

(↓) there exists a cell $\lambda$ of $\tilde{X}_\Gamma$ such that $\lambda \subseteq \tau_j$, $j = 0, \ldots, m$,

(↑) there exists a cell $\mu$ of $\tilde{X}_\Gamma$ such that $\tau_j \subseteq \mu$, $j = 0, \ldots, m$.

**Proof.** First of all, note that since $\mathrm{C}(\tilde{X}_\Gamma)$ is a cubical complex, the vertices $v_0, \ldots, v_m$ induce a simplex in $\link{v, \mathrm{C}(\tilde{X}_\Gamma)}$ if and only if there exists a cube $Q$ of $\mathrm{C}(\tilde{X}_\Gamma)$ containing $v, v_0, \ldots, v_m$.

Assume that they induce a simplex, and let $Q$ be the corresponding cube. From Lemma 4.4 we know that $Q$ has a unique vertex of minimal height, and a unique vertex of maximal height. Let $\lambda, \mu$ be the dual cells. Lemma 4.4 then implies that $\lambda, \mu$ satisfy the conditions (↓) and (↑) in the statement.

Vice versa suppose that the conditions (↓) and (↑) are satisfied. Notice that we have $\lambda \subseteq \tau_j \subseteq \mu$ for all $j = 0, \ldots, m$. Let $C_\lambda = \tilde{g}_X(\lambda)$ and $C_\mu = \tilde{g}_X(\mu)$ be the corresponding cubes of $X$, under the map $\tilde{g}_X = g_X \circ \pi : \tilde{X}_\Gamma \to X_\Gamma \to X$. Notice that $\link{\lambda, \mu} \cong \link{C_\lambda, C_\mu}$ by Lemma 3.17. In particular, we see that in $\tilde{X}_\Gamma$ there is a collection of cells containing $\lambda$ and contained in $\mu$ (among which we find the cells $\tau_j$) that gives rise to a cube $Q$ in $\mathrm{C}(\tilde{X}_\Gamma)$ containing the vertices $v, v_0, \ldots, v_m$. Therefore $v_0, \ldots, v_m$ induce a simplex in $\link{v, \mathrm{C}(\tilde{X}_\Gamma)}$, as desired. \qed

**Remark 4.7.** When the conditions (↓) and (↑) from Lemma 4.6 are satisfied, the cells $\lambda, \mu$ can be chosen to be the lower and upper cell provided by Lemma 3.19.

**Remark 4.8.** A cell of dimension at least 2 in $\tilde{X}_\Gamma$ always admits infinitely many codimension-1 cells (see Figure 11). Lemma 4.6 implies that the link of the dual vertex is neither compact nor connected. In particular the cubical complex $\mathrm{C}(\tilde{X}_\Gamma)$ is not locally compact. As a result, even though $\mathrm{C}(\tilde{X}_\Gamma)$ is constructed as a sort of dual cubical barycentric subdivision with respect to the combinatorial decomposition of $\tilde{X}_\Gamma$ into cells, $\mathrm{C}(\tilde{X}_\Gamma)$ is not homeomorphic to $\tilde{X}_\Gamma$. Namely, $\tilde{X}_\Gamma$ is locally compact, while $\mathrm{C}(\tilde{X}_\Gamma)$ is not locally compact.

As recalled above, if $v \in \mathrm{C}(\tilde{X}_\Gamma)$ is a vertex, then the vertices appearing in $\link{v, \mathrm{C}(\tilde{X}_\Gamma)}$ correspond to vertices of $\mathrm{C}(\tilde{X}_\Gamma)$ that are adjacent to $v$, and these vertices have height equal to $h(v) \pm 1$. We find it useful to decompose $\link{v, \mathrm{C}(\tilde{X}_\Gamma)}$ into
two subcomplexes: we denote by $lk_1 \left(v, \mathcal{C}(\tilde{X}_T)\right)$ the full subcomplex of $lk \left(v, \mathcal{C}(\tilde{X}_T)\right)$ generated by vertices of height $h(v) - 1$, and by $lk^T \left(v, \mathcal{C}(\tilde{X}_T)\right)$ the full subcomplex of $lk \left(v, \mathcal{C}(\tilde{X}_T)\right)$ generated by vertices of height $h(v) + 1$. As we will see, their geometry is controlled respectively by a certain Helly property for orthogonal hyperplanes in $\mathbb{H}^n$, and by the non–positive curvature of $X$. The following statement provides the Helly property. Notice that orthogonality is a key feature here: without the orthogonality requirement, the statement already fails for three geodesics in $\mathbb{H}^2$. On the other hand, the interested reader will notice that the argument generalizes to a collection of pairwise orthogonal and totally geodesic hypersurfaces in a simply connected complete manifold of non–positive curvature. We will not need this generality in this paper.

Lemma 4.9 (Helly property for orthogonal hyperplanes in $\mathbb{H}^n$). Let $\mathcal{V}$ be a collection of pairwise orthogonal hyperplanes in $\mathbb{H}^n$. Then $|\mathcal{V}| \leq n$, and for all $k \in \{2, \ldots, n\}$ all the $k$–fold intersections are non-empty.

Proof. We begin with some preliminary observation about orthogonal subspaces. Let $V_1, \ldots, V_k \in \mathcal{V}$ be a collection of hyperplanes from $\mathcal{V}$, and let $N = \cap_{j=1}^k V_j$ be their intersection. For $x \in N$, let $T_x(\mathbb{H}^n)$ denote the tangent space of $\mathbb{H}^n$ at $x$, and let $v_j \in T_x(\mathbb{H}^n)$ be a unit vector orthogonal to $V_j$ (i.e. to all vectors in the tangent space $T_x(V_j)$). The fact that $V_i$ and $V_j$ are orthogonal hyperplanes means that $v_i$ and $v_j$ are orthogonal vectors for all $i \neq j$. Then a direct computation shows that if $\{n_1, \ldots, n_m\}$ is an orthonormal basis for the tangent space of $T_x(N)$, then $\{n_1, \ldots, n_m, v_1, \ldots, v_k\}$ is an orthonormal basis for $T_x(\mathbb{H}^n)$. This shows in particular that $k \leq n$.

To prove the statement about non–emptiness of intersections, we notice that the case $k = 2$ is exactly the hypothesis that any pair of hyperplanes from $\mathcal{V}$ intersect. For $k \geq 3$, we argue that if all the $h$–fold intersections of elements from $\{V_1, \ldots, V_k\}$ are non–empty for all $h < k$, then the $k$–fold intersection is non–empty too.

Let $N_j = \cap_{i \neq j} V_i$. By assumption we have $N_j \neq \emptyset$. Assume by contradiction that $V_1 \cap \cdots \cap V_k = \emptyset$. Then for any choice of indices $j_1 \neq j_2$ we have that $N_{j_1} \cap N_{j_2} = \emptyset$. In particular, $N_2$ and $N_3$ are non–empty disjoint subspaces of $V_1$ (see Figure 26). Let $\gamma_1$ be the common perpendicular between them in $V_1$, and let $x_k \in N_k$ be its endpoint for $k = 2, 3$. Now in the tangent space $T_{x_2}(\mathbb{H}^n)$ we consider an orthonormal basis $\{n_1, \ldots, n_m, v_1, v_3, \ldots, v_k\}$ constructed as above by adding to an orthonormal basis $\{n_1, \ldots, n_m\}$ for $T_{x_2}(N_2)$ unit vectors $v_1, v_3, \ldots, v_k$ orthogonal to $V_1, V_3, \ldots, V_k$. If $w$ denotes a tangent vector at $x_2$ along $\gamma_1$, then a direct computation shows that $w$ is orthogonal to $\{n_1, \ldots, n_m\}$, because $\gamma_1$ is orthogonal to $N_2$, and it is also orthogonal to $v_1$, because $\gamma_1 \subseteq V_1$. Therefore $w$ is in the subspace of $T_{x_2}(\mathbb{H}^n)$ generated by $v_3, \ldots, v_k$. If we define $W_2 = \cap_{j=3}^k V_j$, then this means that $\gamma_1$ is orthogonal to $W_2$ at $x_2$. Arguing in the same way at the point $x_3$, we find that $\gamma_1$ is orthogonal at $x_3$ to the subspace $W_3 = \cap_{j=2, j \neq 3}^k V_j$. Note that $W_2 \cap W_3 = \cap_{j=2}^k V_j = N_1$ is non–empty. Moreover, as observed above, it is disjoint from $N_2$ and from $N_3$. Therefore we can connect $x_2$ (respectively $x_3$) to a point $x_1$ in $N_1$ with a geodesic arc $\gamma_2$ contained in $W_2$ (respectively $\gamma_3$ contained in $W_3$). Since all the spaces involved are totally geodesic, the arcs $\gamma_1, \gamma_2, \gamma_3$ are geodesic arcs in $\mathbb{H}^n$, so we have obtained a geodesic triangle with two right angles, which leads to the desired contradiction. \qed
The next statement completes our investigation of the combinatorial geometry of $\mathcal{C}(\widetilde{X}_\Gamma)$. Thanks to Gromov’s link condition (see Lemma 2.2), it already implies that $\mathcal{C}(\widetilde{X}_\Gamma)$ is locally CAT(0). We will show in Theorem 4.29 that it is actually CAT(0).

**Proposition 4.10.** Let $\sigma$ be a $k$-cell of $\widetilde{X}_\Gamma$, and let $v$ be the dual vertex in $\mathcal{C}(\widetilde{X}_\Gamma)$. Then the following hold:

1. $\text{lk}_\downarrow (v, \mathcal{C}(\widetilde{X}_\Gamma))$ is a flag simplicial complex.
2. $\text{lk}_\uparrow (v, \mathcal{C}(\widetilde{X}_\Gamma))$ is a flag simplicial complex.
3. $\text{lk} (v, \mathcal{C}(\widetilde{X}_\Gamma))$ is a flag simplicial complex.

