# The Apollonian structure of integer superharmonic matrices 

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#### Abstract

We prove that the set of quadratic growths attainable by integer-valued superharmonic functions on the lattice $\mathbb{Z}^{2}$ has the structure of an Apollonian circle packing. This completely characterizes the PDE that determines the continuum scaling limit of the Abelian sandpile on the lattice $\mathbb{Z}^{2}$.


## 1. Introduction

1.1. Main results. This paper concerns the growth of integer-valued superharmonic functions on the lattice $\mathbb{Z}^{2}$; that is, functions $g: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ with the property that the value at each point $x$ is at least the average of the values at the four lattice neighbors $y \sim x$. In terms of the Laplacian operator, these are the functions $g$ satisfying

$$
\begin{equation*}
\Delta g(x):=\sum_{y \sim x}(g(y)-g(x)) \leq 0 \tag{1.1}
\end{equation*}
$$

for all $x \in \mathbb{Z}^{2}$. Our goal is to understand when a given quadratic growth at infinity, specified by a $2 \times 2$ real symmetric matrix $A$, can be attained by an integer-valued superharmonic function. For technical reasons, it is convenient to replace 0 by 1 in the inequality above. (It is straightforward to translate between the two versions using the function $f(x)=\frac{1}{2} x_{1}\left(x_{1}+1\right)$, which has $\Delta f \equiv 1$.) So we seek to determine, for each matrix $A$, whether there exists a function $g: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ such that

$$
\begin{equation*}
g(x)=\frac{1}{2} x^{t} A x+o\left(|x|^{2}\right) \quad \text { and } \quad \Delta g(x) \leq 1 \quad \text { for all } x \in \mathbb{Z}^{2} . \tag{1.2}
\end{equation*}
$$

When this holds, we say that the matrix $A$ is integer superharmonic, and we call $g$ an integer superharmonic representative for $A$.

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We will relate the set of integer superharmonic matrices to an Apollonian circle packing. Recall that every triple of pairwise tangent general circles (circles or lines) in the plane has exactly two Soddy general circles tangent to all three. An Apollonian circle packing is a minimal collection of general circles containing a given triple of pairwise-tangent general circles that is closed under the addition of Soddy general circles. Let $\mathcal{B}_{k}(k \in \mathbb{Z})$ denote the Apollonian circle packing generated by the vertical lines through the points $(2 k, 0)$ and $(2 k+2,0)$ in $\mathbb{R}^{2}$ together with the circle of radius 1 centered at $(2 k+1,0)$. The band packing is the union $\mathcal{B}=\bigcup_{k \in \mathbb{Z}} \mathcal{B}_{k}$, which is a circle packing of the whole plane, plus the vertical lines.

To each circle $C \in \mathcal{B}$ with center $\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}$ and radius $r>0$, we associate the matrix

$$
A_{C}=\frac{1}{2}\left[\begin{array}{cc}
r+x_{1} & x_{2}  \tag{1.3}\\
x_{2} & r-x_{1}
\end{array}\right] .
$$

We use the semi-definite order on the space $S_{2}$ of $2 \times 2$ real symmetric matrices, which sets $A \leq B$ if and only if $B-A$ is positive semidefinite. Our main result relates integer superharmonic matrices to the band packing.

Theorem 1.1. $A \in S_{2}$ is integer superharmonic if and only if either $A \leq A_{C}$ for some circle $C \in \mathcal{B}$, or $\operatorname{trace}(A) \leq 2$.

This theorem implies that the boundary of the set of integer superharmonic matrices looks like the surface displayed in Figure 1.1: it is a union of slope- 1 cones whose bases are the circles $C \in \mathcal{B}$ and whose peaks are the matrices $A_{C}$. Here we have identified $S_{2}$ with $\mathbb{R}^{3}$ with coordinates $\left(x_{1}, x_{2}, r\right)$,


Figure 1.1. One $2 \mathbb{Z}^{2}$-period of the boundary of the set of integer superharmonic matrices, as characterized by Theorem 1.1.
and $\mathcal{B}$ lies in the $r=0$ plane. Note that this embedding of $S_{2}$ with $\mathbb{R}^{3}$ has the property that the downset of a matrix $A$ is the slope- 1 downward cone in $\mathbb{R}^{3}$ whose vertex corresponds to $A$. In particular, the intersection of the cone

$$
\left\{A \in S_{2} \mid A \leq A_{C}\right\}
$$

with the $r=0$ plane is the closed disk bounded by $C$. Note that the condition $\operatorname{trace}(A) \leq 2$ in the theorem pulls in a Cantor set of matrices in the trace $(A)=$ 2 plane that are in the topological closure of the downset of the set of matrices $A_{C}(C \in \mathcal{B})$, but not in the downset itself.

To prove Theorem 1.1 we will recursively construct an integer superharmonic representative $g_{C}$ for each matrix $A_{C}(C \in \mathcal{B})$. Let

$$
\begin{equation*}
L_{C}=\left\{v \in \mathbb{Z}^{2} \mid A_{C} v \in \mathbb{Z}^{2}\right\} . \tag{1.4}
\end{equation*}
$$

By the extended Descartes circle theorem of Lagarias, Mallows and Wilks [14], each $A_{C}$ has rational entries, so $L_{C}$ is a full-rank sublattice of $\mathbb{Z}^{2}$.

The following theorem, from which Theorem 1.1 will follow, encapsulates the essential properties of the $g_{C}$ 's we construct.

Theorem 1.2. For each circle $C \in \mathcal{B}$, there exists an integer superharmonic representative $g_{C}$ for $A_{C}$ that satisfies the periodicity condition

$$
\begin{equation*}
g_{C}(x+v)=g_{C}(x)+x^{t} A_{C} v+g_{C}(v) \tag{1.5}
\end{equation*}
$$

for all $v \in L_{C}$ and $x \in \mathbb{Z}^{2}$. Moreover, $g_{C}$ is maximal in the sense that $g-g_{C}$ is bounded whenever $g: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ satisfies $\Delta g \leq 1$ and $g \geq g_{C}$.

We can always choose some $b \in \mathbb{R}^{2}$ so that $g\left(v_{i}\right)+b \cdot v_{i}=g\left(-v_{i}\right)-b \cdot v_{i}$ for two basis vectors $v_{1}, v_{2}$ of $L_{C}$. In particular, choosing $g(v-v)=g(0)=0$, this implies with (1.5) that $g_{C}(x)-\frac{1}{2} x^{t} A_{C} x-b^{t} x$ is $L_{C}$-periodic for some $b \in \mathbb{R}^{2}$. In particular, $g$ witnesses that $A_{C} \in \Gamma$. We call an integer superharmonic representative with this property an odometer for $A_{C}$. The construction of $g_{C}$ is explicit but rather elaborate. In Section 2 we outline the steps of this construction and give the derivation of Theorem 1.1 from Theorem 1.2. We now briefly survey a few connections to our work.
1.2. Hexagonal tilings of the plane by $90^{\circ}$ symmetric tiles. Figure 1.2 shows the Laplacians $\Delta g_{C}$ for a triple of circles in $\mathcal{B}$ and their two Soddy circles. In each case $\Delta g_{C}$ is periodic, and we have outlined a fundamental domain $T_{C}$ on whose boundary $\Delta g_{C}=1$. A major component of our paper is the construction of these $T_{C}$, which turn out to have a remarkable tiling property.

To state it precisely, for $x \in \mathbb{Z}^{2}$, write $s_{x}=\left\{x_{1}, x_{1}+1\right\} \times\left\{x_{2}, x_{2}+1\right\} \subseteq \mathbb{Z}^{2}$ and $\bar{s}_{x}=\left[x_{1}, x_{1}+1\right] \times\left[x_{2}, x_{2}+1\right] \subseteq \mathbb{R}^{2}$. If $T$ is a set of squares $s_{x}$, we call $T$


Figure 1.2. The periodic pattern $\Delta g_{C}$ of the integer superharmonic representative Theorem 1.2, shown for five different circles $C$. Black, patterned, and white cells correspond to sites $x \in \mathbb{Z}^{2}$ where $\Delta g_{C}(x)$ equals 1,0 , and -2 , respectively. (In general, -1 can also occur as a value.) In each case the fundamental tile $T_{C}$ is identified by a white outline. Clockwise from the top left is $\Delta g_{C}$ for the circle $(153,17,120)$, its three parents $(4,1,4),(9,1,6)$, and $(76,7,60)$, and its Soddy precursor $(25,1,20)$. The circles themselves are drawn in Figure 1.3.
a tile if the set

$$
\begin{equation*}
I(T):=\bigcup_{s_{x} \in T} \bar{s}_{x} \subseteq \mathbb{R}^{2} \tag{1.6}
\end{equation*}
$$

is a topological disk. A tiling of $\mathbb{Z}^{2}$ is a collection of tiles $\mathcal{T}$ such that every square $s_{x}\left(x \in \mathbb{Z}^{2}\right)$ belongs to exactly one tile from $\mathcal{T}$.


Figure 1.3. The child circle $C_{0}=(153,17,120)$ and its parents $C_{1}=(76,7,60), C_{2}=(4,1,4)$, and $C_{3}=(9,1,6)$ in $\mathcal{B}$. The other Soddy circle $C_{4}=(25,1,20) \in \mathcal{B}$ for the parents of $C_{0}$ is the Soddy precursor of $C_{0}$. The triples of integers are curvature coordinates, defined in Section 1.3.

Theorem 1.3. For every circle $C \in \mathcal{B}$, there is a tile $T_{C} \subseteq \mathbb{Z}^{2}$ with $90^{\circ}$ rotational symmetry, such that $T_{C}+L_{C}$ is a tiling of $\mathbb{Z}^{2}$. Moreover, except when $C$ has radius 1 , each tile in $T_{C}+L_{C}$ borders exactly 6 other tiles.

The tiles $T_{C}$ have the peculiar feature of being more symmetric than their plane tilings (which are only $180^{\circ}$ symmetric). We expect this to be a strong restriction. In particular, call a tiling regular if it has the form $T+L$ for some tile $T$ and lattice $L \subseteq \mathbb{Z}^{2}$, and hexagonal if each tile borders exactly six other tiles. For regular tilings $\mathcal{T}, \mathcal{T}^{\prime}$ of $\mathbb{Z}^{2}$, write $\mathcal{T}^{\prime} \prec \mathcal{T}$ if each tile in $T \in \mathcal{T}$ is a union of tiles from $\mathcal{T}^{\prime}$, and call the regular tiling $\mathcal{T}$ primitive if $\mathcal{T}^{\prime} \prec \mathcal{T}$ implies that either $\mathcal{T}^{\prime}=\mathcal{T}$, or $\mathcal{T}^{\prime}$ is the tiling of $\mathbb{Z}^{2}$ by squares $s_{x}$.

Conjecture 1.4. If $\mathcal{T}$ is a primitive, regular, hexagonal tiling of $\mathbb{Z}^{2}$ by $90^{\circ}$ symmetric tiles, then $\mathcal{T}=T_{C}+L_{C}+v$ for some $C \in \mathcal{B}$ and some $v \in \mathbb{Z}^{2}$.
1.3. Apollonian circle packings. Our proof of Theorem 1.2 is a recursive construction that mimics the recursive structure of Apollonian circle packings. Identifying a circle with center $\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}$ and radius $\frac{1}{c} \in \mathbb{R}$ with its curvature coordinates ( $c, c x_{1}, c x_{2}$ ) (some care must be taken in the case of lines), the Soddy circles $C_{0}$ and $C_{4}$ of a pairwise tangent triple of circles $C_{1}, C_{2}, C_{3}$ satisfy the linear equality

$$
\begin{equation*}
C_{0}+C_{4}=2\left(C_{1}+C_{2}+C_{3}\right) . \tag{1.7}
\end{equation*}
$$

This is a consequence of the Extended Descartes Theorem of Lagarias, Mallows, and Wilks [14] and can be used, for example, to prove that every circle in $\mathcal{B}$ has integer curvature coordinates. Pairwise tangent circles $C_{1}, C_{2}, C_{3}, C_{4}$ constitute a Descartes quadruple. Under permutation of indices, (1.7) gives


Figure 1.4. The tile of the circle $(153,17,120)$, decomposed into the tiles of its parents and Soddy precursor.
four different ways of producing a new Descartes quadruple sharing three circles in common with the original. These four transformations correspond to the four generators of the Apollonian group of Graham, Lagarias, Mallows, Wilks, and Yan [11], which acts on the Descartes quadruples of a circle packing. Our proof works by explicitly determining the action of the same Apollonian group on the set of maximal superharmonic representatives, by giving an operation on our family of integer superharmonic representatives analogous to the operation (1.7) for circles. In particular, referring to Figure 1.4, the example tile $T_{0}$ is seen to decompose into two copies each of $T_{1}, T_{2}, T_{3}$, with the copies of $T_{1}$ overlapping on a copy of $T_{4}$; note, for example, that the tile areas must therefore satisfy $\left|T_{0}\right|+\left|T_{4}\right|=2\left(\left|T_{1}\right|+\left|T_{2}\right|+\left|T_{3}\right|\right)$.

It follows from (1.7) that if the three generating general circles of an Apollonian circle packing have integer curvatures, then every general circle in the packing has integer curvature. In such a packing, the question of which integers arise as curvatures has attracted intense interest over the last decade [11], [22], [12], [2], [3]; see [10] for a survey.

Theorem 1.1 can be regarded as a new characterization of the circles appearing in the band packing: the circles in $\mathcal{B}$ correspond to integer superharmonic matrices that are maximal in the semidefinite order. In constructing $g_{C}$ and $T_{C}$ by an analogue of the Descartes rule (1.7), we follow ideas of Stange [23], who associates circles to "lax lattices" and proves a Descartes rule relating the bases of lattices corresponding to four mutually tangent general circles. Our proof also associates to each circle in $\mathcal{B}$ more detailed arithmetic information: Theorems 1.2 and 1.3 associate to each circle $C \in \mathcal{B}$ an integer superharmonic representative $g_{C}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ and a fundamental tile $T_{C}$. The curvature of $C$ can be recovered as the area of $T_{C}$.
1.4. Abelian sandpile. We briefly describe the Abelian sandpile model of Bak, Tang, and Wiesenfeld [1] that motivated our work. Put $n$ chips at the origin of $\mathbb{Z}^{2}$. In the sandpile model, a vertex having at least four chips topples by sending one chip to each lattice neighbor. Dhar [6] observed that the final configuration of chips does not depend on the order of topplings. This final
configuration $s_{n}$ displays impressive large-scale patterns [17], [8], [4], [21] and has been proposed by Dhar and Sadhu as a model of "proportionate growth" [7]: as $n$ increases, patterns inside $s_{n}$ scale up in proportion to the size of the whole. The second and third authors [20] have shown that as $n \rightarrow \infty$, the sandpile $s_{n}$ has a scaling limit on $\mathbb{R}^{d}$. The sandpile PDE that describes this limit depends on the set of integer superharmonic $d \times d$ matrices. In the companion paper [15] we use Theorem 1.1 to construct exact solutions to the sandpile PDE. An intriguing open problem is to describe the set of integer superharmonic matrices and analyze the associated sandpile processes for periodically embedded graphs other than $\mathbb{Z}^{2}$; see [19].

The weak-* limiting sandpile $s: \mathbb{R}^{2} \rightarrow \mathbb{R}$ appears to have the curious property that it is locally constant on an open neighborhood of the origin. Regions of constancy in $s$ correspond to regions of periodicity in $s_{n}$. Ostojic [17] proposed classifying which periodic patterns occur in $s_{n}$. Caracciolo, Paoletti, and Sportiello [4] and Paoletti [18] give an experimental protocol that recursively generates 2-dimensional periodic "backgrounds" and 1-dimensional periodic "strings." While this protocol makes no explicit reference to Apollonian circle packings, we believe that the 2-dimensional backgrounds it generates are precisely the Laplacians $\Delta g_{C}$ for $C \in \mathcal{B}$. Moreover, periodic regions in $s_{n}$ empirically correspond to the Laplacians $\Delta g_{C}$ of our odometers $g_{C}(C \in \mathcal{B})$, up to error sets of asymptotically negligible size.
1.5. Supplementary Materials. One view of this paper is as a proof that a certain recursive algorithm is correct. That is, we verify that a recursive construction produces odometers. We explicitly coded this algorithm in Matlab/Octave code and this is available on the second author's website and the arXiv. We have also included a larger version of the appendix in the arXiv submission.

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## 2. Overview of the proof

Here we outline the main elements of the proof of Theorem 1.2 and derive Theorem 1.1 from Theorem 1.2. The example in Figure 1.4 motivates the following recursive strategy for constructing the odometers $g_{C}$ :
(1) Construct fundamental domains (tiles) $T_{C}$ for each circle $C \in \mathcal{B}$ such that
(a) $T_{C}$ tiles the plane periodically; and
(b) $T_{C}$ decomposes into copies of smaller tiles $T_{C^{\prime}}$ with specified overlaps.
(2) Use the decomposition of $T_{C}$ to recursively define a tile odometer on $T_{C}$.
(3) Extend the tile odometer to $\mathbb{Z}^{2}$ via the periodicity condition (1.5), and check that the resulting extension $g_{C}$ satisfies $\Delta g_{C} \leq 1$ and is maximal in the sense stated in Theorem 1.2.

In Section 4, we carry out this strategy completely for two especially simple classes of circles.

In Section 5, for the general case of Step 1, we begin by associating to each pair of tangent circles $C, C^{\prime} \in \mathcal{B}$ a pair of Gaussian integers $v\left(C, C^{\prime}\right), a\left(C, C^{\prime}\right) \in$ $\mathbb{Z}[\mathbf{i}]$. The $v\left(C, C^{\prime}\right)$ 's will generate our tiling lattices, while the $a\left(C, C^{\prime}\right)$ 's will describe affine relationships among tile odometers. Recursive relations for these $v$ 's and $a$ 's, collected in Lemma 5.1, are used extensively throughout the paper.

In Section 7 we construct the tile $T_{C}$ recursively by gluing copies of tiles of the parent circles of $C$ as pictured in Figure 1.4, using the vectors $v\left(C, C^{\prime}\right)$ to specify the relative locations of the subtiles. The difficulty now lies in checking that the resulting set is in fact a topological disk and tiles the plane. This is proved in Lemma 7.3, which recursively establishes certain compatibility relations for tiles corresponding to tangent circles in $\mathcal{B}$. The argument uses a technical lemma, proved in Section 6, which allows one to prove that a collection of tiles forms a tiling from its local intersection properties.

Section 9 carries out Step 2 by constructing a tile odometer (an integervalued function on $T_{C}$ ) for each $C \in \mathcal{B}$. The recursive construction of the tile odometers mirrors that of the tiles (compare Definitions 7.1 and 9.3) using the additional data of the Gaussian integers $a\left(C, C^{\prime}\right)$, but we must check that it is well defined where subtiles intersect. A crucial difference between the tiles and the tile odometers is that the former have $90^{\circ}$ rotational symmetry, a fact that is exploited in the proof of the key Lemma 7.3. As tile odometers are merely $180^{\circ}$ symmetric, we require an extra ingredient to assert compatibility of some pairs of tile odometers. In particular, we show in Section 8 that tile boundaries can be suitably decomposed into smaller tiles, a fact that is used in the induction to verify some of the compatible tile odometer pairs.

The final step is carried out in Section 10. In Lemma 10.1 we use the compatibility relations among tile odometers to show that each tile odometer has a well-defined extension $g_{C}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ satisfying (1.5). We then show $\Delta g_{C} \leq 1$ by giving an explicit description of the values of the Laplacian $\Delta g_{C}$ on subtile intersections (Lemma 10.3). In Lemma 10.5, we use this same explicit description of $\Delta g_{C}$ to prove maximality (The fact that $\Delta g_{C} \equiv 1$ on the boundary of $T_{C}$ will be crucial here.) The proof of Theorem 1.2 is completed by showing that the lattice $L_{C}$ of (1.4) is the same as the lattice $\Lambda_{C}$ generated by the vectors $v\left(C, C_{i}\right)$, where $C_{1}, C_{2}, C_{3}$ are the parent circles of $C$. Strangely, our proof of this elementary statement depends on everything else in this paper, as it uses the structure of the constructed odometers. It seems reasonable to expect that a self-contained proof that $L_{C}=\Lambda_{C}$ can be found.

## 3. Integer functions on the lattice

In this section we establish properties of integer functions on $\mathbb{Z}^{2}$. In particular, we prove that $\Gamma$ is both downward and topologically closed and derive Theorem 1.1 from Theorem 1.2. Two key properties of the set of integer superharmonic matrices are used for this derivation: that the set is downward closed, and that the set is topologically closed. To begin, we recall some basic properties of the Laplacian.

Proposition 3.1. The Laplacian agrees with its continuum analogue on quadratic polynomials, is linear, and it monotone. That is, the following properties hold:
(1) If $A \in \mathbb{R}_{\text {sym }}^{2 \times 2}, b \in \mathbb{R}^{2}, c \in \mathbb{R}$ and $u(x)=\frac{1}{2} x \cdot A x$, then $\Delta u \equiv \operatorname{trace} A$.
(2) $\Delta(u+v)=\Delta u+\Delta v$.
(3) $\Delta \max \{u, v\} \leq \max \{\Delta u, \Delta v\}$, where $\max$ is point-wise.

In particular, this proposition tells us the adding and affine function $u(x)=b \cdot x+c$ does not change the Laplacian. We next prove that quadratic polynomials with nonnegative definite Hessians can be approximated by integer-valued functions with nonnegative Laplacians.

Lemma 3.2. If $A \in S_{2}$ satisfies $A \geq \mathbf{0}$, then there exists $g: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ such that

$$
\Delta g(x) \geq 0 \quad \text { and } \quad g(x)=\frac{1}{2} x^{t} A x+O(1+|x|)
$$

for all $x \in \mathbb{Z}^{2}$.
Together with the linearity of the Laplacian, this lemma establishes that the set of integer superharmonic matrices is downward closed in the semidefinite order by giving a suitable integer approximation of the difference $x \mapsto \frac{1}{2} x^{t}(B-A) x$.

Proof. Define $q(x)=\frac{1}{2} x^{t} A x$ and

$$
\begin{equation*}
g(x)=\sup _{p \in \mathbb{Z}^{2}} \inf _{y \in \mathbb{Z}^{2}}|q(y)+|y|\rfloor+p \cdot(x-y) . \tag{3.1}
\end{equation*}
$$

If we set $p=\lfloor D q(x)\rfloor$ in the supremum, then we obtain

$$
\begin{aligned}
g(x) & \geq \inf _{y \in \mathbb{Z}^{2}} q(y)+|y|+\lfloor D q(x)\rfloor \cdot(x-y)-1 \\
& \geq \inf _{y \in \mathbb{Z}^{2}} q(y)+D q(x) \cdot(x-y)-|x|-1 \\
& \geq q(x)-|x|-1,
\end{aligned}
$$

where we used the triangle inequality in the second step and the convexity of $q$ in the third. If we set $y=x$ in the infimum, then we obtain

$$
g(x) \leq \sup _{p \in \mathbb{Z}^{2}} q(x)+|x|+1=q(x)+|x|+1 .
$$

Finally, since $g$ is a pointwise supremum of affine functions (since the $y$ value realizing the infemum in (3.1) is independent of $x$ ), we obtain $\Delta g \geq 0$ by the monotonicity of $\Delta$.

For topological closure, we will use the following consequence of a discrete Harnack inequality [13]:

Proposition 3.3. There is a $C>0$ such that, if $u: \mathbb{Z}^{2} \rightarrow \mathbb{R}, u(x) \geq 0$, and $\Delta u(x) \leq 1$, then $u(x) \leq C\left(u(0)+|x|^{2}\right)$.

We can now prove
Lemma 3.4. The set of integer superharmonic matrices is topologically closed.

Proof. Suppose $A_{k} \in \Gamma$ and $\lim _{k \rightarrow \infty} A_{k}=A$. We must show $A \in \Gamma$.
Step 1 . We find $u_{k}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ and $b_{k} \in \mathbb{R}^{2}$ such that

$$
\Delta u_{k}(x) \leq 1, \quad u_{k}(x) \geq \frac{1}{2} x \cdot A_{k} x+b_{k} \cdot x, \quad u_{k}(0) \leq 1, \quad\left|b_{k}\right| \leq 1 .
$$

Replacing $A_{k}$ with $A_{k}-\frac{1}{k} I$ and adding a constant to each $u_{k}$, we may assume that $u_{k}(x) \geq \frac{1}{2} x \cdot A_{k} x$. In fact, we may assume that $0 \leq u_{k}\left(x_{k}\right)-\frac{1}{2} x_{k} \cdot A x_{k} \leq 1$ holds for some $x_{k} \in \mathbb{Z}^{2}$. Replacing $u$ by $x \mapsto u\left(x+x_{k}\right)-\frac{1}{2} x_{k} \cdot A x_{k}$ and setting $b_{k}=A_{k} x_{k}$, we obtain the first three properties. Subtracting from $u_{k}$ the integer harmonic linear function $x \mapsto \tilde{b}_{k} x$, for some integer rounding $\tilde{b}_{k}$ of $b_{k}$, we may replace $b_{k}$ by $b_{k}-\tilde{b}_{k}$ to obtain the fourth property.

Step 2. We conclude by showing the $u_{k}$ are precompact. Since $\left|A_{k}\right|$ is bounded independently of $k$, the Harnack inequality implies that

$$
v_{k}(x)=u_{k}(x)-\frac{1}{2} x \cdot A_{k} x-b_{k} \cdot x
$$

satisfies

$$
\left|v_{k}(x)\right| \leq C\left(1+|x|^{2}\right)
$$

Combined with the boundedness of $\left|A_{k}\right|$ and $\left|b_{k}\right|$, we obtain

$$
\left|u_{k}(x)\right| \leq C\left(1+|x|^{2}\right) .
$$

In particular, for fixed $x \in \mathbb{Z}^{2}$, there are only finitely many possibilities for $u_{k}(x)$. By compactness, we may assume that $u_{k} \rightarrow u$ pointwise as $k \rightarrow \infty$. Since the $u_{k}$ are integer valued, this means that $k \mapsto u_{k}$ is eventually constant on any finite set $X \subseteq \mathbb{Z}^{2}$. Since $\Delta u(x)$ only depends on nearest neighbors, we obtain in the limit

$$
\Delta u(x) \leq 1 .
$$

By compactness, we may assume $b_{k} \rightarrow b$ and obtain in the limit

$$
u(x) \geq \frac{1}{2} x \cdot A x+b \cdot x .
$$

It follows that $A \in \Gamma$.