**Proof.** Throughout this proof, $w_j$ will denote a vertex in $\text{lk}_\downarrow (v, \mathcal{C}(\widetilde{X}_\Gamma))$, $v_j$ the corresponding vertex of $\mathcal{C}(\widetilde{X}_\Gamma)$ adjacent to $v$, and $\tau_j$ the cell of $\widetilde{X}_\Gamma$ dual to $v_j$. Notice that two vertices $w_i, w_j$ are adjacent in $\text{lk}_\downarrow (v, \mathcal{C}(\widetilde{X}_\Gamma))$ precisely when $v, v_i, v_j$ are contained in a square of $\mathcal{C}(\widetilde{X}_\Gamma)$.

We first prove (1). Let $w_0, \ldots, w_p$ be pairwise adjacent vertices in $\text{lk}_\downarrow (v, \mathcal{C}(\widetilde{X}_\Gamma))$. Notice that $\tau_0, \ldots, \tau_p$ are all cells of codimension 1 in the boundary of $\sigma$. For each $i \neq j$, $v, v_i, v_j$ are contained in a square of $\mathcal{C}(\widetilde{X}_\Gamma)$. By Lemma 4.3, the fourth vertex of the square is dual to a cell contained in $\tau_i \cap \tau_j$. This shows that the cells $\tau_j$ intersect pairwise. We claim that actually $\tau_0 \cap \cdots \cap \tau_p \neq \emptyset$. To see this, embed $\sigma$ into a hyperbolic space of dimension $\dim \sigma$ (as in §3.3). The family of hyperplanes supporting the cells $\tau_j$ is a collection of pairwise orthogonal hyperplanes. By Lemma 4.9 their intersection is non-empty, hence the intersection of the cells $\tau_j$’s is non-empty by convexity. By Lemma 3.19, their intersection consists of a single cell $\lambda \subseteq \tau_j$. We use Lemma 4.6 with this cell $\lambda$ and $\mu = \sigma$ to conclude that $w_0, \ldots, w_k$ span a simplex.

We argue via a dual argument to prove (2). Let $w_0, \ldots, w_p$ be pairwise adjacent vertices in $\text{lk}_\uparrow (v, \mathcal{C}(\widetilde{X}_\Gamma))$. Notice that $\tau_0, \ldots, \tau_p$ are cells containing $\sigma$ as a cell of codimension 1 in their boundary. For each $i \neq j$, $v, v_i, v_j$ are contained in a
square of $\mathcal{C}(\tilde{X}_\Gamma)$. By Lemma 4.3, the fourth vertex of the square is dual to a cell containing $\tau_i \cup \tau_j$. So $\tau_i$, $\tau_j$ are adjacent in $\text{lk} \left( \sigma, \tilde{X}_\Gamma \right)$. By (1) in Lemma 3.17 this link is isomorphic to the link of the corresponding cube in $X$. Since $X$ is non-positively curved, this link is a flag simplicial complex. Therefore there is a cell $\mu$ containing all the cells $\tau_j$; this can actually be taken to be the upper cell provided by Lemma 3.19. We use Lemma 4.6 with this cell $\mu$ and $\lambda = \sigma$ to conclude that $w_0, \ldots, w_k$ span a simplex.

Finally, in order to prove (3), let $w_0, \ldots, w_p$ be pairwise adjacent vertices in $\text{lk} \left( v, \mathcal{C}(\tilde{X}_\Gamma) \right)$, ordered so that for some $m$ we have $w_0, \ldots, w_m \in \text{lk}_j \left( v, \mathcal{C}(\tilde{X}_\Gamma) \right)$ and $w_{m+1}, \ldots, w_p \in \text{lk}_j \left( v, \mathcal{C}(\tilde{X}_\Gamma) \right)$. By (1) we know that $w_0, \ldots, w_m$ span a simplex, hence by Lemma 4.6 there exists a cell $\lambda$ in $\cap_{j=0}^m \tau_j$. Similarly, by (2) we know that $w_{m+1}, \ldots, w_p$ span a simplex, hence by Lemma 4.6 there exists a cell $\mu$ containing $\tau_m, \ldots, \tau_p$. Notice that $\lambda \subseteq \tau_i \subseteq \sigma \subseteq \tau_j \subseteq \mu$ for all $i = 0, \ldots, m$ and $j = m+1, \ldots, p$. In particular we have $\lambda \subseteq \tau_j \subseteq \mu$ for all $j = 0, \ldots, p$. Using Lemma 4.6 again we obtain that $w_0, \ldots, w_p$ span a simplex. □

4.2. Efficiency. In this section we study a notion of complexity for edge–paths in $\mathcal{C}(\tilde{X}_\Gamma)$, which is based on the height function, and use it to find suitable representatives of homotopy classes of edge–paths and edge–loops. Recall that if $p$ is an edge–path in the 1–skeleton of a cubical complex, an elementary homotopy of $p$ is a homotopy which is contained in the 2–skeleton and is obtained by a finite sequence of the following two moves:

- remove a backtracking subpath, i.e. replace $(v_1, v_2, v_1)$ with $v_1$;
- slide across a square, i.e. replace $(v_1, v_2, v_3, v_4)$ with $(v_1, v_4, v_3)$ if $v_1, v_2, v_3, v_4$ appear in this order on the boundary of a square (as in Figure 23).

An edge–path $p = (v_0, \ldots, v_s)$ is said to be efficient if $\exists k \in \{0, \ldots, s\}$ such that the height strictly increases from $v_0$ to $v_k$ and strictly decreases from $v_k$ to $v_s$, i.e. $v_k$ is the unique point of maximum for the height along $p$. We allow $k = 0$ or $k = s$, i.e. that the height is strictly monotone along $p$. In any case, $h(p) = h(v_k)$, and the cell dual to $v_k$ contains the cells dual to all the other vertices of $p$. This implies that an efficient edge–path is contained in the union of at most two cubes which share at least a vertex. In particular, an efficient edge–loop is entirely contained in a single cube. These observations motivate the following definitions and constructions.

If $\tau$ is a tile of $\tilde{X}_\Gamma$, we define the dual tile $\mathcal{C}(\tau)$ to be the full subcomplex of $\mathcal{C}(\tilde{X}_\Gamma)$ whose vertices are dual to the cells of $\tau$. If $v$ is the vertex of $\mathcal{C}(\tilde{X}_\Gamma)$ which is dual to $\tau$, then $\mathcal{C}(\tau)$ consists of all the cubes of $\mathcal{C}(\tilde{X}_\Gamma)$ that contain $v$, i.e. $\mathcal{C}(\tau)$ is the combinatorial 1–neighborhood of $v$. Notice that $v$ is the only vertex of height $n$ in $\mathcal{C}(\tau)$ (see Figure 27). We say that an edge–path $p$ in $\mathcal{C}(\tilde{X}_\Gamma)$ stays in a tile if there exists a tile $\tau$ of $\tilde{X}_\Gamma$ such that $p \subseteq \mathcal{C}(\tau)$.

**Lemma 4.11.** Let $p$ be an edge–path in $\mathcal{C}(\tilde{X}_\Gamma)$. If $p$ stays in a tile, then there is an elementary homotopy relative to endpoints between $p$ and an efficient path.

**Proof.** Let $p = (v_0, \ldots, v_s)$. First of all, notice that if $s = 0, 1$ then $p$ is already efficient. Moreover, by an elementary homotopy relative to endpoints, we can assume that $p$ has no backtracking subpath. Since $p$ stays in a tile, $p$ goes through at most one vertex of height $n$ (possibly several times, possibly at the endpoints $v_0, v_s$).
For $0 < j < s$, we say $v_j$ is a local minimum (with respect to the height function along $p$) if $h(v_{j+1}) = h(v_j) + 1$, and we consider the following quantity
\[ h(p) = \min\{h(v_j) \mid v_j \text{ is a local minimum}\}. \]

If there is no local minimum, set $h(p) = \infty$; in this case $p$ is already efficient. So let us assume that there are some local minima, i.e., $h(p) < \infty$. Notice that $h(p)$ is in general larger than the minimum of the height along $p$. If $h(p) = n$ then $p$ is constant, hence efficient. If $h(p) = n - 1$, then $p$ has a backtracking subpath, because $p$ goes through at most one vertex of height $n$. By an elementary homotopy relative to endpoints we can remove this local minimum. Repeating this process, we obtain a path $p'$ with $h(p') = n$, and we reduce to the previous case. So let us assume in the following that $h(p) \leq n - 2$.

We now claim that, by deforming $p$ locally at local minima, we can produce an elementary homotopy relative to endpoints to a path $p'$ such that $h(p') \geq h(p) + 1$.

To prove the claim, let $v_j$ be a local minimum, and let its height be $h(v_j) = h_j$ for some $0 < j < s$. Consider the subpath $(v_{j-1}, v_j, v_{j+1})$, and note that the cells dual to $v_{j-1}, v_{j+1}$ meet along the cell dual to $v_j$. Since $p$ stays in a tile, there is a cell containing all these cells, namely the tile itself. By Lemma 4.6 we get that $(v_{j-1}, v_j, v_{j+1})$ is part of a square in $\mathcal{C}(X_T)$, whose fourth vertex is some $v'_j$, of height $h(v'_j) = h_j + 2$. Then we can homotope $(v_{j-1}, v_j, v_{j+1})$ to the other side $(v_{j-1}, v'_j, v_{j+1})$ of the square via an elementary homotopy relative to endpoints (see Lemma 4.3). This process can be applied to all local minima at the same time, since no two local minima can be adjacent along $p$. Then we remove all backtracking subpaths, if needed, keeping endpoints fixed. The result is an elementary homotopy relative to endpoints between $p$ and an edge–path $p'$ with $h(p') \geq h(p) + 1$. It is even possible that $h(p') = \infty$ but in any case this proves the claim.