We are now ready to derive Theorem 1.1 from Theorem 1.2.
Proof of Theorem 1.1. The downward closure of integer superharmonic matrices now makes the implication $\left(A \leq A_{C} \Longrightarrow A\right.$ superharmonic $)$ from Theorem 1.1 an immediate consequence of Theorem 1.2. The implication ( $\operatorname{trace}(A) \leq 2 \Longrightarrow A$ superharmonic) follows from Lemma 3.4, since the circle packing $\mathcal{B}$ is dense in the points exterior to its circles.

For the "only if" direction, observe that points of tangency are dense on each circle of the band packing $\mathcal{B}$. In particular, any circle in the plane is either enclosed by some circle of $\mathcal{B}$ or strictly encloses some circle of $\mathcal{B}$. Therefore any matrix $A \in S_{2}$ with $\operatorname{trace}(A)>2$ satisfies either $A \leq A_{C}$ for some $C \in \mathcal{B}$ or $A_{C}<A$ (that is, $A-A_{C}$ is positive definite) for some $C \in \mathcal{B}$. In the latter case, the existence of an integer superharmonic representative $g$ for $A$ would contradict the maximality of $g_{C}$ : by (1.2), some constant offset $g+c$ of $g$ would satisfy $g+c \geq g_{C}$, but $A_{C}<A$ implies that $g+c-g_{C}$ is unbounded.

## 4. Two degenerate cases

4.1. Explicit constructions. Before defining odometers for general circles, we consider two subfamilies of $\mathcal{B}$, the Ford and Diamond circles. The odometers for these subfamilies have simpler structure than the general case, and this allows us to give a more explicit description of their construction. The purpose of carrying this out is twofold. First, these cases, particularly the Ford circles, serve as a concrete illustration of our general construction that avoids most of the technical complexity. Second, by taking these subfamilies as our base-case in the general construction, we avoid several tedious degeneracies. In particular, the general construction operates quite smoothly when the circle under consideration has parents with distinct nonzero curvatures; this condition fails precisely when the circle is either Ford or Diamond.
4.2. The Ford circles. It is a classical fact that the reduced fractions, namely, the pairs of integers $(p, q) \in \mathbb{Z}^{2}$ such that $q \geq 1$ and $\operatorname{gcd}(q, p)=1$, are in algebraic bijection with the circles in $\mathcal{B}_{0}$ that are tangent to the vertical line through the origin. In curvature coordinates, this bijection is given by

$$
(p, q) \mapsto\left(q^{2}, 1,2 p q\right)
$$

and the circles are called the Ford circles. We write $C_{p q}$ for the Ford circle associated to the reduced fraction $(p, q)$.

The tangency structure of the Ford circles is famously simple. If $q \geq 2$, then the parents of the Ford circle $C_{p q}$ in the packing $\mathcal{B}_{0}$ (Figure 5.1) are the vertical line through the origin, and the unique Ford parents $C_{p q}^{1}:=C_{p_{1} q_{1}}$ and $C_{p q}^{2}:=C_{p_{2} q_{2}}$ determined by the constraints $p_{1} q-q_{1} p=1, p_{2} q-q_{2} p=-1$, $0 \leq q_{1}<q$ and $0<q_{2} \leq q$. The Ford parents are the only Ford circles tangent to $C_{p q}$ having smaller curvature (except that when $q=1$, we have $q_{1}=0$ and


Figure 4.1. The Ford circles $C_{1,3}, C_{2,5}$, and $C_{3,8}$ and several periods of the Laplacian of their odometers $g_{p q}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$.
thus the Ford parents are not both Ford circles). We observe that the definition of parents implies $(p, q)=\left(p_{1}, q_{1}\right)+\left(p_{2}, q_{2}\right)$ and $p_{1} / q_{1}>p / q>p_{2} / q_{2}$, which gives the classical relationship between the Ford circles and the Farey fractions. This connection between the circles and Farey fractions was used by Ford [9] (and later, by Nichols [16]) to prove results about Diophantine approximation.

The odometer $g_{p q}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ associated to the Ford circle $C_{p q}$ also enjoys a simple description. To understand the structure of $g_{p q}$, we consider the Ford circle $C_{3,8}$ and its parents $C_{2,5}$ and $C_{1,3}$. These circles have the tangency structure and periodic odometer Laplacians $\Delta g_{p q}$ displayed in Figure 4.1.

Examining the patterns in Figure 4.1, we see that $[1, q]^{2}$ is a fundamental domain for the periodic Laplacian $\Delta g_{p q}$. Moreover, the fundamental tile of $\Delta g_{3,8}$ decomposes (with a few errors on points of overlap) into two copies each of the fundamental tiles of $\Delta g_{1,3}$ and $\Delta g_{2,5}$.

These observations lead to the following construction. For a general Ford circle $C_{p q}$, we define an odometer $g_{p q}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ by specifying its values on the domain $[0, q]^{2} \subseteq \mathbb{Z}^{2}$ and then extending periodically. The values of $g_{p q}$ on $[0, q]^{2}$ are determined recursively by copying data from the parent odometers $g_{p_{1} q_{1}}$ and $g_{p_{2} q_{2}}$ onto the subdomains $E_{1}:=\left[0, q_{1}\right]^{2}, E_{2}:=\left[q_{2}, q\right]^{2}, E_{3}:=\left[q_{1}, q\right] \times$ $\left[0, q_{2}\right]$, and $E_{4}:=\left[0, q_{2}\right] \times\left[q_{1}, q\right]$. These subdomains always have one of the three overlapping structures displayed in Figure 4.2, depending on the relative sizes of $q_{1}$ and $q_{2}$. This construction is encoded precisely in the following lemma. Note that our conditions below are redundant, since one only needs two translations to generate a lattice.

Lemma 4.1. There is a unique family of functions $g_{p q}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$, indexed by the Ford circles $C_{p q}$, satisfying
$g_{p q}(x)=\left\lceil\frac{p}{q} x_{1} x_{2}\right\rceil$
for $x \in[0, q]^{2} \backslash[2, q-2]^{2}$,
$g_{p q}(x+(0,-q))=g_{p q}(x)+(-p, 0) \cdot x$
for $x \in \mathbb{Z}^{2}$,

|  |  |
| :---: | :---: |
| $E_{4}$ | $E_{2}$ |
|  |  |
|  |  |
| $E_{1}$ | $E_{3}$ |
|  |  |



Figure 4.2. The three possible overlapping structures of the subtiles $E_{i}$ of a Ford tile $[0, q]^{2}$. (We are drawing outlines in the dual lattice; so for example, $E_{1} \cap E_{3}$ is a just a line of points in $\mathbb{Z}^{2}$.)
$g_{p q}\left(x+\left(q, q_{1}\right)\right)=g_{p q}(x)+\left(p_{1}, p\right) \cdot x+q_{1} p \quad$ for $x \in \mathbb{Z}^{2}$,
$g_{p q}\left(x+\left(-q, q_{2}\right)\right)=g_{p q}(x)+\left(p_{2},-p\right) \cdot x-q p_{2} \quad$ for $x \in \mathbb{Z}^{2}$
and, when $q \geq 2$,
$g_{p q}(x)=g_{p_{1} q_{1}}(x) \quad$ for $x \in E_{1}$,
$g_{p q}(x)=g_{p_{1} q_{1}}\left(x-\left(q_{2}, q_{2}\right)\right)+\left(p_{2}, p_{2}\right) \cdot x-p_{2} q_{2}+1 \quad$ for $x \in E_{2}$,
$g_{p q}(x)=g_{p_{2} q_{2}}\left(x-\left(q_{1}, 0\right)\right)+\left(0, p_{1}\right) \cdot x \quad$ for $x \in E_{3}$,
$g_{p q}(x)=g_{p_{2} q_{2}}\left(x-\left(0, q_{1}\right)\right)+\left(p_{1}, 0\right) \cdot x \quad$ for $x \in E_{4}$,
where $E_{1}:=\left[0, q_{1}\right]^{2}, E_{2}:=\left[q_{2}, q\right]^{2}, E_{3}:=\left[q_{1}, q\right] \times\left[0, q_{2}\right]$, and $E_{4}:=\left[0, q_{2}\right] \times$ $\left[q_{1}, q\right]$.

Proof. We prove by induction on $q \geq 1$ that there is a function $g_{p q}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ satisfying the above conditions. For $p \in \mathbb{Z}$, we observe that

$$
g_{p, 1}(x):=\frac{1}{2} x_{1}\left(x_{1}-1\right)+p x_{1} x_{2}
$$

satisfies the first four rules. Thus, we may assume that $q \geq 2$ and the induction hypothesis holds when $q^{\prime}<q$.

We first observe that any map $g_{p q}:[0, q]^{2} \rightarrow \mathbb{Z}$ that satisfies the first condition has a unique extension to all of $\mathbb{Z}^{2}$ that satisfies the next three conditions. Since the translations of $[0, q]^{2}$ by the lattice $L_{p q}$ generated by $\left\{(0,-q),\left(q, q_{1}\right),\left(-q, q_{2}\right)\right\}$ cover $\mathbb{Z}^{2}$, the three translation conditions prescribe values for $g_{p q}$ on the rest of $\mathbb{Z}^{2}$. Since $[0, q-1]^{2}+L_{p q}$ is a partition of $\mathbb{Z}^{2}$ and the composition of the three translations is the identity, it suffices to check the consistency of the translation conditions on the set $[0, q]^{2} \backslash[1, q-1]^{2}$. This follows from the first condition.

Thus, to construct $g_{p q}$, it suffices to check the consistency of the first and the last four conditions, which specify the values of $g_{p q}$ on $[0, q]^{2}$. We observe that the first condition prescribes values on the doubled boundary set
$B:=[0, q]^{2} \backslash[2, q-2]^{2}$ and the last four conditions prescribe values on the sets $E_{i}$ described above. We must check the consistency of these five prescriptions.

Case 1: consistency for the intersections $E_{1} \cap E_{3}, E_{1} \cap E_{4}, E_{2} \cap E_{3}$, and $E_{2} \cap E_{4}$. To check the first intersection, we must verify that

$$
g_{p_{1} q_{1}}(x)=g_{p_{2} q_{2}}\left(x-\left(q_{1}, 0\right)\right)+\left(0, p_{1}\right) \cdot x
$$

for $x \in E_{1} \cap E_{3}=\left\{q_{1}\right\} \cap\left[0, \min \left\{q_{1}, q_{2}\right\}\right]$. Applying the inductive version of the first condition, this reduces to

$$
\frac{p_{1}}{q_{1}} x_{1} x_{2}=\frac{p_{2}}{q_{2}}\left(x_{1}-q_{1}\right) x_{2}+p_{1} x_{2} .
$$

Since $x_{1}=q_{1}$, this is easily seen to be true. The other three intersections can by checked by symmetric arguments.

Case 2: consistency for the intersections $E_{1} \cap E_{2}$ and $E_{3} \cap E_{4}$. To check the first intersection, we must verify that

$$
g_{p_{1} q_{1}}(x)=g_{p_{1} q_{1}}\left(x-\left(q_{2}, q_{2}\right)\right)+\left(p_{2}, p_{2}\right) \cdot x-p_{2} q_{2}+1
$$

for all $x \in E_{1} \cap E_{2}=\left[q_{2}, q_{1}\right]^{2}$. This is nontrivial if and only if $q_{1} \geq q_{2}$, in which case $C_{p_{2} q_{2}}$ is a parent of $C_{p_{1} q_{1}}$. Let $C_{p_{3} q_{3}}$ denote the other parent. Since $p_{3} / q_{3}>p_{1} / q_{1}>p_{2} / q_{2}$, the inductive version of the last four conditions gives

$$
g_{p_{1} q_{1}}(x)=g_{p_{3} q_{3}}\left(x-\left(q_{2}, q_{2}\right)\right)+\left(p_{2}, p_{2}\right) \cdot x-p_{2} q_{2}+1
$$

and, since $q_{3}=q_{1}-q_{2}$,

$$
g_{p_{1} q_{1}}\left(x-\left(q_{2}, q_{2}\right)\right)=g_{p_{3} q_{3}}\left(x-\left(q_{2}, q_{2}\right)\right)
$$

for all $x \in\left[q_{2}, q_{1}\right]^{2}$. Substituting this into the above equation, we easily see that equality holds. The other intersection can be checked by a symmetric argument.

Case 3: consistency for the intersections $B \cap E_{i}$. When $i=1$, this amounts to verifying

$$
g_{p_{1} q_{1}}(x)=\left\lceil\frac{p}{q} x_{1} x_{2}\right\rceil
$$

for $x \in\left[0, q_{1}\right] \times[0,1]$. By induction, this reduces to checking

$$
\left\lceil\frac{p_{1}}{q_{1}} x_{1} x_{2}\right\rceil=\left\lceil\frac{p}{q} x_{1} x_{2}\right\rceil
$$

for $x \in\left[0, q_{1}\right] \times[0,1]$. When $x_{1} x_{2}=0$, this is trivial, so we may assume $x_{1}>0$ and $x_{2}=1$. Since $\operatorname{gcd}(p, q)=1$, we see that the distance between $\frac{p}{q} x_{1}$ and the nearest integer is at least $\frac{1}{q}$. Using the relation $q p_{1}-p q_{1}=1$, we see that $\left|\frac{p}{q} x_{1}-\frac{p_{1}}{q_{1}} x_{1}\right|=\frac{1}{q q_{1}}\left|x_{1}\right| \leq \frac{1}{q}$ and therefore $\left\lceil\frac{p_{1}}{q_{1}} x_{1} x_{2}\right\rceil=\left\lceil\frac{p}{q} x_{1} x_{2}\right\rceil$, as desired. The other intersections can be checked by symmetric arguments.

In contrast to our construction of general circle odometers in later sections, the description of the Ford circle odometers $g_{p q}$ in Lemma 4.1 has two enormous advantages: the fundamental tiles are square shaped and, more importantly, there is a closed formula for the odometer on the outer two layers of the fundamental tile. This renders geometric questions about the fundamental tile moot and makes it possible to compute the Laplacian $\Delta g_{p q}$ in a relatively straightforward manner.

Define the boundary of a subset $X \subseteq \mathbb{Z}^{2}$ to be

$$
\partial X=\left\{x \in X:|x-y|<2 \text { for some } y \in \mathbb{Z}^{2} \backslash X\right\}
$$

and call $X \backslash \partial X$ the interior of $X$. Thus, $\partial X$ consists of the points in $X$ whose distance to $\mathbb{Z}^{2} \backslash X$ is 1 or $\sqrt{2}$. In the case of a square $[0, q]^{2}$, we have $\partial[0, q]^{2}=[0, q]^{2} \backslash[1, q-1]^{2}$.

Lemma 4.2. The Laplacian $\Delta g_{p q}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ satisfies

$$
\Delta g_{p q}(x)=\Delta g_{p q}(x+(0, q))=\Delta g_{p q}\left(x+\left(q, q_{1}\right)\right) \quad \text { for } x \in \mathbb{Z}^{2}
$$

and

$$
\Delta g_{p q}(x)=1 \quad \text { for } x \in \partial[0, q]^{2} .
$$

Moreover, if $q \geq 2$ and $x \in[1, q-1]^{2}$ lies in 2,3 , or 4 of the boundaries $\partial E_{i}$, then $\Delta g_{p q}(x)=1,0$, or -2 , respectively.

Proof. To obtain the periodicity conditions, we simply compute the Laplacian of the periodicity conditions in Lemma 4.1, recalling that the Laplacian of an affine function vanishes.

We next check $\Delta g_{p q}(x)=1$ for $x \in[0, q]^{2} \backslash[1, q-1]^{2}$. We claim that, for $0 \leq a \leq q$,

$$
\begin{gathered}
g_{p q}(a, 0)=g_{p q}(0, a)=0, \\
g_{p q}(a, 1)=g_{p q}(1, a)=\left\lceil\frac{p}{q} a\right\rceil,
\end{gathered}
$$

and

$$
g_{p q}(a,-1)=g_{p q}(-1, a)=-\left\lfloor\frac{p}{q} a\right\rfloor .
$$

The first two are immediate from Lemma 4.1, so we check the third. Using the periodicity, we compute

$$
g_{p q}(a,-1)=g_{p q}(a, q-1)-(p, 0) \cdot(a,-1)=\left\lceil\frac{p}{q} a(q-1)\right\rceil-p a=-\left\lfloor\frac{p}{q} a\right\rfloor .
$$

Similarly, we compute

$$
\begin{aligned}
g_{p q}(-1, a) & =g_{p q}\left(q-1, a+q_{1}\right)-\left(p_{1}, p\right) \cdot(-1, a)-q_{1} p \\
& =\left\lceil\frac{p}{q}(q-1)\left(a+q_{1}\right)\right\rceil+p_{1}-p a-q_{1} p \\
& =-\left\lfloor\frac{p}{q}\left(a+q_{1}\right)\right\rfloor+p_{1} \\
& =-\left\lfloor\frac{p}{q} a-\frac{1}{q}\right\rfloor \\
& =-\left\lfloor\frac{p}{q} a\right\rfloor
\end{aligned}
$$

for $0 \leq a \leq q_{2}$ and

$$
\begin{aligned}
g_{p q}(-1, a) & =g_{p q}\left(q-1, a-q_{2}\right)-\left(-p_{2}, p\right) \cdot(-1, a)+p_{2} q \\
& =\left\lceil\frac{p}{q}(q-1)\left(q-q_{2}\right)\right\rceil+p_{2}-p a+p_{2} q \\
& =-\left\lfloor\frac{p}{q}\left(a-q_{2}\right)\right\rfloor+p_{2} \\
& =-\left\lfloor\frac{p}{q} a-\frac{1}{q}\right\rfloor \\
& =-\left\lfloor\frac{p}{q} a\right\rfloor
\end{aligned}
$$

for $q_{2} \leq a \leq q$. Using $\operatorname{gcd}(p, q)=1$, we then obtain

$$
\Delta g_{p q}(0, a)=\Delta g_{p q}(a, 0)=\left\lceil\frac{p}{q} a\right\rceil-\left\lfloor\frac{p}{q} a\right\rfloor=1
$$

for $0 \leq a<q$. That $\Delta g_{p q}=1$ on the rest of $[0, q]^{2} \backslash[1, q-1]^{2}$ follows by periodicity.

The moreover clause can be handled similarly, using Lemma 4.1 to explicitly evaluate the Laplacian at the intersections of the boundaries of the sets. For example, suppose that $x \in\left(E_{1} \cap E_{3}\right) \backslash\left(E_{2} \cup E_{4}\right)$. In this case, $x=\left(q_{1}, h\right)$ for some $0<h<\min \left\{q_{1}, q_{2}\right\}$. Using $g_{p q}=g_{p_{1} q_{1}}$ on $E_{1}$, we compute

$$
g_{p q}(x+(0, k))=\left\lceil\frac{p_{1}}{q_{1}} q_{1}(h+k)\right\rceil=p_{1}(h+k) \quad \text { for } k=-1,0,1
$$

and

$$
g_{p q}(x-(1,0))=\left\lceil\frac{p_{1}}{q_{1}}\left(q_{1}-1\right) h\right\rceil=p_{1} h-\left\lfloor\frac{p_{1}}{q_{1}} h\right\rfloor .
$$

Using $q_{p q}(y)=g_{p_{2} q_{2}}\left(y-\left(q_{1}, 0\right)\right)+\left(0, p_{1}\right) \cdot x$ for $y \in E_{3}$, we compute

$$
\begin{aligned}
g_{p q}(x+(1,0)) & =g_{p_{2} q_{2}}(1, h)+p_{1} h \\
& =\left\lceil\frac{p_{2}}{q_{2}} h\right\rceil+p_{1} h \\
& =\left\lceil\frac{p_{1}}{q_{1}} h-\frac{1}{q_{1} q_{2}} h\right\rceil+p_{1} h \\
& =\left\lceil\frac{p_{1}}{q_{1}} h\right\rceil+p_{1} h .
\end{aligned}
$$

Putting these together, we obtain $\Delta g_{p q}(x)=1$.
The above lemma tells us how to compute $\Delta g_{p q}$. Indeed, suppose we wish to compute $\Delta g_{p q}(x)$ for some $x \in \mathbb{Z}^{2}$. We first reduce to the case $x \in[0, q]^{2}$ using the periodicity of $\Delta g_{p q}$. Now, if $x$ lies on the boundary of $[0, q]^{2}$ or least two of the boundaries of the $E_{i}$, then we can read off $\Delta g_{p q}(x)$ from Lemma 4.2. Otherwise, $x$ lies in the interior of one of the $E_{i}$, and we can proceed recursively to the corresponding parent $\Delta g_{p_{i} q_{i}}(x)=\Delta g_{p q}(x)$.

Following this line of reasoning, we prove Theorem 1.2 in the special case of the Ford circles. Observe that the peak associated to the Ford circle $(p, q)$ is the matrix

$$
A_{p q}:=\frac{1}{q^{2}}\left[\begin{array}{cc}
1 & p q \\
p q & 0
\end{array}\right],
$$

which has lattice $L_{p q}:=\left\{x \in \mathbb{Z}^{2}: A_{p q} x \in \mathbb{Z}^{2}\right\}$ generated by $\left\{(0,-q),\left(q, q_{1}\right)\right\}$.

Call $X \subset \mathbb{Z}^{2}$ connected if for all $x, y \in X$ there is a sequence $x=$ $x_{0}, \ldots, x_{k}=y$ with $x_{i} \in X$ and $\left|x_{i}-x_{i-1}\right|=1$ for all $i=1, \ldots, k$. Call a bounded set $X$ simply connected if $\mathbb{Z}^{2} \backslash X$ is connected.

Proposition 4.3. For every Ford circle $C_{p q}$, the odometer $g_{p q}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ satisfies

$$
\Delta g_{p q}(x) \leq 1 \quad \text { and } \quad g_{p q}(x+v)=g_{p q}(x)+x \cdot A_{p q} v+g_{p q}(v)
$$

for all $x \in \mathbb{Z}^{2}$ and $v \in L_{p q}$. Moreover, the only infinite connected subset $X \subseteq \mathbb{Z}^{2}$ such that $\Delta\left(g_{p q}+\mathbf{1}_{X}\right) \leq 1$ is $X=\mathbb{Z}^{2}$.

Proof. The periodicity condition is immediate from Lemma 4.1, after we observe $A_{p q}(0, q)=(p, 0), A_{p q}\left(q, q_{1}\right)=\left(p_{1}, p\right), g_{p q}(0, q)=0$, and $g_{p q}\left(q, q_{1}\right)=q_{1} p$. That $\Delta g_{p q} \leq 1$ is immediate from Lemma 4.2.

To check the moreover clause, we suppose $X \subseteq \mathbb{Z}^{2}$ is infinite and connected and let $X^{C}=\mathbb{Z}^{2} \backslash X$ denote its complement. If $\tilde{X}$ is the complement of any connected component of $X^{C}$, then $\Delta \mathbf{1}_{\tilde{X}} \leq \Delta \mathbf{1}_{X}$. In particular, we may assume that both $X$ and its complement $X^{C}$ are connected.

If $X^{C}$ intersects an $L_{p q}$ translation of the boundary of $[0, q]^{2}$, then, since $X$ is infinite and connected, the set $\left\{x \in X^{C}: \Delta \mathbf{1}_{X}(x)>0\right\}$ must intersect an $L_{p q}$ translation of the boundary of $[0, q]^{2}$. By Lemma 4.2, $\Delta g_{p q}(x)=1$ at any such point. Thus, we may assume that $X^{C}$ is contained in the interior of $[0, q]^{2}$. In particular, it suffices to prove the following claim.

Claim. Suppose $Y \subseteq[0, q]^{2}$ is simply connected, $Y \cap[1, q-1]^{2}$ is nonempty, $Y \backslash[1, q-1]^{2}$ is connected, and $\left(Y \backslash[1, q-1]^{2}\right) \cap E_{i}$ is nonempty for at most one $E_{i}$. Then there is a point $x \in Y \cap[1, q-1]^{2}$ such that $\Delta g_{p q}(x)-\Delta \mathbf{1}_{Y}(x)>1$.

We proceed by induction on $q$. This is trivial when $q=1$, since $[1, q-1]^{2}$ is empty. When $q=2$ we have $[1, q-1]=\{(1,1)\}$ and each of the four nearest neighbors of $(1,1)$ belongs to two of the $E_{i}$, so $Y=\{(1,1)\}$ and, by Lemma 4.2, we must have $\Delta g_{p q}(1,1)=-2$. We may therefore assume $q>2$, which in particular implies $q_{1} \neq q_{2}$. By symmetry, we may assume $q_{1}>q_{2}$, so the $E_{i}$ enjoy the intersection structure displayed on the left of Figure 4.2. We divide the analysis into several cases.

Case 1. Suppose $Y$ is contained in the interior of some $E_{i}$. We apply the induction hypothesis to find $x$.

Case 2. Suppose $x \in \partial Y \cap[1, q-1]^{2}$ lies in the boundary of exactly two $E_{i}$. Since Lemma 4.2 gives $\Delta g_{p q}(x)=1$, we have $\Delta\left(g_{p q}-\mathbf{1}_{Y}\right)(x)>1$.

Case 3. Suppose $x \in \partial Y \cap[1, q-1]^{2}$ lies in the boundary of exactly three $E_{i}$. Since we are not in the previous case, there are least two neighbors not in $Y$ and thus $\Delta \mathbf{1}_{Y}(x)<-2$. By Lemma 4.2, $\Delta g_{p q}(x)=0$.

Case 4. Suppose none of the above cases hold. Assuming $q_{1}>q_{2}$, then $Y$ does not intersect $E_{3}$ or $E_{4}$. Moreover, it must fail to intersect either $E_{1} \backslash[1, q-1]^{2}$ or $E_{2} \backslash[1, q-1]^{2}$. By symmetry, we may assume the former. If $Y \subseteq E_{1}$, then we can apply the induction hypothesis to $Y$. Otherwise, we can apply the induction hypothesis to $Y \cap E_{2}$.

The proof of Theorem 1.2 for the Ford circles is completed by the following lemma, which shows that the moreover clause of Proposition 4.3 suffices to show maximality.

Lemma 4.4. If $g: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ is superharmonic, $\Delta g$ is doubly periodic, and the only infinite connected set $X \subseteq \mathbb{Z}^{2}$ such that $\Delta\left(g+\mathbf{1}_{X}\right) \leq 1$ is $X=\mathbb{Z}^{2}$, then $g$ is maximal in the sense of Theorem 1.2.