We repeat this process of elevating local minima, and after a finite number of steps we obtain a path $p''$ with $h(p'') \geq n - 1$ (again, possibly $h(p'') = \infty$). Hence, we reduce to the previously discussed cases to conclude that $p''$ (hence $p$) admits an elementary homotopy relative to endpoints to an efficient path.

In the previous lemma we allow $p$ to be an edge–loop, i.e., $v_0 = v_s$. All the homotopies in it are relative to the base point $v_0 = v_s$. In the following statement we consider free homotopies, i.e., homotopies that are not required to fix any point.

**Corollary 4.12.** Let $p$ be an edge–loop in $\mathcal{C}(X_T)$. If $p$ stays in a tile, then there is an elementary homotopy between $p$ and a constant path.

**Proof.** Pick a basepoint $v_0$ on $p$ to be a vertex of maximal height on $p$, and write $p = (v_0, \ldots, v_s)$, for $v_0 = v_s$. Apply the previous argument (from Lemma 4.11) to $p$. At every iteration we allow ourselves to change the basepoint on $p$ to always be a vertex of maximal height. At the end there can be no local minimum, hence the path is constant.

A simple way for an edge–loop to satisfy the condition of Corollary 4.12 is to be short. Recall from § 4.1 that the length $\ell(p)$ of an edge–path $p$ is defined to be the number of edges of $p$.

**Corollary 4.13.** Let $p$ be an edge–loop in $\mathcal{C}(X_T)$. Then $\ell(p)$ is even. Moreover, if $\ell(p) \leq 4$ then $p$ stays in a tile, and there is an elementary homotopy between $p$ and a constant path.
Proof. The first statement follows from the fact that if an edge $e$ has endpoints $v, w$ then $|h(v) - h(w)| = 1$, so if an edge–path has odd length then the endpoints have different height.

Suppose now $\ell(p) \leq 4$. If $\ell(p) = 2$ then $p = (v, w, v)$ for two adjacent vertices $v, w$. In particular the cell dual to $v$ contains the one dual to $w$, or vice versa. If $\ell(p) = 4$ then $p$ is the boundary path of a square. It follows from Lemma 4.3 that $p$ contains a unique point of maximal height, and that the cell dual to it contains every other cell. In either case, there is a cell containing all the cells dual to the vertices of $p$. If $\tau$ is a tile of $\tilde{X}_\Gamma$ containing that cell, then $p$ is entirely contained in $\mathcal{C}(\tau)$ by construction. In particular, $p$ stays in a tile, so the statement follows from Corollary 4.12. 

\[\begin{figure}[h]
\centering
\includegraphics[width=\textwidth]{figure27.png}
\caption{A dual tile in $\mathcal{C}(\tilde{X}_\Gamma)$, and a long edge–path that stays in a tile.}
\end{figure}\]

Remark 4.14. From now on, our main goal in this section will be to show that every edge–loop in $\mathcal{C}(\tilde{X}_\Gamma)$ can be written as a product of nullhomotopic edge–loops, i.e. $\mathcal{C}(\tilde{X}_\Gamma)$ is simply connected. A naive approach would consist in splitting an edge–loop along mirrors into shorter edge–paths, until they are short enough to be contracted (in the sense of Corollary 4.12). However, there are arbitrarily long edge–paths that stay in a tile (see Figure 27). Therefore, an inductive argument based on length alone would not suffice, and this idea requires some additional tools which we develop in §4.3, before returning to the problem of simple connectedness of $\mathcal{C}(\tilde{X}_\Gamma)$ in §4.4.

Remark 4.15. Given two cells $\sigma, \sigma'$ contained in the same tile $\tau$, let $\mu = \mu(\sigma, \sigma')$ be their upper cell (i.e. the smallest cell that contains both of them, as defined in Lemma 3.19). If $v, v'$ and $w$ are the vertices dual to $\sigma, \sigma'$ and $\mu$ respectively, then an edge–path of minimal length in $\mathcal{C}(\tilde{X}_\Gamma)$ from $v$ to $v'$ can be obtained as an efficient path $p$ in $\mathcal{C}(\tau)$ going through $w$. Such an efficient edge–path is not unique, but the length of any such path is given by

$$\ell(p) = 2h(w) - h(v) - h(v') = 2\dim \mu - \dim \sigma - \dim \sigma'.$$
It should be noted that if $\mu \subseteq \tau$ then there are edge–paths from $v$ to $v'$ which are strictly longer than $p$ but still efficient.

4.3. Mirror complexity. Here we define an additional notion of complexity for an edge–path, based on the relative position in $\tilde{X}_\Gamma$ between mirrors and the cells dual to the vertices of the edge–path. We start with the following definition, in analogy to that of a dual tile. If $M$ is a mirror of $\tilde{X}_\Gamma$, we define the dual mirror $\mathcal{C}(M)$ to be the full subcomplex of $\mathcal{C}(\tilde{X}_\Gamma)$ whose vertices are dual to the cells of $M$. Since we have not proved yet that $\mathcal{C}(\tilde{X}_\Gamma)$ is simply connected, a priori it is not clear that a dual mirror enjoys properties reminiscent of those of a mirror of $\tilde{X}_\Gamma$; for instance, it is not clear yet whether it is convex. Nevertheless, we can obtain the following statement about separation (analogous to Proposition 3.37).

**Lemma 4.16.** Let $M$ be a mirror of $\tilde{X}_\Gamma$ and let $\mathcal{C}(M)$ be the dual mirror in $\mathcal{C}(\tilde{X}_\Gamma)$. Let $z_1, z_2$ be two points in $\tilde{X}_\Gamma \setminus M$, let $\sigma_1, \sigma_2$ be cells in $\tilde{X}_\Gamma$ such that $z_k \in \sigma_k$, and let $v_k$ be the vertex of $\mathcal{C}(\tilde{X}_\Gamma)$ dual to $\sigma_k$. Then $M$ separates $z_1$ and $z_2$ if and only if $\mathcal{C}(M)$ separates $v_1$ and $v_2$. In particular, $\mathcal{C}(M)$ separates $\mathcal{C}(\tilde{X}_\Gamma)$.

**Proof.** Suppose $M$ separates $z_1$ and $z_2$, and assume by contradiction that there is an edge–path $p$ in $\mathcal{C}(\tilde{X}_\Gamma)$ from $v_1$ to $v_2$ avoiding $\mathcal{C}(M)$. Then the union of the cells dual to the vertices of $p$ contains a path–connected subspace of $\tilde{X}_\Gamma \setminus M$ that contains both $z_1$ and $z_2$. This is in contradiction with the fact that $M$ separates $z_1$ from $z_2$. 

Vice versa, suppose $\mathcal{C}(M)$ separates $v_1$ and $v_2$, and assume by contradiction that there is a path $\gamma$ in $\tilde{X}_\Gamma$ from $z_1$ to $z_2$ avoiding $M$. By a small perturbation, we can assume that $\gamma$ intersects the strata of $\tilde{X}_\Gamma$ in such a way that the sequence of the minimal cells that it visits gives rise to an edge–path in $\mathcal{C}(\tilde{X}_\Gamma)$ (i.e. their dimension jumps by 1 at a time along $\gamma$). By construction, such an edge–path connects $v_1$ to $v_2$ in the complement of $\mathcal{C}(M)$, which is not possible.

In particular, it follows that $\mathcal{C}(M)$ separates $\mathcal{C}(\tilde{X}_\Gamma)$, because $M$ separates $\tilde{X}_\Gamma$ by Proposition 3.37.

This provides a correspondence between complementary components of a mirror $M$ in $\tilde{X}_\Gamma$ and complementary components of the dual mirror $\mathcal{C}(M)$ in $\mathcal{C}(\tilde{X}_\Gamma)$.

4.3.1. Crossings. Let $p = (v_0, \ldots, v_s)$ be an edge–path (possibly an edge–loop) in $\mathcal{C}(\tilde{X}_\Gamma)$, and let $\sigma_0, \ldots, \sigma_s$ be the cells of $\mathcal{C}(\tilde{X}_\Gamma)$ dual to its vertices. Let $M$ be a mirror in $\tilde{X}_\Gamma$, and let $\mathcal{C}(M)$ be the dual mirror in $\mathcal{C}(\tilde{X}_\Gamma)$. Recall from Lemma 4.16 that $\mathcal{C}(\tilde{X}_\Gamma) \setminus \mathcal{C}(M)$ is disconnected. We say that $p$ crosses $M$ if $p \cap \mathcal{C}(M) \neq \emptyset$ and there are at least two connected components $C_1, C_2$ of $\mathcal{C}(\tilde{X}_\Gamma) \setminus \mathcal{C}(M)$ such that $p \cap C_k \neq \emptyset$. This means that among the cells $\sigma_0, \ldots, \sigma_s$, some are contained in $M$, but at least two of them are such that their interiors are contained in different complementary components of $M$. (Recall that in our setting cells are closed and complementary components of mirrors are open.) Let $q = (v_j, \ldots, v_{j+m})$ be a subpath of $p$. We say that $q$ is a $(p, M)$–crossing if $v_j, \ldots, v_{j+m} \in \mathcal{C}(M)$, but $v_{j-1}$ and $v_{j+m+1}$ lie in different connected components of $\mathcal{C}(\tilde{X}_\Gamma) \setminus \mathcal{C}(M)$. In other words, $q$ is a subpath of $p \cap \mathcal{C}(M)$ which is maximal among subpaths of $p \cap \mathcal{C}(M)$ contained in the closure of a single connected component of $\mathcal{C}(\tilde{X}_\Gamma) \setminus \mathcal{C}(M)$. (See Figure 28 for some examples.) We denote by $m(p, M)$ the number of $(p, M)$–crossings. The
mirror complexity of $p$ is defined by taking into account the family $M$ of all mirrors of $\tilde{X}_\Gamma$, i.e. by the following formula:

$$m(p) = \sum_{M \in M} m(p, M).$$

Figure 28. An edge–path $p$ in $\mathcal{C}(\tilde{X}_\Gamma)$ crossing some mirrors. Mirror crossings are highlighted. We have $m(p, M) = 2$, $m(p, N_1) = 1$, $m(p, N_2) = 3$, and $m(p, N_3) = 3$. In particular, notice that even if $p$ intersects $\mathcal{C}(N_1)$ twice, there is only one $(p, N_1)$-crossing.