Proof. Suppose for contradiction that $h: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ satisfies $h \geq g$ and $\Delta h \leq 1$, and that $h-g$ is unbounded. Replacing $h$ with $h-\min _{x \in \mathbb{Z}^{2}}(h-g)$, we may assume that $\{h-g=0\} \subseteq \mathbb{Z}^{2}$ is nonempty. Using the monotonicity of the Laplacian, we see that $f=\min \{h, g+1\}$ satisfies $\Delta f \leq 1$ and $f-g \in\{0,1\}$. Since $h-g$ is unbounded and $\Delta(h-g)$ is bounded above (here we use the periodicity of $\Delta g$ ), the set $\{f-g=1\} \subseteq \mathbb{Z}^{2}$ must contain connected components of arbitrarily large size. Note that, if $X$ is any connected component of $\{f-g=1\}$, then $\Delta\left(g+\mathbf{1}_{X}\right) \leq 1$. Using the periodicity of $\Delta g$ and compactness, we can select an infinite connected $X \subsetneq \mathbb{Z}^{2}$ such that $\Delta\left(g+\mathbf{1}_{X}\right) \leq 1$. Indeed, let $X_{n}$ be a sequence of connected components of $\{f-g=1\}$ such that $\left|X_{n}\right| \rightarrow \infty$. Since $\{f-g=0\} \subseteq \mathbb{Z}^{2}$ is nonempty, we may also select $x_{n} \in\{f-g=0\}$ that is adjacent to $X_{n}$. Using the fact that $\Delta g$ is periodic, we can translate so that $x_{n}$ always lies in the same period of $\Delta g$. Now, we can pass to a subsequence $n_{k}$ such that $x_{n_{k}}=x^{*}$ is a constant sequence. Now we can pass to a further subsequence such that for each $r \in \mathbb{N}$, the set $X_{n_{k}} \cap\left\{y \in \mathbb{Z}^{2}:\left|y-x^{*}\right| \leq r\right\}$ does not depend on $k$ for $k \geq r$. Then $\lim _{k \rightarrow \infty} X_{n_{k}}$ exists and has an infinite connected component $X^{*}$ adjacent to $x^{*}$. Note that $X^{*} \neq \mathbb{Z}^{2}$ (since $\left.x^{*} \notin X^{*}\right)$ and $\Delta g+\Delta 1_{X^{*}} \leq 1$, contradicting our hypothesis.
4.3. The diamond circles. The diamond circles are the circles $\left(c, c x_{1}, c x_{2}\right)$ $\in \mathcal{B}_{0}$ that are tangent to $(1,1,0)$ and $(1,1,2)$ and satisfy $0<x_{1}<1$. In curvature coordinates, the diamond circles can be parametrized via

$$
C_{k}:=\left(2 k(k+1), 2 k^{2}-1,2 k(k+1)\right)
$$

for $k \in \mathbb{Z}^{+}$. The peak matrix associated to the diamond circle $C_{k}$ is

$$
A_{k}:=\frac{1}{2}\left[\begin{array}{cc}
\frac{k}{k+1} & 1 \\
1 & \frac{1-k}{k}
\end{array}\right],
$$

which has lattice $L_{k}:=\left\{x \in \mathbb{Z}^{2}: A_{k} x \in \mathbb{Z}^{2}\right\}$ generated by $\{(0,-2 k),(k+1, k)\}$.

To better understand the structure of the diamond circle odometers $g_{k}$ : $\mathbb{Z}^{2} \rightarrow \mathbb{Z}$, we examine their Laplacians. Figure 4.3 displays the patterns associated to the first four diamond circles. We observe that the periodicity of each pattern is exactly $L_{k}$. Moreover, the internal structure of the fundamental tiles are similar enough that we can immediately conjecture what the general case should be. In fact, for the diamond circles, we have a closed formula for the odometer on each fundamental tile, making this family even simpler than the Ford circles.

Lemma 4.5. For each $k \in \mathbb{Z}_{+}$, there is a unique function $g_{k}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ such that

$$
\begin{array}{ll}
g_{k}(x+(0,-2 k))=g_{k}(x)+(-k, k-1) \cdot x-k(k-2) & \text { for } x \in \mathbb{Z}^{2}, \\
g_{k}(x+(k+1, k))=g_{k}(x)+(k, 1) \cdot x+\frac{1}{2} k(k+1) & \text { for } x \in \mathbb{Z}^{2}, \\
g_{k}(x+(-k-1, k))=g_{k}(x)+(0,-k) \cdot x-\frac{1}{2} k(k+1) & \text { for } x \in \mathbb{Z}^{2},
\end{array}
$$

and

$$
g_{k}(x)=\frac{1}{2}\left|x_{1}\right|\left(\left|x_{1}\right|-1\right)-\left\lfloor\frac{1}{4}\left(x_{1}-x_{2}\right)^{2}\right\rfloor \quad \text { for } x \in T_{k} \text {, }
$$

where

$$
T_{k}:=\left\{x \in \mathbb{Z}^{2}: \max \left\{\left|x_{1}\right|,\left|x_{2}-k\right|,\left|x_{1}\right|+\left|x_{2}-k\right|-1\right\} \leq k\right\} .
$$

Proof. This is similar to the corresponding proof for the Ford circles above. In particular, it suffices to check the consistency of the translation conditions on the set $T_{k} \backslash T_{k}^{\prime}$, where

$$
T_{k}^{\prime}:=\left\{x \in \mathbb{Z}^{2}:\left|x_{1}\right|+\left|x_{2}-k\right| \leq k-1\right\}
$$



Figure 4.3. Several periods of the Laplacian of the odometers $g_{k}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ of the first four diamond circles. Here, the black, dark patterned, light patterned, and white cells correspond to Laplacian values $1,0,-1$, and -2 , respectively.
is the interior of $T_{k}$. We claim that this follows from the fourth condition. If $x, x+(k+1, k) \in T_{k}$, then either $x=(-j, j)$ for $j=1, \ldots, k-1$ or $x=(-1-j, j)$ for $j=0, \ldots, k-1$. We then compute

$$
g_{k}(x+(k+1, k))=g_{k}(x)+(k, 1) \cdot x+\frac{1}{2} k(k+1)
$$

in either case. This implies the consistency of the second condition. The first and third conditions follow similarly.

Lemma 4.6. For each $k \in \mathbb{Z}_{+}$, the Laplacian $\Delta g_{k}$ satisfies

$$
\begin{array}{ll}
\Delta g_{k}(x+(0,2 k))=\Delta g_{k}(x+(k+1, k))=\Delta g_{k}(x) & \text { for } x \in \mathbb{Z}^{2}, \\
\Delta g_{k}(x)=1 & \text { for } x \in T_{k} \backslash T_{k}^{\prime}, \\
\Delta g_{k}(x)=(-1)^{x_{1}+x_{2}}-\mathbf{1}_{\{0\}}\left(x_{1}\right) & \text { for } x \in T_{k}^{\prime} .
\end{array}
$$

Proof. The first two conditions follow by computations analogous to those in the proof of Lemma 4.2. The formula in the third condition follows by inspection of the formula for $g_{k}$ on $T_{k}$.

Proposition 4.7. For every diamond circle $C_{k}$, the odometer $g_{k}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ satisfies

$$
\Delta g_{k}(x) \leq 1 \quad \text { and } \quad g_{k}(x+v)=g_{k}(x)+x \cdot A_{k} v+g_{k}(v)
$$

for all $x \in \mathbb{Z}^{2}$ and $v \in L_{k}$. Moreover, the only infinite connected subset $X \subseteq \mathbb{Z}^{2}$ such that $\Delta\left(g_{k}+\mathbf{1}_{X}\right) \leq 1$ is $X=\mathbb{Z}^{2}$.

Proof. This follows from the above lemma and the first part of the argument of Proposition 4.3. We can avoid the recursive argument in this case because of the explicit computation of $\Delta g_{k}$ on $T_{k}^{\prime}$ in Lemma 4.6.

## 5. Circles and lattices

5.1. Periodicity conditions. When defining odometers for the Ford and diamond circles above, we extended a finite amount of data to all of $\mathbb{Z}^{2}$ via periodicity conditions of the form

$$
g(x+v)=g(x)+a \cdot x+k \quad \text { for } x \in \mathbb{Z}^{2}
$$

where $v, a \in \mathbb{Z}^{2}$ and $k \in \mathbb{Z}$. In each of Lemmas 4.1 and 4.5, three pairs of vectors $\left(v_{i}, a_{i}\right)$ appear in such conditions, and these vectors have several suggestive properties. If we view each vector as a Gaussian integer $\mathbb{Z}[\mathbf{i}]$ via the usual identification $x \mapsto x_{1}+\mathbf{i} x_{2}$, then we have

$$
v_{1}+v_{2}+v_{3}=0 \quad \text { and } \quad a_{1}+a_{2}+a_{3}=0 .
$$

In this section, we generalize these vectors, associating pairs $\left(v_{i}, a_{i}\right)$ of vectors to each circle in $\mathcal{B}$. Our calculations largely follow those of Stange [23], who, motivated by data on our lattices $L_{C}$ for $C \in \mathcal{B}$ and Conway's association of


Figure 5.1. The Apollonian Band packing $\mathcal{B}_{0}$ with complex curvature coordinates, in the region $0<\operatorname{Im}(z)<2$. The circles tangent to the line $(0,-1)$ are Ford circles.
lattices to quadratic forms [5], studied ways to associate lattices to circles in an Apollonian circle packing.
5.2. Action of the Apollonian group. For the rest of this paper, we identify $\mathbb{C}$ with $\mathbb{R}^{2}$ and $\mathbb{Z}[\mathbf{i}]$ with $\mathbb{Z}^{2}$ in the usual way. Curvature coordinates are made complex by writing $C=(c, c z) \in \mathbb{R} \times \mathbb{C}$ for a circle with radius $c^{-1}$ and center $z \in \mathbb{C}$. Part of the band packing in complex curvature coordinates is shown in Figure 5.1. As noted in the introduction, a pairwise-tangent triple ( $C_{1}, C_{2}, C_{3}$ ) is related linearly to its Soddy circles $C_{0}, C_{4}$ in curvature coordinates by

$$
C_{0}+C_{4}=2\left(C_{1}+C_{2}+C_{3}\right) .
$$

This relation works also for lines, with the convention that the curvature coordinates of a line $\ell$ are $(0, z)$ where $z$ is the unit normal vector to the line, oriented away from the component of $\mathbb{R}^{2} \backslash \ell$ containing the other circles in the triple. In particular, all lines in the circle packing $\mathcal{B}$ have coordinates $(0,-1)$ or $(0,1)$.

A Descartes quadruple is a list of four circles such that any three form a pairwise-tangent triple. As any pairwise tangent triple of circles has exactly two Soddy circles, any pairwise tangent triple of circles can likewise be completed to exactly two Descartes quadruples, up to permutation. We call a Descartes quadruple $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ proper if the curvatures satisfy
$c_{0}>\max \left\{c_{1}, c_{2}, c_{3}\right\}$ and points of tangency between $C_{0}$ and $C_{1}, C_{2}, C_{3}$ are clockwise around $C_{0}$.

Note that, if $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ is a proper Descartes quadruple, then so is the parent rotation $\left(C_{0}, C_{2}, C_{3}, C_{1}\right)$ and the successor $\left(2\left(C_{0}+C_{2}+C_{3}\right)-\right.$ $C_{1}, C_{0}, C_{2}, C_{3}$ ). Each circle $C_{0} \in \mathcal{B}$ determines a Descartes quadruple ( $C_{0}, C_{1}$, $\left.C_{2}, C_{3}\right)$ up to parent rotation. Moreover, any pairwise tangent triple $\left(C_{1}, C_{2}, C_{3}\right)$ $\in \mathcal{B}^{3}$ can be completed to a proper Descartes quadruple ( $C_{0}, C_{1}, C_{2}, C_{3}$ ) in at most one way.
5.3. Lattice vectors. We now assign vectors $v\left(C, C^{\prime}\right), a\left(C, C^{\prime}\right) \in \mathbb{Z}[\mathbf{i}]$ to each pair of tangent circles $C, C^{\prime} \in \mathcal{B}$. In analogy to the two degenerate cases, the vectors $v\left(C, C^{\prime}\right)$ and $a\left(C, C^{\prime}\right)$ will generate periodicity conditions of the odometers associated to $C$ and $C^{\prime}$.

To describe our recursive construction, let us call a Descartes quadruple semi-proper if it is either proper or a parent rotation of

$$
\left(C_{0}, C_{1}, C_{2}, C_{3}\right)=((1,1+2 z),(1,1+2 z+2 \mathbf{i}),(0,1),(0,-1))
$$

for some $z \in \mathbb{Z}[\mathbf{i}]$. These latter quadruples are associated to the large circles in $\mathcal{B}$ and fail to be proper because one of the parents has the same curvature as the child. Note that, modulo parent rotations, the set of semi-proper quadruples is a directed forest in which each node has exactly one parent and three children. We induct along this tree structure to define our vectors.

For each semi-proper Descartes quadruple $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$, we define $v\left(C_{j}, C_{i}\right)$ and $a\left(C_{j}, C_{i}\right)$ for all rotations $(i, j, k)$ of $(1,2,3)$. For the base case, we consider the clockwise Descartes quadruple $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)=((1,1+2 z),(1,1+2 z+2 \mathbf{i}),(0,1),(0,-1)) \quad$ for $\quad z \in \mathbb{Z}[\mathbf{i}]$,
and we set

$$
\begin{array}{ll}
v\left(C_{3}, C_{2}\right)=0, & a\left(C_{3}, C_{2}\right)=1, \\
v\left(C_{2}, C_{1}\right)=1, & a\left(C_{2}, C_{1}\right)=z, \\
v\left(C_{1}, C_{3}\right)=-1, & a\left(C_{1}, C_{3}\right)=-1-z .
\end{array}
$$

For the induction step, we fix a semi-proper Descartes quadruple ( $C_{0}, C_{1}, C_{2}, C_{3}$ ) and suppose that $v\left(C_{j}, C_{i}\right)$ and $a\left(C_{j}, C_{i}\right)$ have been defined for all rotations $(i, j, k)$ of $(1,2,3)$. For the successor quadruple $\left(2\left(C_{0}+C_{2}+C_{3}\right)-C_{1}, C_{0}, C_{2}, C_{3}\right)$, we define

$$
\begin{aligned}
& v\left(C_{2}, C_{0}\right)=v\left(C_{2}, C_{1}\right)-\mathbf{i} v\left(C_{3}, C_{2}\right), \\
& v\left(C_{0}, C_{3}\right)=v\left(C_{1}, C_{3}\right)+\mathbf{i} v\left(C_{3}, C_{2}\right), \\
& a\left(C_{2}, C_{0}\right)=a\left(C_{2}, C_{1}\right)+\mathbf{i} a\left(C_{3}, C_{2}\right), \\
& a\left(C_{0}, C_{3}\right)=a\left(C_{1}, C_{3}\right)-\mathbf{i} a\left(C_{3}, C_{2}\right),
\end{aligned}
$$

so that $v\left(C_{j}, C_{i}\right)$ and $a\left(C_{j}, C_{i}\right)$ are defined for all rotations $(i, j, k)$ of $(0,2,3)$. Since every pair of tangent circles $C, C^{\prime} \in \mathcal{B}$ can be completed to a proper

Descartes quadruple of the form $\left(C_{0}, C, C^{\prime}, C_{3}\right)$, this recursive construction generates vectors for all pairs of tangent circles in $\mathcal{B}$. Since the quadruple of the big circle $(1,1+2 z)$ is a successor of the quadruple of $(1,1+2 z+2 \mathbf{i})$, we must also check everything is well defined, but this is immediate from the definition.

Lemma 5.1. If $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ is a proper Descartes quadruple and we write $C_{i}=\left(c_{i}, c_{i} z_{i}\right), v_{i j}=v\left(C_{i}, C_{j}\right)$, and $a_{i j}=a\left(C_{i}, C_{j}\right)$, then the following hold:

$$
\begin{array}{ll}
v_{32}+v_{13}+v_{21}=0, & a_{32}+a_{13}+a_{21}=0 \\
v_{10}=v_{13}-\mathbf{i} v_{21}, & a_{10}=a_{13}+\mathbf{i} a_{21} \\
v_{01}=\mathbf{i} v_{10}, & a_{01}=-\mathbf{i} a_{10} \\
v_{32}^{2}=c_{3} c_{2}\left(z_{3}-z_{2}\right), & 2 v_{32} a_{32}=c_{3} z_{3}+c_{2} z_{2} \\
\bar{v}_{13} v_{21}+v_{13} \bar{v}_{21}=-2 c_{1} . & \tag{5.1e}
\end{array}
$$

Properties (5.1a), (5.1b), and (5.1c) give inductive relationships among the vectors. The property (5.1d) expresses the vectors $v\left(C, C^{\prime}\right)$ and $a\left(C, C^{\prime}\right)$ up to sign in terms of the circles $C, C^{\prime} \in \mathcal{B}$. Finally, (5.1e) implies that the determinant of the lattice generated by $v_{31}$ and $v_{21}$ is $c_{1}$. Observe that since each statement holds also for any rotation of the Descartes quadruple, we have, e.g., that $v_{02}=\mathbf{i} v_{20}$, etc.

Proof of Lemma 5.1. The inductive construction gives (5.1a)-(5.1c) immediately. To check (5.1d), we first observe that it holds for the base case and its immediate successors. By induction, it suffices to assume (5.1d) holds for the proper Descartes quadruples $\left(C_{0}, C_{1}, C_{2}, C_{3}\right),\left(C_{1}, C_{4}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ and check the conditions for each of the three successors of $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$. We write $w_{i}=c_{i} z_{i}$ and compute

$$
\begin{aligned}
v_{10}^{2} & =\left(v_{13}-\mathbf{i} v_{21}\right)^{2} \\
& =\left(v_{13}-v_{12}\right)^{2} \\
& =2 v_{13}^{2}+2 v_{12}^{2}-\left(v_{13}+v_{12}\right)^{2} \\
& =2 v_{13}^{2}+2 v_{12}^{2}-v_{14}^{2} \\
& =2\left(c_{3} w_{1}-c_{1} w_{3}\right)+2\left(c_{2} w_{1}-c_{1} w_{2}\right)-\left(c_{4} w_{1}-c_{1} w_{4}\right) \\
& =\left(2\left(c_{1}+c_{2}+c_{3}\right)-c_{4}\right) w_{1}-c_{1}\left(2\left(w_{1}+w_{2}+w_{3}\right)-w_{4}\right) \\
& =c_{0} w_{1}-c_{1} w_{0}
\end{aligned}
$$

using the induction hypotheses and the linear Soddy relation. We also compute

$$
\begin{aligned}
v_{10} a_{10} & =\left(v_{13}-\mathbf{i} v_{21}\right)\left(a_{13}+\mathbf{i} a_{21}\right) \\
& =v_{13} a_{13}+v_{21} a_{21}+\mathbf{i}\left(v_{13} a_{21}-v_{21} a_{13}\right) \\
& =v_{13} a_{13}+v_{21} a_{21}-v_{13} a_{12}-v_{12} a_{13}
\end{aligned}
$$

$$
\begin{aligned}
& =2 v_{13} a_{13}+2 v_{21} a_{21}-\left(v_{13}+v_{12}\right)\left(a_{13}+a_{12}\right) \\
& =2 v_{13} a_{13}+2 v_{21} a_{21}-v_{14} a_{14} \\
& =\frac{1}{2}\left(w_{1}+2\left(w_{1}+w_{2}+w_{3}\right)-w_{4}\right) \\
& =\frac{1}{2}\left(w_{1}+w_{0}\right)
\end{aligned}
$$

The relations $v_{i 0}^{2}=c_{i} c_{0}\left(z_{i}-z_{0}\right)$ and $2 v_{i 0} a_{i 0}=c_{i} z_{i}+c_{0} z_{0}$ for $i=2,3$ follow by analogous computations.

The relation (5.1e) can be obtained as follows. We make the observation that $v_{i j}^{2}=c_{i} c_{j}\left(z_{i}-z_{j}\right)$ implies $\left|v_{i j}\right|^{2}=c_{i}+c_{j}$ and therefore

$$
\begin{aligned}
v_{13} v_{21}\left(\bar{v}_{13} v_{21}+v_{13} \bar{v}_{21}\right) & =\left(c_{1}+c_{3}\right) v_{21}^{2}+\left(c_{2}+c_{1}\right) v_{13}^{2} \\
& =c_{1}\left(v_{21}^{2}+v_{13}^{2}\right)+c_{3}\left(c_{1} w_{2}-c_{2} w_{1}\right)+c_{2}\left(c_{3} w_{1}-c_{1} w_{3}\right) \\
& =c_{1}\left(v_{21}^{2}+v_{13}^{2}\right)-c_{1}\left(c_{3} w_{2}-c_{2} w_{3}\right) \\
& =c_{1}\left(v_{21}^{2}+v_{13}^{2}-v_{32}^{2}\right) \\
& =-2 c_{1} v_{21} v_{13}
\end{aligned}
$$

Since at most one of the $C_{i}$ is a line when we are not in the base case, we have $v_{i j} \neq 0$ and thus obtain the next relation.

Given a circle $C_{0}$ in the band packing $\mathcal{B}$ with the proper Descartes quadruple ( $C_{0}, C_{1}, C_{2}, C_{3}$ ), we define the lattice

$$
\Lambda_{C_{0}}=\mathbb{Z} v\left(C_{1}, C_{0}\right)+\mathbb{Z} v\left(C_{2}, C_{0}\right)+\mathbb{Z} v\left(C_{3}, C_{0}\right) .
$$

(Any two of the three summands suffice since $\sum_{i=1}^{3} v\left(C_{i}, C_{0}\right)=0$.) Next we compare $\Lambda_{C}$ with the lattice $L_{C}$ of (1.4).

Lemma 5.2. $\Lambda_{C} \subseteq L_{C}$ for all $C \in \mathcal{B}$.
Proof. If $C_{0} \in \mathcal{B}$ has radius 1 , then $\Lambda_{C_{0}}=L_{C_{0}}=\mathbb{Z}^{2}$ and there is nothing to prove. Otherwise, we may assume that $C_{0}$ is part of a proper Descartes quadruple ( $C_{0}, C_{1}, C_{2}, C_{3}$ ) $\in \mathcal{B}^{4}$. Using Lemma 5.1, we compute for $i=1,2,3$

$$
\begin{align*}
a_{i 0} & =\frac{1}{2}\left(c_{i} z_{i}+c_{0} z_{0}\right) v_{i 0}^{-1} \\
& =\frac{1}{2}\left(c_{0}^{-1} v_{i 0}^{2}+\left(c_{i}+c_{0}\right) z_{0}\right) v_{i 0}^{-1}  \tag{5.2}\\
& =\frac{1}{2} c_{0}^{-1} v_{i 0}+\frac{1}{2} z_{0} \overline{v_{i 0}} \\
& =A_{C_{0}} v_{i 0} .
\end{align*}
$$

Since $A_{C_{0}} v_{i 0}=a_{i 0} \in \mathbb{Z}[\mathbf{i}]$, we conclude that $v_{i 0} \in L_{C_{0}}$ for $i=1,2,3$.
In fact $\Lambda_{C}=L_{C}$ for all $C \in \mathcal{B}$, but our proof of the inclusion $L_{C} \subseteq \Lambda_{C}$ uses our construction of the odometers $g_{C}$ and is thus postponed to the end of Section 10.

Observe that (5.1b) can be rewritten as $v_{0 i}=v_{j i} \pm v_{k i}$, where $(i, j, k)$ is a rotation of $(1,2,3)$ and sign of $v_{k i}$ depends on whether or not $C_{i}$ is a parent of $C_{k}$. Inductively, this implies that a vector $v\left(C, C_{0}\right)$ lives in the lattice of the circle $C_{0}$ :

Lemma 5.3. Given a proper Descartes quadruple $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ and any $C \in \mathcal{B}$ that is tangent to and smaller than $C_{0}$, we have that $v\left(C, C_{0}\right)$ $\in \Lambda_{C_{0}}$.

Finally, we check that our general lattice vector construction agrees with the degenerate cases.

Proposition 5.4. Every Ford circle $C_{0}=\left(q^{2}, 1+2 p q \mathbf{i}\right)$ is part of a proper quadruple $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$, where $C_{1}=\left(q_{1}^{2}, 1+2 p_{1} q_{1} \mathbf{i}\right), C_{2}=\left(q_{2}^{2}, 1+2 p_{2} q_{2} \mathbf{i}\right)$, and $C_{3}=(0,-1)$. Moreover, we have

$$
\begin{array}{ll}
v\left(C_{1}, C_{0}\right)=q+q_{1} \mathbf{i}, & a\left(C_{1}, C_{0}\right)=p_{1}+p \mathbf{i}, \\
v\left(C_{2}, C_{0}\right)=-q+q_{2} \mathbf{i}, & a\left(C_{2}, C_{0}\right)=p_{2}-p \mathbf{i}, \\
v\left(C_{3}, C_{0}\right)=-q \mathbf{i}, & a\left(C_{3}, C_{0}\right)=-p \mathbf{i},
\end{array}
$$

which are exactly the vectors appearing in Lemma 4.1.
Proof. That $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$ is proper follows from the discussion in Section 4.2. Since (5.1a) and (5.1d) determine the tuple ( $v_{10}, v_{20}, v_{30}, a_{10}, a_{20}, a_{30}$ ) up to sign, it is enough to show $v_{30}=-q$ i. Using (5.1b) and (5.1c), we inductively compute

$$
v_{30}=v_{32}-\mathbf{i} v_{13}=v_{32}+\mathbf{i} v_{13}=-\mathbf{i} q_{2}-\mathbf{i} q_{1}=-\mathbf{i} q .
$$

We conclude after reading off the base case $C_{0}=(1,1+2)$ from the beginning of this section.