The relevance of this notion of complexity with respect to Remark 4.14 is showcased by the following two lemmas.

**Lemma 4.17.** Let $p$ be an edge–path in $\mathcal{C}(\tilde{X}_\Gamma)$. Then $p$ stays in a tile if and only if $p$ does not cross any mirror.

**Proof.** Suppose that $p$ stays in a tile, i.e. there exists a tile $\tau$ such that $p \subseteq \mathcal{C}(\tau)$. Assume by contradiction that $p$ crosses a mirror $M$. So there are two vertices $v_1, v_2$ of $p$ which are separated by $\mathcal{C}(M)$. Let $\sigma_k$ be the cell dual to $v_k$. By Lemma 4.16 $M$ separates the interior of $\sigma_1$ from the interior of $\sigma_2$. In particular, there is no tile of $\tilde{X}_\Gamma$ that contains both of them, which contradicts the hypothesis that $p$ stays in a tile.

Vice versa, suppose $p$ does not cross any mirror, and assume by contradiction that there are two vertices $v_1, v_2$ on $p$ such that the dual cells are not contained in the same tile. Let $\tau_1, \tau_2$ be different tiles containing them. Up to choosing $v_1, v_2$ closer to each other along $p$, we can assume that the tiles are adjacent, i.e. $\tau_1 \cap \tau_2 \neq \emptyset$. In particular, $\sigma = \tau_1 \cap \tau_2$ is a cell and it is contained in some mirror $M$. Then $p$ intersects $M$ between $v_1$ and $v_2$. Moreover the tiles $\tau_1, \tau_2$ provide a framing in the sense of §3.7. Proposition 3.37 implies that the interiors of $\tau_1, \tau_2$ are separated by $M$. The same holds for the interiors of the cells dual to $v_1, v_2$. 


So by Lemma 4.16 we have that \( v_1, v_2 \) are separated by \( \scr C(M) \), i.e. \( p \) crosses \( M \), a contradiction. \( \Box \)

Lemma 4.18. Let \( p \) be an edge–path in \( \scr C(\tilde X_\Gamma) \), and let \( M \) be a mirror in \( \tilde X_\Gamma \). Then the following hold.

1. \( \mathfrak m(p, M) = 0 \) if and only if \( p \) does not cross \( M \).
2. \( \mathfrak m(p) = 0 \) if and only if \( p \) stays in a tile.
3. If \( p \) is a loop and \( \mathfrak m(p, M) \geq 1 \), then \( \mathfrak m(p, M) \geq 2 \).
4. If \( p \) has finite length, then \( \mathfrak m(p, M) \) and \( \mathfrak m(p) \) are finite.

Proof. For (1), note that \( \mathfrak m(p, M) \) is by definition the number of \( (p, M) \)-crossings. For (2), note that \( \mathfrak m(p) \) is a sum of non-negative numbers, so it is zero if and only if \( \mathfrak m(p, M) = 0 \) for every mirror \( M \). By (1) this is equivalent to saying that \( p \) does not cross any mirror. Then the statement follows from Lemma 4.17. To prove (3), note that if \( p \) is a loop that crosses \( M \) at least once, then it must cross it at least twice, because \( \scr C(M) \) separates \( \scr C(\tilde X_\Gamma) \) by Lemma 4.16.

Finally, to prove (4) notice that each \( (p, M) \)-crossing contributes to at least one vertex of \( p \), dual to a cell of \( M \). Since \( p \) has finite length, there can be only finitely many \( (p, M) \)-crossings. Then the finiteness of \( \mathfrak m(p) \) follows from the fact that \( X \) (hence the collection of mirrors \( M \)) is locally finite. \( \Box \)

Remark 4.19. In this framework, Corollary 4.13 can be restated by saying that \( \ell(p) \leq 4 \) implies \( \mathfrak m(p) = 0 \).

4.3.2. Bridges. Let \( p = (v_0, \ldots, v_s) \) be an edge–path in \( \scr C(\tilde X_\Gamma) \), and let \( \sigma_0, \ldots, \sigma_s \) be the cells of \( \scr C(\tilde X_\Gamma) \) dual to its vertices. Let \( M \) be a mirror in \( \tilde X_\Gamma \), and let \( \scr C(M) \) be the dual mirror in \( \scr C(\tilde X_\Gamma) \). We say that \( p \) is a bridge if there exists a mirror \( M \) of \( \tilde X_\Gamma \) such that \( v_0, v_s \in \scr C(M) \), but \( p \not\subseteq \scr C(M) \). In other words, \( \sigma_0, \sigma_s \subseteq M \) but some of the other cells \( \sigma_1, \ldots, \sigma_{s-1} \) are not contained in \( M \). In this case, we say that \( p \) is supported by \( M \). We say \( p \) is a minimal bridge if none of its subpaths is a bridge (see Figure 29).

Lemma 4.20. Let \( p \) be an edge–path in \( \scr C(\tilde X_\Gamma) \). If \( p \) is a bridge, then there exists a subpath \( q \subseteq p \) that is a minimal bridge.

Proof. Let us consider the collection of subpaths of \( p \) which are are bridges. Notice that this collection contains \( p \) itself, it is partially ordered by inclusion, and it is finite. Therefore it contains a minimal element. \( \Box \)

Lemma 4.21. Let \( p \) be a minimal bridge supported by a mirror \( M \), and let \( N \) be a mirror such that \( \mathfrak m(p, N) > 0 \). Then the following hold.

1. \( \mathfrak m(p, N) = 1 \).
2. \( \scr C(M) \cap \scr C(N) \neq \emptyset \) and \( M \cap N \neq \emptyset \).

Proof. The assumption that \( \mathfrak m(p, N) > 0 \) means that \( p \) crosses \( N \) at least once. If \( p \) crossed \( N \) twice, then any subpath between two consecutive \( (p, N) \)-crossings would be a bridge supported by \( N \). But this would contradict minimality, hence \( \mathfrak m(p, N) = 1 \), which proves (1). In particular the endpoints of \( p \) lie in different connected components of \( \scr C(\tilde X_\Gamma) \setminus \scr C(N) \). Since they also live on the same dual mirror \( \scr C(M) \), which is connected, and dual mirrors separate by Lemma 4.16, we can conclude that \( \scr C(M) \cap \scr C(N) \neq \emptyset \). Finally, the cell dual to a vertex in their intersection is contained in \( M \cap N \), hence we obtain (2). \( \Box \)
Recall that for a mirror $M$ in $\widetilde{X}_{\Gamma}$ we have a nearest point projection $\pi_M : \widetilde{X}_{\Gamma} \to M$, as discussed in §3.6. If $p$ is a minimal bridge supported on $M$, then we can use $\pi_M$ to induce a projection of $p$ to $\mathcal{C}(M)$, as established by the next results.

**Lemma 4.22.** Let $M, N$ be mirrors in $\widetilde{X}_{\Gamma}$ and let $\tau$ be a tile in $\widetilde{X}_{\Gamma}$.

1. If $M \cap N \neq \emptyset$, then $\pi_M(N) = M \cap N$.
2. If $M \cap \tau \neq \emptyset$, then $\pi_M(\tau) = M \cap \tau$.

**Proof.** We start by proving $\pi_M(N) \subseteq M \cap N$. By contradiction, let $x \in N$ be such that $\pi_M(x) \notin N$. Let $y = \pi_N(\pi_M(x)) \in M \cap N$, where $\pi_N$ denotes the nearest point projection to the mirror $N$. Since $\pi_M(x) \notin N$, $y \neq \pi_M(x)$, so we can consider the geodesic triangle with vertices $x, \pi_M(x)$ and $y$. By convexity of $M$, the geodesic $[\pi_M(x), y]$ is contained in $M$. By convexity of $N$, the geodesic $[y, x]$ is contained in $N$. Moreover, since $\pi_M$ is the nearest point projection to $M$, the angle between $[x, \pi_M(x)]$ and $[\pi_M(x), y]$ at $\pi_M(x)$ is at least $\frac{\pi}{2}$. Analogously, the angle between $[\pi_M(x), y]$ and $[y, x]$ at $y$ is at least $\frac{\pi}{2}$ too. We obtained a geodesic triangle with two angles larger than $\frac{\pi}{2}$, which is impossible in the CAT(0) space $\widetilde{X}_{\Gamma}$. Vice versa, if $x \in M \cap N$, then $x = \pi_M(x)$, so $x \in \pi_M(N)$ already.

The second statement can be obtained via an analogous argument. Indeed, recall that $\tau$ is isometric to a convex subset of $\mathbb{H}^n$ bounded by orthogonal hyperplanes (see Lemma 3.11). In particular, the nearest point projection to a boundary face of $\tau$ is entirely contained in $\tau$. \qed

The next lemma is a combinatorial statement about the stratification of $\widetilde{X}_{\Gamma}$ introduced in §3.5, and will be needed in the following lemma.

**Lemma 4.23.** Let $\tau, \tau'$ be non-disjoint tiles of $\widetilde{X}_{\Gamma}$. Let $W_1, \ldots, W_q$ be the collection of mirrors of $\widetilde{X}_{\Gamma}$ that separate $\tau$ and $\tau'$. Then we have that

1. $W_1, \ldots, W_q$ coincides with the collection of mirrors of $\widetilde{X}_{\Gamma}$ that contain $\tau \cap \tau'$.
2. $\tau \cap W_1 \cap \cdots \cap W_q = \tau \cap \tau' = \tau' \cap W_1 \cap \cdots \cap W_q$.