Proposition 5.5. Every diamond circle

$$
C_{0}=\left(2 k(k+1), 2 k^{2}-1+2 k(k+1) \mathbf{i}\right)
$$

is part of a proper quadruple $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$, where $C_{1}=(1,1+2 \mathbf{i}), C_{2}=$ $(1,1)$, and $C_{3}=\left(2(k-1) k, 2(k-1)^{2}-1+2(k-1) k \mathbf{i}\right)$. Moreover, we have

$$
\begin{array}{ll}
v\left(C_{1}, C_{0}\right)=k+1+k \mathbf{i}, & a\left(C_{1}, C_{0}\right)=k+\mathbf{i}, \\
v\left(C_{2}, C_{0}\right)=-k-1+k \mathbf{i}, & a\left(C_{2}, C_{0}\right)=-k \mathbf{i}, \\
v\left(C_{3}, C_{0}\right)=-2 k \mathbf{i}, & a\left(C_{3}, C_{0}\right)=-k+(k-1) \mathbf{i},
\end{array}
$$

which are exactly the vectors appearing in Lemma 4.5.
Proof. That $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$ is proper follows from the discussion in Section 4.3. As in the previous proposition, it is enough to show $v\left(C_{1}, C_{0}\right)=$
$k+1+k \mathbf{i}$. Using (5.1b) and (5.1c), we inductively compute

$$
v_{10}=v_{13}-\mathbf{i} v_{21}=(k+(k-1) \mathbf{i})+(1+\mathbf{i})=k+1+k \mathbf{i} .
$$

We conclude after reading off the base case $C_{0}=(4,1+4 \mathbf{i})$ from the previous proposition.
5.4. Symmetry reduction. The band packing $\mathcal{B}$ is invariant under the operations

$$
\begin{aligned}
& (c, z) \mapsto(c,-z), \\
& (c, z) \mapsto(c, \bar{z}), \\
& (c, z) \mapsto(c, z+2 c \mathbf{i}), \\
& (c, z) \mapsto(c, z+2 c) .
\end{aligned}
$$

These operations can be extended to the vectors $v\left(C, C^{\prime}\right)$ and $a\left(C, C^{\prime}\right)$ in the obvious way. For example, if we apply the shift $(c, z) \mapsto(c, z+2 c w)$ for some $w \in \mathbb{Z}[\mathbf{i}]$, then $v\left(C, C^{\prime}\right)$ is unchanged while $a\left(C, C^{\prime}\right)$ is replaced by $a\left(C, C^{\prime}\right)+$ $c c^{\prime} v\left(C, C^{\prime}\right)^{-1}$.

We can extend the results of Section 4 to the orbit of the Ford and diamond circles under these symmetries. For example, suppose that $C=(c, z)$ is a Ford circle and we want to construct an odometer for the shifted circle $C^{\prime}=$ $(c, z+2 c w)$ for $w=a+b \mathbf{i} \in \mathbb{Z}[\mathbf{i}]$. Observe that

$$
A_{C^{\prime}}-A_{C}=\left[\begin{array}{cc}
a & b \\
b & -a
\end{array}\right]
$$

and that the function $h_{w}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ given by

$$
h_{w}(x)=\frac{a}{2} x_{1}\left(x_{1}+1\right)-\frac{a}{2} x_{2}\left(x_{2}+1\right)+b x_{1} x_{2}
$$

satisfies $\Delta h_{w} \equiv 0$. In particular, $g_{C^{\prime}}=g_{C}+h_{w}$ is an odometer for $C^{\prime}$.

## 6. A topological lemma

In the course of our general tile construction, it is necessary to translate local knowledge of tile compatibility to global knowledge regarding the intersection structure of the collection of tiles. This section concerns a technical lemma that is our primary tool for checking that certain collections of tiles form tilings using local conditions.

We continue to identify $\mathbb{Z}[\mathbf{i}]$ with $\mathbb{Z}^{2}$ in the natural way. Thus a tile $T$ is a set of squares $s_{x}=\{x, x+1, x+\mathbf{i}, x+1+\mathbf{i}\} \subseteq \mathbb{Z}[\mathbf{i}]$, where $I(T)$ (adapted from (1.6) to $\mathbb{C}$ in the obvious way) is a topological disc. We let $\mathbb{Z}[\mathbf{i}]$ inherit the standard degree-4 square lattice graph of $\mathbb{Z}^{2}$. For any tile, define the footprint $\mathrm{f}(T)$ to be the subgraph of $\mathbb{Z}[\mathbf{i}]$ induced by the union of squares $s_{x} \in T$. For two tiles $T_{1}, T_{2}$, write $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$ for the induced subgraph on the intersection of the vertex sets of $\mathrm{f}\left(T_{1}\right)$ and $\mathrm{f}\left(T_{2}\right)$. In an abuse of terminology, we say that
tiles $T_{1}$ and $T_{2}$ intersect if $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right) \neq \varnothing$ and overlap if $T_{1} \cap T_{2} \neq \varnothing$. If $T_{1} \cap T_{2}=\varnothing$, they are nonoverlapping.

Recall from Section 1.2 that we define a tiling as a collection of tiles such that every square $s_{x}$ of $\mathbb{Z}[\mathbf{i}]$ belongs to a unique tile. Thus, the tiles in a tiling are permitted to intersect, but not to overlap.

Lemma 6.1. Suppose $\mathcal{T}$ is an (infinite) collection of tiles and $\mathcal{E}$ is a set of two-element subsets of $\mathcal{T}$ satisfying the following four hypotheses:
(1) The graph $G=(\mathcal{T}, \mathcal{E})$ is a 3 -connected planar triangulation.
(2) $\mathcal{T}$ and $\mathcal{E}$ are invariant under translation by some full-rank lattice $L \subseteq \mathbb{Z}[\mathbf{i}]$, and $\sum_{T \in \mathcal{T} / L}|T|=|\operatorname{det} L|$.
(3) If $\left\{T_{1}, T_{2}\right\} \in \mathcal{E}$, then the intersection $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$ contains at least two vertices.
(4) We can select for each face $F=\left\{T_{1}, T_{2}, T_{3}\right\}$ of $G$ a point $\rho(F) \in \mathrm{f}\left(T_{1}\right) \cap$ $\mathrm{f}\left(T_{2}\right) \cap \mathrm{f}\left(T_{3}\right)$ such that, for the face $F^{\prime}=\left\{T_{1}, T_{2}, T_{4}\right\}$ of $G, \rho(F)$ and $\rho\left(F^{\prime}\right)$ lie on some path contained in $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$.
Then $\mathcal{T}$ is a tiling. Moreover, if $T_{1}, T_{2} \in \mathcal{T}$ and $T_{1}$ and $T_{2}$ intersect, then $\left\{T_{1}, T_{2}\right\} \in \mathcal{E}$, and $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$ is a path joining $\rho(F)$ and $\rho\left(F^{\prime}\right)$ for the faces $F, F^{\prime}$ with $F \cap F^{\prime}=\left\{T_{1}, T_{2}\right\}$.

The proof is essentially a homotopy argument. The task is complicated slightly by the fact that we assume in (4) merely that each intersection $f\left(T_{1}\right) \cap$ $\mathrm{f}\left(T_{2}\right)$ contains a path, and not that it is equal to a path. Although this weaker assumption prevents us from deducing the Lemma from standard topological facts, it will significantly simplify our inductive argument in Section 7; see Remark 7.10.

Before commencing, recall that the winding number of a closed walk in $\mathbb{Z}[\mathbf{i}]$ about a point $z^{*} \in \mathbb{Z}[\mathbf{i}]^{*}$ is the number of times the walk circles the point $z^{*}$ counterclockwise. (Here, $\mathbb{Z}[\mathbf{i}]^{*}$ is the dual lattice to $\mathbb{Z}[\mathbf{i}]$, whose vertices are centered in the faces of $\mathbb{Z}[\mathbf{i}]$.) Thus in terms of the vertices $z_{1}, z_{2}, \ldots, z_{t}=z_{1}$ of the walk, the winding number is given by

$$
\operatorname{wind}\left(z_{1}, \ldots, z_{t} ; z^{*}\right):=\frac{1}{2 \pi} \sum_{i=1}^{t-1} \arg \left(\frac{z_{i+1}-z^{*}}{z_{i}-z^{*}}\right) .
$$

We say that the closed walk $W$ encloses a point $z^{*}$ if the winding number of $W$ about $z^{*}$ is nonzero; similarly, we say that $W$ encloses a square $s_{x}=$ $\{x, x+1, x+\mathbf{i}, x+1+\mathbf{i}\} \subseteq \mathbb{Z}[\mathbf{i}]$ if the winding number of the walk about $s_{x}^{*}=x+\frac{1}{2}+\frac{\mathbf{i}}{2}$ is nonzero. Note that the definition of a tile (via (1.6)) implies that any tile $T$ is a set of squares enclosed by a simple cycle $\partial T$ in $\mathbb{Z}^{2}$, which is the set of points and edges of $T$ that each also lie in a square $s_{x} \notin T$.

Proof of Lemma 6.1. Our first goal is to show that $\mathcal{T}$ is a tiling. Note that it suffices to show that every square $s_{x}=\{x, x+1, x+\mathbf{i}, x+1+\mathbf{i}\}$ lies in
some tile in $\mathcal{T}$, since hypothesis (2) implies that the average number of tiles a square lies in is 1 . The idea is to use the periodicity to draw a large cycle that surrounds a given square $s_{x}$ and then use the graph structure to contract this cycle to the boundary of a single tile.

We fix an $L$-periodic drawing of $G=(\mathcal{T}, \mathcal{E})$ and work with the dual graph $G^{*}$, whose vertex set $\mathcal{F}$ is the set of triangles in the $\operatorname{graph}(\mathcal{T}, \mathcal{E})$. For $F \in \mathcal{F}$, we have that $\rho(F)$ is a point in the intersection of the three tiles in $F$, and for adjacent $F, F^{\prime}$, we let $\rho\left(F, F^{\prime}\right)$ be a choice of path from $\rho(F)$ to $\rho\left(F^{\prime}\right)$ in the intersection of the two tiles in $F \cap F^{\prime}$, guaranteed to exist by hypothesis (4). Given a sequence $F^{0}, F^{1}, \ldots, F^{n}=F^{0}$ with $\left|F^{i} \cap F^{i+1}\right|=2$ (indices evaluated modulo $n$ ), we define the closed walk $\ell\left(F^{0}, \ldots, F^{n}\right)$ in $\mathbb{Z}^{2}$ as the consecutive concatenation of the paths $\rho\left(F^{i}, F^{i+1}\right)(0 \leq i<n)$.

Since the graph $(\mathcal{T}, \mathcal{E})$ is connected and periodic under a nontrivial lattice, we can, given any finite set $Z \subseteq \mathbb{Z}[\mathbf{i}]^{*}$, find a cycle $F^{0}, \ldots, F^{n}=F^{0}$ in $G^{*}$ that wraps around each point $z^{*} \in Z$, just in the sense that

$$
\operatorname{wind}\left(\rho\left(F^{1}\right), \ldots, \rho\left(F^{n}\right) ; z^{*}\right)=1
$$

for all $z^{*} \in Z$. Since hypothesis (2) also implies that the tiles $T \in \mathcal{T}$ are of bounded size, this means that given the point $s_{x}^{*} \in \mathbb{Z}[\mathbf{i}]^{*}$, we can (by choosing $Z$ to be the set of all points within some sufficiently large distance of $s_{x}^{*}$ ) find a cycle $C$ given by $F^{0}, F^{1}, \ldots, F^{n}=F^{0}$ such that $\ell(C)$ encloses $s_{x}^{*}$ and so $s_{x}$. Letting $\mathcal{K}$ denote the set of faces of $G^{*}$ in the region of the plane bounded by $C$ in our fixed embedding (so, each element of $\mathcal{K}$ corresponds to some $T \in \mathcal{T}$ ), we choose this $C$ such that $|\mathcal{K}|$ is as small as possible, subject to the condition that $\ell(C)$ encloses $s_{x}^{*}$.

Note that if $|\mathcal{K}|=1$, then $s_{x}$ is indeed covered by $\mathcal{T}$, since then all vertices in $\ell(C)$ belong to a single tile, which is simply connected and would thus cover $s_{x}$. Moreover, any element of $\mathcal{K}$ whose intersection with $C$ is discontiguous is a cut-vertex $\{f \in \mathcal{K} \mid C \cap f \neq \varnothing\}$, when this collection is viewed as a subgraph of the original graph $G$. (This occus for the third leftmost member of $\mathcal{K}$ in Figure 6.1, for example.) Since not every vertex in a finite graph can be a cut-vertex, we may therefore assume without loss of generality that there is face $C^{\prime \prime}$ of $G^{*}$ that is a member of $\mathcal{K}$ and whose intersection with $C$ is contiguous; in other words, whose boundary cycle is

$$
F^{n}, F^{n-1}, \ldots, F^{\ell}, E_{1}, \ldots, E_{t}, F^{0} \quad(t>0),
$$

where no $E_{i}$ lies on $C$. We now consider the closed walk $\ell\left(C^{\prime}\right)$ in $\mathbb{Z}[\mathbf{i}]$, where $C^{\prime}$ is the cycle

$$
F^{\ell}, F^{\ell-1}, F^{0}, E_{t}, \ldots, E_{1}, F^{\ell}
$$



Figure 6.1. Pulling across $\mathcal{K}$. The triangulation $G$ is drawn. The six circled vertices correspond to faces of $G^{*}$ that belong to $\mathcal{K}$. The curve drawn is the cycle $C$; the faces $F^{0}, \ldots, F^{n-1}$ of $G$ are marked by o's, and the $E_{i}$ 's (there is only one here) are marked by $\boldsymbol{\square}$. The twice circled vertex of $G$ corresponds to the face $C^{\prime \prime}$ of $G^{*}$.
in $G^{*}$. Note that the region bounded by $C^{\prime}$ has exactly one fewer face than that bounded by $C$. We will show that if $\ell(C)$ enclosed $s_{x}$, then so must $\ell\left(C^{\prime}\right)$, contradicting minimality of $C$.