**Proof.** First of all, notice that the collection of mirrors is not empty since $\tau$ and $\tau'$ are different tiles. We start by proving (1). Let $W$ be a mirror containing
τ ∩ τ′. Then the two tiles provide a framing for the cell τ ∩ τ′. In particular we get from Proposition 3.37 that W separates the two tiles, hence W is in the collection \{W₁, ..., W₉\}. Conversely, if τ ∩ τ′ was not inside one Wᵢ, then we could connect the two tiles with a path that goes through the intersection but avoids Wᵢ, contradicting the fact that Wᵢ separates them.

To prove (2) we argue as follows. By (1) we know that τ ∩ τ′ ⊆ W₁ ∩ ⋯ ∩ W₉, so we have that τ ∩ τ′ ⊆ τ ∩ W₁ ∩ ⋯ ∩ W₉. Note now that, by definition of the stratification, if τ ∩ τ′ is a k-cell, then it must be contained in n − k mirrors, so q = n − k. But then the two sides of the inclusion are cells of the same dimension k, so they must be equal. Switching the roles of τ and τ′ proves the second equality in (2).

Lemma 4.24. Let τ be a tile and M be a mirror in ⃗XΓ, such that M ∩ τ ≠ ∅. Let σ be a cell of τ, and let N₁, ..., Nᵣ ≠ M be all the mirrors containing σ and such that M ∩ Nᵢ ≠ ∅ for j = 1, ..., r. (Possibly r = 0 if there are no such mirrors.) Then the following hold.

1. τ ∩ M ∩ N₁ ∩ ⋯ ∩ Nᵣ is an (n − 1 − r)-cell that contains πₘ(σ).
2. The cell τ ∩ M ∩ N₁ ∩ ⋯ ∩ Nᵣ only depends on σ and M.

Proof. We start by proving (1). It follows from Lemma 4.22 that πₘ(σ) ⊆ τ ∩ M ∩ Nᵢ for each j = 1, ..., r. Since σ ⊆ τ ∩ N₁ ∩ ⋯ ∩ Nᵣ, we obtain that πₘ(σ) ⊆ τ ∩ M ∩ N₁ ∩ ⋯ ∩ Nᵣ. To show that this intersection is a cell, note that τ is an n-cell. So, by Lemma 3.21 we have that M ∩ τ is an (n − 1)-cell and then for each j = 1, ..., r we have that τ ∩ M ∩ Nᵢ ∩ ⋯ ∩ Nᵣ is an (n − 1 − j)-cell.

To prove (2) we argue as follows. Suppose τ′ is another tile as in the statement, i.e. σ ⊆ τ′ and M ∩ τ′ ≠ ∅. Let W₁, ..., W₉ be the collection of mirrors of ⃗XΓ that separate τ and τ′. (Note that this collection depends on τ and τ′, while the collection N₁, ..., Nᵣ only depends on σ and M.) Since σ ⊆ τ ∩ τ′, we also have that σ is contained in each Wᵢ thanks to (1) in Lemma 4.23. We now claim that each Wᵢ meets M. This is clear if σ ⊆ M. On the other hand, if σ is not inside M, then we can take an efficient edge-path p in C(⃗XΓ) from the vertex dual to σ to the vertex dual to M ∩ τ which is contained in C(τ) and meets C(M) only at the endpoint M ∩ τ. Take an analogous path p′ in C(τ′), and concatenate p and p′ to obtain a minimal bridge p supported on M. Since Wᵢ separates τ and τ′, we see that p crosses Wᵢ. So by (2) in Lemma 4.21 we conclude that M ∩ Wᵢ ≠ ∅, which proves the claim.

As a result, we have that the collection \{W₁, ..., W₉\} is a subcollection of \{M, N₁, ..., Nᵣ\}. (Note that M could be one of the mirrors separating τ and τ′, but Nᵣ ≠ M by definition.) In particular, using (2) from 4.23, we obtain that

τ ∩ M ∩ N₁ ∩ ⋯ ∩ Nᵣ ⊆ τ ∩ W₁ ∩ ⋯ ∩ W₉ ⊆ τ ∩ τ′

Therefore it follows that τ ∩ M ∩ N₁ ∩ ⋯ ∩ Nᵣ ⊆ τ′ ∩ M ∩ N₁ ∩ ⋯ ∩ Nᵣ. Reversing the roles of τ and τ′ provides the other inclusion, and shows that the cell defined in (1) does not depend on the choice of the tile.

In the notation and setting of Lemma 4.24, if v ∈ C(⃗XΓ) is the vertex dual to σ, then we denote by πₘ(v) the vertex dual to the cell constructed in (1) of the lemma, and call it the projection of v to C(M). This is well defined by (2) in the same lemma. Note that in general 0 ≤ r ≤ n − dim σ, as σ could be contained
in some mirrors that are disjoint from $M$. Nevertheless, this provides the desired projection to $\mathcal{C}(M)$ for vertices of $\mathcal{C}(\tilde{X})$ which are contained in the cubical 2–neighborhood of $\mathcal{C}(M)$, i.e. the union of all the dual tiles corresponding to all the tiles that intersect $M$ in $\tilde{X}$.

The content of the next two lemmas is that a minimal bridge supported by a mirror $M$ is completely contained in such a neighborhood of $\mathcal{C}(M)$ (see Lemma 4.25), so we can define a projection of a minimal bridge to $\mathcal{C}(M)$ (see Lemma 4.26). We note that the minimality assumption is necessary, see the difference between $q$ and $q'$ in Figure 29.

Lemma 4.25. Let $p$ be a minimal bridge supported on a mirror $M$. Then for each vertex $v$ of $p$ there exists a tile $\tau$ such that $v \in \mathcal{C}(\tau)$ and $\tau \cap M \neq \emptyset$.

Proof. Let $p = (v_0, \ldots, v_s)$, and assume by contradiction that some vertices do not satisfy the statement. Let $v_k$ be the first one. Since $p$ is a bridge, its endpoints are on $\mathcal{C}(M)$, so $k \neq 0, s$. Let $\tau_\pm$ be tiles such that $v_{k-1} \in \mathcal{C}(\tau_-)$ and $v_k \in \mathcal{C}(\tau_+)$. By construction, we can choose $\tau_-$ so that $\tau_- \cap M \neq \emptyset$, while necessarily $\tau_+$ is disjoint from $M$.

Consider the cell $\sigma = \tau_- \cap \tau_+$. For any mirror $N$ containing $\sigma$, we claim that $N$ must intersect $M$. Indeed, the tiles $\tau_\pm$ form a framing for $\sigma$ in the sense of §3.7. By Proposition 3.37 we know that $\tau_\pm$ belong to the closure of distinct complementary components of $N$. In particular, a maximal subpath of $p \cap \mathcal{C}(N)$ whose vertices are dual to cells contained in $\sigma$ gives rise to a $(p, N)$-crossing, hence $m(p, N) > 0$. By (2) in Lemma 4.21 we know $M \cap N \neq \emptyset$, as claimed.

Let $N_1, \ldots, N_r$ be the collection of all mirrors containing $\sigma$. Since $\sigma$ is a cell of $\tau_-$, we can write $\sigma = \tau_- \cap N_1 \cap \cdots \cap N_r$. Using (1) in Lemma 4.24, we have that

$$\tau_M(\sigma) \subseteq \tau_- \cap M \cap N_1 \cap \cdots \cap N_r$$

$$= (M \cap \tau_-) \cap (\tau_- \cap N_1 \cap \cdots \cap N_r) \subseteq M \cap \sigma \subseteq M \cap \tau_+$$

This contradicts the fact that $\tau_+$ is disjoint from $M$. $\Box$

In the next lemma we finally obtain a projection of a minimal bridge to a supporting mirror. As it might be expected, such a projection is length–decreasing (see Figure 30).

Lemma 4.26. Let $p$ be a minimal bridge supported on a mirror $M$. Then there exists an edge–path $p^M \subseteq \mathcal{C}(M)$, such that $p^M$ has the same endpoints as $p$ and $\ell(p^M) \leq \ell(p) - 2$.

Proof. Let $p = (v_0, \ldots, v_s)$, and let $\sigma_0, \ldots, \sigma_s$ be the cells dual to its vertices. Since $p$ is a minimal bridge supported on $M$, by Lemma 4.25 we know that for each vertex $v_k$ there exists a tile $\tau_k$ of $\tilde{X}$ such that $v_k \in \mathcal{C}(\tau_k)$ and $\tau_k \cap M \neq \emptyset$. Let $w_k = \pi_M(v_k)$ be the projection of $v_k$ to $\mathcal{C}(M)$, constructed in Lemma 4.24. We claim that for each $k$ the vertices $w_{k-1}$ and $w_k$ are either the same vertex or adjacent vertices.

To see this, consider two vertices $v_{k-1}$ and $v_k$ adjacent along $p$. Without loss of generality we can assume that $\sigma_{k-1}$ is a cell of codimension 1 in $\sigma_k$. In particular, we can take $\tau_{k-1} = \tau_k$, and there is exactly one mirror $\tilde{N}_k$ that contains $\sigma_{k-1}$ but not $\sigma_k$. Let $\{\hat{N}_1, \ldots, \hat{N}_r\}$ be the collection of all the mirrors that contain $\sigma_k$ and intersect $M$, but are different from $M$. Then the analogous collection for $\sigma_{k-1}$ consists of the same mirrors, possibly with the addition of $\tilde{N}_k$. (Note that
since \( p \) is a minimal bridge supported on \( M \), any mirror containing \( \sigma_1, \ldots, \sigma_{s-1} \) is guaranteed to be different from \( M \), while \( \hat{N}_k = M \) for \( k = 1 \).) By (1) in Lemma 4.24 we have that \( \pi_M(\sigma_k) \subseteq \tau_k \cap M \cap N_1 \cap \cdots \cap N_r \) and that either \( \pi_M(\sigma_{k-1}) \subseteq \tau_k \cap M \cap N_1 \cap \cdots \cap N_r \cap \hat{N}_k \) or \( \pi_M(\sigma_{k-1}) \subseteq \tau_k \cap M \cap N_1 \cap \cdots \cap N_r \cap \hat{N}_k \). In the first case we have that \( \pi_M(\sigma_{k-1}) \) and \( \pi_M(\sigma_k) \) are contained in the intersection of the same mirrors, hence \( w_{k-1} = w_k \); in the second case \( \tau_k \cap M \cap N_1 \cap \cdots \cap N_r \cap \hat{N}_k \) is a codimension-1 cell of \( \tau_k \cap M \cap N_1 \cap \cdots \cap N_r \), hence \( w_{k-1} \) is adjacent to \( w_k \). This proves the claim.

Notice in particular that in the case \( k = 1 \) we have \( \hat{N}_k = M \), so we have proved that \( w_0 = w_1 \). Analogously, we also have \( w_s = w_{s-1} \). As a result, \((w_0, \ldots, w_s)\) is an edge-path in \( \mathcal{C}(M) \). Let \( p^M \) be the edge-path obtained from \((w_0, \ldots, w_s)\) by removing all backtracking subpaths and repeated vertices. In particular, since \( w_0 = w_1 \) and \( w_s = w_{s-1} \), we have that \( \ell(p^M) \leq s - 2 = \ell(p) - 2 \). Moreover, since \( p \) is a bridge supported on \( M \), we have that \( \sigma_0, \sigma_s \subseteq M \), so that \( v_0 = w_0, v_s = w_s \), i.e. \( p \) and \( p^M \) have the same endpoints. □

4.4. Surgeries on edge-loops. We are now ready to apply the above technology to the study of edge-loops in \( \mathcal{C}(\tilde{X}_\Gamma) \). The goal is to show that \( \mathcal{C}(\tilde{X}_\Gamma) \) is simply connected. The strategy will be to reduce the length and mirror complexity of an edge-loop enough to ensure that it stays in a tile, so that Corollary 4.12 can be applied. The following statement is the key surgery step. Roughly speaking, we chop an edge-loop \( p \) along a mirror \( M \) that it crosses, use the projection \( p^M \) of \( p \) to \( M \) to introduce a shortcut along \( M \) and obtain two edge-loops \( p_1, p_2 \) such that \( p \) and \( p_1p_2 \) are elementary homotopic, and finally then check that the lengths have dropped.

**Proposition 4.27.** Let \( p \) be an edge-loop in \( \mathcal{C}(\tilde{X}_\Gamma) \). If \( m(p) > 0 \), then there exist two edge-loops \( p_1, p_2 \) in \( \mathcal{C}(\tilde{X}_\Gamma) \) such that \( \ell(p_1), \ell(p_2) < \ell(p) \), and there is an elementary homotopy between \( p \) and \( p_1p_2 \).
Proof. By assumption, there is a mirror $M_0$ that is crossed by $p$, so $m(p, M_0) \geq 1$, hence by (3) in Lemma 4.18 we have that $m(p, M_0) \geq 2$, i.e. $p$ crosses $M_0$ at least twice. It follows from the definitions that any subarc of $p$ between any two $(p, M_0)$-crossings is a bridge supported by $M_0$.

Choose a $(p, M_0)$-crossing and a bridge $q$ supported by $M_0$ (in general this cannot be chosen to be minimal, see Figure 31). Let $q'$ be the complement of $q$ in $p$, i.e. the edge–path such that $p = qq'$. Of course we have
\[
\ell(q) + \ell(q') = \ell(p)
\]
Moreover, without loss of generality we can assume that
\[
\ell(q) \leq \ell(q').
\]
By Lemma 4.20 we can find a minimal bridge $q_1 \subseteq q \subseteq p$. In particular we have
\[
\ell(q_1) \leq \ell(q).
\]
Let $q_2$ be the complement of $q_1$ in $p$, i.e. the edge–path such that $p = q_1q_2$. We can compute that
\[
\ell(q_2) = \ell(p) - \ell(q_1) \leq \ell(q) - \ell(q_1) \geq \ell(q') \geq \ell(q_1).
\]
Let $M$ be a mirror supporting the minimal bridge $q_1$, and let $q_1^M$ be the projection of $q_1$ to $C(M)$, i.e. the edge–path obtained in Lemma 4.26. In particular we have
\[
\ell(q_1^M) \leq \ell(q_1) - 2 < \ell(q_1).
\]
Define the edge–loops $p_1 = q_1q_1^M$ and $p_2 = q_1^Mq_2$, where $q_1^M$ denotes the edge–path $q_1^M$ with the opposite orientation. There is an elementary homotopy between the edge–loops $p = q_1q_2$ and $p_1p_2 = q_1q_1^Mq_1^Mq_2$, obtained by removing the backtracking subpath $q_1^Mq_1^M$. We can compute the desired inequality on the length of $p_1$ and $p_2$ as follows:
\[
\ell(p_1) = \ell(q_1) + \ell(q_1^M) < \ell(q_1) + \ell(q_1) \leq \ell(q_1) + \ell(q_2) = \ell(p).
\]
\[
\ell(p_2) = \ell(q_1^M) + \ell(q_2) < \ell(q_1) + \ell(q_2) = \ell(p).
\]
\[\square\]

Remark 4.28. Note that it is possible to have an edge–loop $p$ for which the surgery from Proposition 4.27 strictly reduces the mirror complexity for only one subloop. For an example see Figure 31, where the mirror complexity of the loop $p_2$ is the same as that of the original loop $p$.

We are now ready to prove the main result of this section.

Theorem 4.29. The complex $C(X_\Gamma)$ is a connected CAT(0) cubical complex.

Proof. By construction, the complex $C(X_\Gamma)$ is a cubical complex. Moreover, Gromov’s link condition from Lemma 2.2 implies that $C(X_\Gamma)$ is non–positively curved, since the link of any vertex is a flag simplicial complex by Proposition 4.10.

Next, $C(X_\Gamma)$ is path–connected because $X_\Gamma$ is path–connected. Indeed, let $v, w$ be vertices in $C(X_\Gamma)$, and let $\sigma_v, \sigma_w$ be the dual cells. Pick any continuous path $\eta$ connecting the two cells in $X_\Gamma$, and keep track of the list of cells that are intersected by $\eta$. By isotoping $\eta$ into lower-dimensional cells, we can ensure that the difference
between the dimension of two consecutive cells in this list is exactly 1. The dual vertices in $\mathcal{C}(\tilde{X}_\Gamma)$ give rise to an edge–path from $v$ to $w$.

To conclude, we need to show that $\mathcal{C}(\tilde{X}_\Gamma)$ is simply connected. To do this, we argue that edge–loops are nullhomotopic by induction on their length. Let $p$ be an edge–loop in $\mathcal{C}(\tilde{X}_\Gamma)$, homotopically non trivial and of minimal length. If $p$ does not cross any mirror, then by Lemma 4.17 $p$ stays in a tile. Hence by Corollary 4.12 there is an elementary homotopy between $p$ and a constant path. So let us assume that $m(p) > 0$. Then by Proposition 4.27 there exist two edge–loops $p_1, p_2$ in $\mathcal{C}(\tilde{X}_\Gamma)$ such that $p$ is homotopic to $p_1p_2$ and $\ell(p_1), \ell(p_2) < \ell(p)$. By minimality, both $p_1$ and $p_2$ are homotopically trivial, and so is $p$. □

We conclude this section by noting that the action of the hyperbolized group $\Gamma_X = \pi_1(X_\Gamma)$ on $\tilde{X}_\Gamma$ induces an action on the dual cubical complex $\mathcal{C}(\tilde{X}_\Gamma)$.

**Lemma 4.30.** The group $\Gamma_X = \pi_1(X_\Gamma)$ acts on $\mathcal{C}(\tilde{X}_\Gamma)$ by cubical isometries. Moreover, if $X$ is compact, then the action is cocompact.

*Proof.* The group $\Gamma_X$ acts on $\tilde{X}_\Gamma$ preserving the family of mirrors, hence the stratification defined in §3.5. The action permutes the cells, so $\Gamma_X$ acts on vertices of the dual cubical complex $\mathcal{C}(\tilde{X}_\Gamma)$ described in §4. Moreover, the action of $\Gamma_X$ on $\tilde{X}_\Gamma$
preserves the incidence relations between cells, hence we can extend the action of \( \Gamma_X \) on vertices to a combinatorial action of \( \Gamma_X \) on the entire \( \mathcal{C}(\tilde{X}_\Gamma) \). Since \( \Gamma_X \) acts on \( \mathcal{C}(\tilde{X}_\Gamma) \) combinatorially, it preserves the standard cubical metric.

When \( X \) is compact, \( X_\Gamma \) is compact as well, by (3) in Lemma 2.5. The action of \( \Gamma_X \) on \( \tilde{X}_\Gamma \) has finitely many orbits of cells, so its action on \( \mathcal{C}(\tilde{X}_\Gamma) \) has finitely many orbits of vertices. Since \( \mathcal{C}(\tilde{X}_\Gamma) \) is finite-dimensional (see Lemma 4.5), the quotient has finitely many cubes, so it is compact. \( \square \)

5. Special cubulation

The purpose of this section is to study the action of the hyperbolized group \( \Gamma_X = \pi_1(X_\Gamma) \) on the dual cubical complex \( \mathcal{C}(\tilde{X}_\Gamma) \) (see Lemma 4.30), and prove that the group \( \Gamma_X \) is virtually compact special in the sense of [HW08]. When \( X \) is compact and admissible, \( \Gamma_X \) is a Gromov hyperbolic group and \( \mathcal{C}(\tilde{X}_\Gamma) \) is a CAT(0) cubical complex (see (4) in Proposition 3.5, and Theorem 4.29). Therefore, one could hope to obtain virtual specialness directly from Agol’s result from [Ago13]. However, the action of \( \Gamma_X \) on \( \mathcal{C}(\tilde{X}_\Gamma) \) is not proper (see Remark 5.3). To remedy this, we will use a result of Groves and Manning from [GM18, Theorem D] that deals with improper actions. This requires a study of stabilizers of cubes.

In §5.1 we show that cube stabilizers for the action of \( \Gamma_X \) on \( \mathcal{C}(\tilde{X}_\Gamma) \) coincide with cell stabilizers for the action of \( \Gamma_X \) on \( \tilde{X}_\Gamma \). Then in §5.2 we show that such stabilizers are quasiconvex and virtually compact special. The complex \( X \) is always assumed to be admissible in the sense of §3. In some statements (such as Theorem 5.15) we also assume that it is compact.

Remark 5.1 (Why we consider the action on \( \mathcal{C}(\tilde{X}_\Gamma) \) instead of \( \tilde{X} \)). It is worth noting that when \( X \) is admissible, \( \tilde{X} \) is already a CAT(0) cube complex. Moreover the map \( g_X : X_\Gamma \to X \) from Proposition 3.5 induces a surjection \( \Gamma_X \to \pi_1(X) \) that can be used to obtain an action of \( \Gamma_X \) on \( \tilde{X} \). However, this action has a very large kernel. For example, in the case in which \( X \) is already simply connected the map \( \Gamma_X \to \pi_1(X) \) is trivial, but \( \Gamma_X \) is an infinite group; indeed, it retracts to \( \Gamma_{\square^n} \), as discussed in Remark 3.9.

5.1. Cube stabilizers for the action of \( \Gamma_X \) on \( \mathcal{C}(\tilde{X}_\Gamma) \). In this section we relate the cube stabilizers for the action of \( \Gamma_X \) on \( \mathcal{C}(\tilde{X}_\Gamma) \) to the cell stabilizers for the action of \( \Gamma_X \) on \( \tilde{X}_\Gamma \).

Lemma 5.2. The stabilizer of a vertex in \( \mathcal{C}(\tilde{X}_\Gamma) \) coincides with the stabilizer of its dual cell in \( \tilde{X}_\Gamma \).

Remark 5.3. It follows from Lemma 5.2 that the action of \( \Gamma_X \) on \( \mathcal{C}(\tilde{X}_\Gamma) \) is in general not proper. Namely, the stabilizer of a vertex dual to a cell of dimension at least 2 is infinite (compare Remark 4.8 and Figure 11).

We now proceed to the study of stabilizers of higher-dimensional cubes for the action of \( \Gamma_X \) on \( \mathcal{C}(\tilde{X}_\Gamma) \). Recall that the dual cubical complex \( \mathcal{C}(\tilde{X}_\Gamma) \) is equipped with a \( \Gamma_X \)-invariant height function: the vertex dual to a \( k \)-cell has height \( k \). We proved in Lemma 4.4 that every cube in \( \mathcal{C}(\tilde{X}_\Gamma) \) has a unique vertex of minimal height.
Lemma 5.4. The stabilizer of a cube in $\mathcal{C}(\tilde{X}_{\Gamma})$ coincides with the stabilizer of its vertex of minimal height in $\tilde{X}_{\Gamma}$.

Proof. Let $C$ be a cube in $\mathcal{C}(\tilde{X}_{\Gamma})$, let $v$ be its vertex of minimal height, and let $\sigma$ be the dual cell in $\tilde{X}_{\Gamma}$. Let $g \in \Gamma_X$ be an element that stabilizes $C$. Since the height function is invariant, $g$ must fix $v$, by uniqueness of the vertex of minimal height.

Conversely let $g$ fix $v$. By Lemma 5.2 we get that $g$ stabilizes $\sigma$, i.e. $g.\sigma = \sigma$. Let $w$ be another vertex of $C$ and let $\tau$ be the dual cell. By Lemma 4.4, we have that $\sigma \subseteq \tau$. It follows that $\sigma = g.\sigma \subseteq g.\tau$, so that both $\tau$ and $g.\tau$ appear in the link of $\sigma$ in the combinatorial structure of $\tilde{X}_{\Gamma}$ (see (3) in Lemma 3.17). Since the covering projection $\pi : \tilde{X}_{\Gamma} \to X_{\Gamma}$ induces isomorphisms on links, if $\tau$ and $g.\tau$ were distinct, then in $X_{\Gamma}$ we would see a tile $\pi(\tau) = \pi(g.\tau)$ isometric to a copy of $\Box^n_\mathbb{H}$ with some identification along the boundary (namely along the subspace corresponding to $\pi(\sigma)$). However, by (1) in Lemma 2.5 we know that tiles of $X_{\Gamma}$ are embedded copies of $\Box^n_\mathbb{H}$, so we must have $g.\tau = \tau$. By Lemma 5.2, this means $g.w = w$. Therefore $g$ fixes $C$ pointwise. □

Remark 5.5. In the proof of Lemma 5.4 we established that the stabilizer of a cell in $\mathcal{C}(\tilde{X}_{\Gamma})$ actually fixes the cell pointwise.

5.2. Cell stabilizers are quasiconvex and virtually compact special. The goal of this section is to study the stabilizers of cells for the action of $\Gamma_X$ on $\tilde{X}_{\Gamma}$ by covering transformations. In particular, note that by Lemma 3.11 stabilizers of tiles (i.e. top–dimensional cells) are isomorphic to the fundamental group of the hyperbolizing cube $\Gamma_{\Box^n} = \pi_1(\Box^n_\mathbb{H})$. More precisely, our goal is to show that cell stabilizers for the action of $\Gamma_X$ on $\tilde{X}_{\Gamma}$ are quasiconvex in $\Gamma_X$, and virtually compact special.

5.2.1. Quasiconvexity. In the following we say that an action of a group on a metric space is geometric if it is proper, cocompact and isometric. We will make use of the following standard fact.

Lemma 5.6. Let $Z$ be a proper Gromov hyperbolic geodesic space, and let $Y \subseteq Z$ be a quasiconvex subset. Let $G$ be a finitely generated group acting geometrically on $Z$, and let $H$ be the stabilizer of $Y$ in $G$. If $H$ acts cocompactly on $Y$, then $H$ is quasiconvex in $G$.

We apply this lemma to the cases $G = \Gamma_X$, $H = \Gamma_{\Box^n}$ and $G = \Gamma$, $H = \Gamma_{\Box^n}$. As noted, $\Gamma_X$ is a Gromov hyperbolic group when $X$ is compact. In both cases, before using the lemma we need to ensure that $H$ is a subgroup of $G$. This is not obvious, because a priori $H$ is just defined as the fundamental group of the hyperbolizing cube $\Box^n_\mathbb{H}$.

Lemma 5.7. Let $X$ be compact. Then $\Gamma_{\Box^n}$ is a quasiconvex subgroup of $\Gamma_X$.

Proof. By Lemma 2.5, we know that a hyperbolized complex retracts to each of its tiles, each of which is homeomorphic to the hyperbolizing cell. In our setting this means that $X_{\Gamma}$ retracts to $\Box^n_\mathbb{H}$, so in particular the inclusion $\Box^n_\mathbb{H} \hookrightarrow X_{\Gamma}$ as a tile induces an injection $\Gamma_{\Box^n} \hookrightarrow \Gamma_X$. Since $X$ is compact, the group $\Gamma_X$ acts geometrically on $\tilde{X}_{\Gamma}$. Moreover, the subgroup $\Gamma_{\Box^n}$ acts geometrically on a tile, which is a closed convex subspace by Lemma 3.16. Therefore $\Gamma_{\Box^n}$ is quasiconvex by Lemma 5.6. □
Lemma 5.8. The group $\Gamma_{\square^n}$ is a quasiconvex subgroup of $\Gamma$.

Proof. First of all we will prove that $\Gamma_{\square^n}$ naturally injects in $\Gamma$, by showing that there exists a (normal) cover $Y_T$ of $M_T = \mathbb{H}^n / \Gamma$ which retracts to $\square^n_T$ (see Figure 32). This would provide the desired injection

$$\Gamma_{\square^n} = \pi_1(\square^n_T) \hookrightarrow \pi_1(Y_T) \hookrightarrow \pi_1(M_T) = \Gamma.$$

To construct this cover, consider the admissible cubical complex $Y$ given by the standard cubulation of $\mathbb{R}^n$ with vertices on $\mathbb{Z}^n$. Notice that $Y$ admits a standard folding $f : Y \to \square^n$, and that $Y$ is an admissible cubical complex. Therefore we can consider the hyperbolized complex $Y'$. As in the proof of Lemma 5.7, Lemma 2.5 implies that $Y_T$ retracts onto any of its tiles, hence $\Gamma_{\square^n}$ injects in $\pi_1(Y_T)$.

We now claim that $Y_T$ is a (normal) covering space of $M_T$. For each $i = 1, \ldots, n$ consider the mirror of $Y$ given by $M_i = \{y_i = 0\}$, and the hyperplane of $Y$ given by $H_i = \{y_i = \frac{1}{2}\}$. Let $m_i$ and $h_i$ be the reflections in $M_i$ and $H_i$ respectively, i.e.

$$m_i : Y \to Y, m_i(y_1, \ldots, y_i, \ldots, y_n) = (y_1, \ldots, -y_i, \ldots, y_n)$$

$$h_i : Y \to Y, h_i(y_1, \ldots, y_i, \ldots, y_n) = (y_1, \ldots, 1 - y_i, \ldots, y_n)$$

For each $i = 1, \ldots, n$, the group $D_i = \langle m_i, h_i \rangle$ is an infinite dihedral group of cubical isometries of $Y$. The group $D = \langle m_1, h_1, \ldots, m_n, h_n \rangle$ is isomorphic to the direct product $D_1 \times \cdots \times D_n$, and admits a representation into the group $B_n$ of Euclidean isometries of the standard cube $\square^n$, in which $m_i$ acts trivially and $h_i$ acts as the standard reflection of $\square^n$ in the $i$-th coordinate. By Lemma 3.1 we have an action of $B_n$ on $\square^n_T$ by isometries, hence we can induce an action of $D$ on $\square^n_T$ by isometries such that the Charney–Davis map $g : \square^n_T \to \square^n$ is equivariant. Moreover, the standard folding $f : Y \to \square^n$ is clearly $D$–equivariant too, because it can be obtained by reflecting in the mirrors of $Y$. Since the two maps in the pullback square defining $Y_T$ are $D$–equivariant (see Figure 32), we obtain an action of $D$ on $Y_T$ by isometries.

Note that $t_i = h_i m_i$ is the unit integer translation of $Y$ in the $i$-th direction. As a result, $D$ contains a (normal) subgroup $T$ isomorphic to the group of integer translations $\mathbb{Z}^n$. The action of $D$ on $Y_T$ restricts to a free and properly discontinuous action of $T$ on $Y_T$. A fundamental domain for this action is given by a single tile. Each tile is isometric to a hyperbolizing cube $\square^n_T$, and the induced action identifies corresponding cells on opposite mirrors, recovering $M_T$ (see §3.1.1 for more details about the construction of $\square^n_T$.) In particular $\mu_T : Y_T \to Y_T / T \cong M_T$ realizes the desired covering map, which covers the standard universal covering map $\mu : Y = \mathbb{R}^n \to (S^1)^n$ (see Figure 32).

Finally, let us prove that $\Gamma_{\square^n}$ is quasiconvex in $\Gamma$. We know $\Gamma$ acts geometrically on $\mathbb{H}^n$, permuting the stratification induced by the coordinate mirrors and their translates. The subgroup $\Gamma_{\square^n}$ stabilizes a $\Gamma$–cell, i.e. the closure of a connected component of the complement of such collection. This is a closed convex subspace, on which $\Gamma_{\square^n}$ acts geometrically (see §3.3 for details). In particular $\Gamma_{\square^n}$ is quasiconvex in $\Gamma$ by Lemma 5.6. \qed

Remark 5.9. In Lemma 5.8 we have constructed a normal covering space of $M_T$ by producing an action of $T = \mathbb{Z}^n$ by deck transformations on the hyperbolization $Y_T$ of the standard integral cubulation of $\mathbb{R}^n$. This covering space can also be defined as the covering space of $M_T$ corresponding to the kernel of the homomorphism $\Gamma = \pi_1(Y_T)$.
5.2.2. Virtual specialness. A cubical complex is special if it admits a local isometry into the Salvetti complex of a right-angled Artin group (see [HW08]). A group $G$ is virtually compact special if there exist a finite index subgroup $G' \subseteq G$ and a compact special cubical complex $B$ such that $G' = \pi_1(B)$. This property passes from a Gromov hyperbolic group to its quasiconvex subgroups, as established in the following statement. This kind of arguments has appeared in the literature (see for instance [HW08, Corollary 7.8]). We include a proof for the reader’s convenience.

**Lemma 5.10.** Let $G$ be a Gromov hyperbolic group, and let $H$ be a quasiconvex subgroup. If $G$ is virtually compact special, then so is $H$.

**Proof.** Let $G'$ be a finite index subgroup of $G$ and $B$ a compact special cubical complex such that $G' = \pi_1(B)$. By [HW08, Remark 3.4, Lemma 3.13] we can assume without loss of generality that $B$ is also non–positively curved. The universal cover $\tilde{B}$ is a CAT(0) cubical complex. It is finite dimensional, uniformly locally finite, and Gromov hyperbolic, because $G'$ acts geometrically on it by covering transformations.

Let $H' = H \cap G'$. This is a finite–index subgroup of $H$, and a quasiconvex subgroup of $G'$. Since $G'$ is Gromov hyperbolic and acts geometrically on $B$, it follows that $H'$–orbits are quasiconvex. By [Hag08, Theorem H, or Corollary 2.29] or [SW15, Theorem 1.2], there exists a cocompact convex core for $H'$, i.e. a convex subcomplex $\tilde{Y} \subseteq \tilde{B}$ on which $H'$ acts cocompactly. Moreover, $H'$ acts by deck transformations, and the quotient $Y = \tilde{Y}/H'$ is a compact non–positively curved cubical complex with $\pi_1(Y) = H'$. The convex embedding $\tilde{Y} \hookrightarrow \tilde{B}$ descends to a local isometry $Y \rightarrow B$. Since $B$ is special, by [HW08, Corollary 3.9] we obtain that $Y$ is special too. Therefore $H$ is virtually compact special, as desired.

**Remark 5.11.** In the previous proof we have the Gromov hyperbolic group $H'$ acting geometrically on the CAT(0) cubical complex $\tilde{Y}$, so the fact that $H'$ is virtually compact special also follows from the celebrated theorem of Agol in [Ago13]. However here everything happens inside the special group $G'$, so one does not need Agol’s result.

\[ \begin{array}{cccc}
M_\Gamma & \xrightarrow{\mu_\Gamma} & Y_\Gamma & \xrightarrow{f_\Gamma} & \box^n \\
\downarrow g_0 & & \downarrow g_Y & & \downarrow g \\
(S^1)^n & \xleftarrow{\mu} & Y & \xrightarrow{f} & \box^n 
\end{array} \]

**Figure 32.** $Y_\Gamma$ as a covering space of $M_\Gamma$ that retracts to $\box^n$. 

$\pi_1(M_\Gamma) \rightarrow \mathbb{Z}^n$ induced by the collapse map $g_0 : M_\Gamma \rightarrow (S^1)^n$ obtained by applying the Pontryagin–Thom construction to $M_\Gamma$ with respect to suitable codimension–1 submanifolds (see §3.1.1 for details, and Figure 32). Compact quotients of $Y_\Gamma$, provide examples of closed hyperbolized manifolds which are finite covers of $M_\Gamma$. These are all genuine arithmetic hyperbolic manifolds.
We now apply the previous lemma to the cell stabilizers for the action of $\Gamma_{\square^n}$ on $\tilde{X}_G$, starting with the stabilizer of a tile.

**Lemma 5.12.** The group $\Gamma_{\square^n}$ is virtually compact special.

**Proof.** $\Gamma_{\square^n}$ is a quasiconvex subgroup of $\Gamma$ by Lemma 5.8 and $\Gamma$ is virtually compact special by [HW12, Theorem 1.6]. Indeed, it is an arithmetic lattice in $\text{SO}_0(n,1)$ by construction (see §3.1 or [CD95] for details). The statement then follows from Lemma 5.10. □

Finally we prove the same result for all cell stabilizers.

**Lemma 5.13.** Let $X$ be compact. Then the cell stabilizers for the action of $\Gamma_X$ on $\tilde{X}_G$ are quasiconvex and virtually compact special.

**Proof.** Let $\sigma$ be a cell in $\tilde{X}_G$ and let $H$ be the stabilizer of $\sigma$ for the action of $\Gamma_X$ on $\tilde{X}_G$. Since $\sigma$ is a convex subset of $\tilde{X}_G$ and $H$ acts geometrically on it, we conclude by Lemma 5.6 that $H$ is quasiconvex in $\Gamma_X$.

Arguing as in the proof of Lemma 5.4, if $\tau$ is a tile containing $\sigma$, and $K$ is its stabilizer, then $H \subseteq K$. Note that the folding map $X_G \to \square^n$ provides an isomorphism of $K \cong \Gamma_{\square^n}$, under which $H$ is isomorphic to a quasiconvex subgroup of $\Gamma_{\square^n}$ (again by Lemma 5.6). We know that $\Gamma_{\square^n}$ is Gromov hyperbolic (by Lemma 5.7 or Lemma 5.8) and virtually compact special (by Lemma 5.12). So it follows from Lemma 5.10 that $H$ is virtually compact special too. □

5.3. **Specialization.** We are now ready to prove that the fundamental group $\Gamma_X$ of the hyperbolized complex $X$ is virtually compact special, when the original cubical complex $X$ is admissible and compact. If the action of $\Gamma_X$ on $C(\tilde{X}_G)$ was proper, this would follow from Theorem 4.29, Lemma 4.30 and Agol’s main result from [Ago13]. However, as observed in Remark 5.3, the action on $C(\tilde{X}_G)$ is not proper. We will use a result by Groves and Manning (see Theorem D in [GM18]), which is designed to deal with this situation. We report here their statement for the reader’s convenience.

**Theorem 5.14.** [GM18, Theorem D] Suppose that $G$ is a Gromov hyperbolic group acting cocompactly on a CAT(0) cubical complex so that cell stabilizers are quasiconvex and virtually compact special. Then $G$ is virtually compact special.

Note that when authors of [GM18] say “virtually special” they imply that the quotient is compact (see page 3 in [GM18]). Also notice that they explicitly do not assume their complexes to be locally compact (see page 2).

**Theorem 5.15.** If $X$ is a compact admissible cubical complex and $\Gamma$ is a hyperbolizing lattice, then $\Gamma_X$ is virtually compact special Gromov hyperbolic group.

**Proof.** First of all, since $X$ is admissible, by Theorem 4.29 the dual cubical complex $C(\tilde{X}_G)$ is a CAT(0) cubical complex. Moreover, since $X$ is compact, $\Gamma_X$ is a Gromov hyperbolic group by (4) in Proposition 3.5. By Lemma 4.30 $\Gamma_X$ acts on $C(\tilde{X}_G)$ cocompactly by isometries.

Let $C$ be a cube in $C(\tilde{X}_G)$ and let $H$ be its stabilizer. By Lemma 5.4 $H$ coincides with the stabilizer of the vertex of minimal height in $C$. By Lemma 5.2 this in turn coincides with the stabilizer of the corresponding dual cell in $\tilde{X}_G$. Therefore by Lemma 5.13 $H$ is a quasiconvex subgroup of $\Gamma_X$ and it is also virtually compact.
special. Finally, by [GM18, Theorem D] the group $\Gamma_X$ is virtually compact special. □

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