To begin, note that the face $C^{\prime \prime}$ of the dual graph $G^{*}$ corresponds to some vertex $T \in \mathcal{T}$, which completely contains the walk $\ell\left(C^{\prime \prime}\right)$. Thus, $\ell\left(C^{\prime \prime}\right)$ does not enclose $s_{x}$ unless $s_{x}$ lies in $T$. Finally, the winding number about any square $s_{x}$ of the loop $\ell\left(C^{\prime}\right)$ is the same as the winding number about the square $s_{x}$ of the concatenation of $\ell(C)$ with $\ell\left(C^{\prime \prime}\right)$. In particular, since $s_{x}$ is not enclosed by $\ell\left(C^{\prime}\right)$, the winding number of $\ell(C)$ and $\ell\left(C^{\prime \prime}\right)$ about $s_{x}$ must be equal. Thus $\ell\left(C^{\prime \prime}\right)$ also encloses $s_{x}$, contradicting the minimality of $C$ with this property.

Having shown that $\mathcal{T}$ is a tiling, we next wish to show that if $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$ $\neq \varnothing$, then $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$ is a path in $\mathbb{Z}[\mathbf{i}]$ from $\rho(F)$ to $\rho\left(F^{\prime}\right)$ for the faces $F, F^{\prime}$ whose intersection is the pair $\left\{T_{1}, T_{2}\right\}$. We begin by showing that the intersection is a path.

If it is not, there are paths $P_{1} \subseteq \partial T_{1}$ and $P_{2} \subseteq \partial T_{2}$ with the same pair of endpoints, whose concatenation $C$ is a cycle enclosing a region $S^{*}$ of $\mathbb{Z}[\mathbf{i}]^{*}$ disjoint from $T_{1}^{*}$ and $T_{2}^{*}$. (The dual $T^{*}$ of $T$ consists of the points of the dual of $\mathbb{Z}^{2}$ that are centers of squares $s_{x} \subseteq T$.) Among all possible pairs $T_{1}, T_{2}$, we may assume we have chosen such that $\left|S^{*}\right|$ is as large as possible. (Note
that there is some absolute bound on $\left|S^{*}\right|$, since, for example, hypothesis (2) implies that tiles have bounded size.)

We let $\mathcal{T}_{S}$ denote the tiles in the region bounded by $C$. Since $(\mathcal{T}, \mathcal{E})$ is 3 -connected, there must be at least three tiles in $\mathcal{T} \backslash \mathcal{T}_{S}$ that are adjacent to tiles in $\mathcal{T}_{S}$. By hypothesis (3), such a tile $T$ must have the property that it shares two vertices with some tile in $\mathcal{T}_{S}$; however, the only candidate points to be shared between a tile $T \notin \mathcal{T}_{S}, T \neq T_{1}, T_{2}$, and a tile in $\mathcal{T}_{S}$ are the two common endpoints of $P_{1}$ and $P_{2}$, and for $C$ to be a (simple) cycle there is, for each of these two endpoints, at most one square of $\mathbb{Z}[\mathbf{i}]$ containing the point and lying outside $C$ and outside of the tiles $T_{1}, T_{2}$. In particular, there must be exactly three tiles in $\mathcal{T} \backslash \mathcal{T}_{S}$ adjacent to tiles in $\mathcal{T}_{S}$; namely, $T_{1}, T_{2}$, and a third tile $T_{3}$ that includes both endpoints of the paths $P_{1}, P_{2}$. But now either the pair $\left\{T_{3}, T_{2}\right\}$ or the pair $\left\{T_{3}, T_{1}\right\}$ contradict the maximality of the choice of $S^{*}$. Thus $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$ is indeed a path.
$\rho(F)$ and $\rho\left(F^{\prime}\right)$ both lie in $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$. If they are not the endpoints of the path, then, without loss of generality, let $e_{1}, e_{2}$ be the two edges of the path $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$ that are incident with $\rho(F)$. Since $T_{1}$ and $T_{2}$ are nonoverlapping, for each $i=1,2, e_{i}$ lies in one square from $T_{1}$ and one square of $T_{2}$. Moreover, their shared endpoint $\rho(F)$ lies also in a third tile $T_{3}$ (which again, is nonoverlapping with $T_{1}, T_{2}$ ). In particular, either $T_{1}$ or $T_{2}$ must contain two diagonally opposite squares about $\rho(F)$ without containing the other two squares, contradicting the definition of a tile.

Finally, we wish to show that if $f\left(T_{1}\right) \cap f\left(T_{2}\right) \neq \emptyset$, then $\left\{T_{1}, T_{2}\right\} \in \mathcal{E}$. We do this by giving a suitable plane drawing of the graph $G=(\mathcal{T}, \mathcal{E})$. For each tile $T \in \mathcal{T}$, we draw a vertex $v_{T}$ corresponding to $T$ at some point of the dual $T^{*}$. Identifying $\mathbb{C}$ with the Euclidean plane, let $\bar{s}_{x}=\{x+s+t \mathbf{i} \mid 0 \leq s, t \leq 1\} \subseteq \mathbb{C}$. Since we know that tiles $T, T^{\prime} \in \mathcal{T}$ that are adjacent in $G$ must share an edge of $\mathbb{Z}[\mathbf{i}]$, we can draw a curve from $v_{T}$ to $v_{T^{\prime}}$ such that every point in the curve lies in the interior of $\bar{s}_{x} \cup \bar{s}_{y}$ where $s_{x}, s_{y}$ each lie in $T$ or $T^{\prime}$; in particular, we can draw all edges of $G$ such that they are pairwise nonintersecting (except at shared endpoints) and such that the edge from $v_{T}$ to $v_{T^{\prime}}$ is disjoint from any $\bar{s}_{x}$ for an $s_{x}$ not contained in $T$ or $T^{\prime}$. With this drawing the curve $C_{T}$ in $\mathbb{C}$ corresponding to the cycle through the neighbors of a tile $T$ is disjoint from $T$ and bounds a region containing $T$. Since $C_{T}$ and $C_{T^{\prime}}$ bound disjoint regions of $\mathbb{C}$ when $T, T^{\prime}$ are nonadjacent in $G$, we see that any nonadjacent tiles are nonintersecting.

## 7. Tiles

In this section, we associate to each circle $C \in \mathcal{B}$ a tile (unique up to translation) that will be $90^{\circ}$ symmetric and tile the plane under translation by the lattice $\Lambda_{C}$. Before beginning our construction, we need a few additional
definitions regarding tiles. We let $\mathrm{c}(T)$ and $|T|$ denote the centroid and area of a tile $T$, which are, respectively, the centroid and area of the real subset $I(T)$ from (1.6). We say that $T_{1}$ and $T_{2}$ touch if they are nonoverlapping and $\partial T_{1} \cap \partial T_{2}$ is a simple path of $\mathbb{Z}^{2}$ with at least two vertices. We say that three tiles form a touching triple of tiles if they are pairwise touching and share exactly one common boundary vertex.

Recall that a tile is a set of squares $s_{x}$ and that the footprint of a tile $T$ is

$$
\mathrm{f}(T):=\bigcup_{s_{x} \in T} s_{x}
$$

It will be convenient to allow a degenerate case of our tile definition. Note that if $T=\varnothing$, then the centroid $\mathrm{c}(T)$ would be undefined. We will allow tiles $T=\varnothing$, whose footprints $\mathrm{f}(T)$ may be any singleton from $\mathbb{Z}[\mathbf{i}]$; this singleton then gives the centroid of $T$. In particular, note that the degenerate tiles are in bijective correspondence with $\mathbb{Z}[\mathbf{i}]$. If $T$ is degenerate, we say that $T, T^{\prime}$ touch if $\mathrm{f}(T) \subseteq \partial T^{\prime}$. We emphasize that $T^{\prime} \backslash T=T^{\prime}$ whenever $T$ is degenerate.

Let a prototile $T$ be a set of squares $s_{x} \subseteq \mathbb{Z}[\mathbf{i}]$ and, if empty, have $\mathrm{f}(T)$ assigned as a singleton in $\mathbb{Z}[\mathbf{i}]$. (Compared with the definition of a tile, we are dropping the requirement that $I(T)$ is a topological disk.) We begin by recursively associating a prototile to each $C \in \mathcal{B}$; in Lemma 7.3 we will verify that these prototiles are in fact tiles.

Definition 7.1. If $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ is a proper Descartes quadruple, then a set $T_{0}$ of squares $s_{x} \subseteq \mathbb{Z}[\mathbf{i}]$ is a prototile for $C_{0}$ if $T_{0}$ has the tile decomposition

$$
\begin{equation*}
T_{0}=T_{1}^{+} \cup T_{1}^{-} \cup T_{2}^{+} \cup T_{2}^{-} \cup T_{3}^{+} \cup T_{3}^{-}, \tag{7.1}
\end{equation*}
$$

with $\mathrm{f}\left(T_{i}^{ \pm}\right) \subseteq \mathrm{f}\left(T_{0}\right)$ even if $T_{i}$ is degenerate where, for each rotation $(i, j, k)$ of $(1,2,3), T_{i}^{ \pm}$is a prototile of $C_{i}$ satisfying

$$
\begin{equation*}
\mathrm{c}\left(T_{i}^{ \pm}\right)-\mathrm{c}\left(T_{0}\right)= \pm \frac{1}{2}\left(v_{k j}-\mathbf{i} v_{k j}\right), \tag{7.2}
\end{equation*}
$$

where $v_{i j}:=v\left(C_{i}, C_{j}\right)$.
The base cases are those circles in $\mathcal{B}$ that are not the first circle of any proper Descartes quadruple: $T_{0}$ is a prototile for $C_{0}=(0, \pm 1)$ if $T_{0}=\varnothing$ and $\mathrm{f}\left(T_{0}\right)=\{x\}$ for any $x \in \mathbb{Z}[\mathbf{i}]$, while $T_{0}$ is a prototile for $C_{0}=(1,1+2 z)$ if $T_{0}=\left\{s_{x}\right\}$ for any $x, z \in \mathbb{Z}[\mathbf{i}]$.

Note that by induction, any circle in $\mathcal{B}$ has at most one prototile up to translation, and any circle's prototiles must necessarily be $180^{\circ}$ symmetric. An example of a decomposition as (7.1) can be seen in large tile in the center of Figure 7.1.

When a prototile $T$ for a circle $C$ satisfies the definition of a tile (i.e., $I(T)$ is a topological disk), we say that $T$ is a tile for $C$. Given a tile $T_{0}$ for a


Figure 7.1. Verifying a tiling by decomposing into parent tiles. The large tile in the center is $T_{0}$.
circle $C_{0} \in \mathcal{B}$ with $c_{0}>1$, we say $T$ is a subtile of $T_{0}$ if $T$ is one of the tiles in the decomposition (7.1) for $T_{0}$. In the proof of Lemma 7.3, we will see that the decomposition of $T$ into subtiles is nonoverlapping, except for prescribed overlap between the largest pair of subtiles.

Our construction in Section 4.2 recursively assigns tiles to each Ford circle with decompositions

$$
T_{0}=T_{1}^{+} \cup T_{1}^{-} \cup T_{2}^{+} \cup T_{2}^{-},
$$

where the $T_{i}^{ \pm}$'s were constructions for the two Ford parents of $C_{0}$. As a general circle in the Apollonian packing, the third parent of a Ford circle is a line $(0,-1)$ with the degenerate tile $\varnothing$; thus to see that $T_{0}$ can be realized as a tile for $C_{0}$ via Definition 7.1, it is only necessary to check, via Proposition 5.4, that assigning $\mathrm{f}\left(T_{3}^{ \pm}\right)=\mathrm{c}\left(T_{0}\right) \pm \frac{1}{2}\left(v_{12}-\mathbf{i} v_{12}\right)$ (from (7.2)) gives that $\mathrm{f}\left(T_{3}^{ \pm}\right) \subseteq \mathrm{f}\left(T_{0}\right)$.

The presence of the coefficient $\frac{1}{2}$ in (7.2) means that even the existence of prototiles for general circles is not quite immediate.

Lemma 7.2. There is a prototile $T_{0}$ for every $C_{0} \in \mathcal{B}$.
Proof. If $C_{0}$ is not equivalent to a Ford circle under the symmetries discussed in Section 5.4, then it is a member of a proper Descartes quadruple


Figure 7.2. Constructing a prototile from the decomposition of its largest parent tile.
$\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$ for $C_{i} \neq(0, \pm 1)$ for each $i$, and $\left(C_{1}, C_{4}, C_{2}, C_{3}\right)$ is also a proper Descartes quadruple, where $C_{4}=2\left(C_{1}+C_{2}+C_{3}\right)-C_{0}$ is the Soddy precursor of $C_{0}$. By induction, $C_{1}$ has a prototile $T_{1}$ falling into one two cases. Either it is a translate of $T_{1}=\left\{s_{0}\right\}$, in which case a direct application of Definition 7.1 verifies that $T_{0}=\left\{s_{-1-\mathbf{i}}, s_{-1}, s_{-\mathbf{i}}, s_{0}\right\}$ is a prototile for $C_{0}$. Otherwise it is given as

$$
T_{1}=S_{4}^{+} \cup S_{4}^{-} \cup S_{2}^{+} \cup S_{2}^{-} \cup S_{3}^{+} \cup S_{3}^{-},
$$

where each $S_{i}^{ \pm}$is a prototile for $C_{i}$ such that

$$
\mathrm{c}\left(S_{i}^{ \pm}\right)-\mathrm{c}\left(T_{1}\right)= \pm \frac{1}{2}\left(v_{k j}-\mathbf{i} v_{k j}\right)
$$

for all rotations $(i, j, k)$ of $(4,2,3)$. In this latter case, we set

$$
\begin{array}{lll}
T_{1}^{+}=T_{1}+v_{32}-\mathbf{i} v_{32}, & T_{2}^{+}=S_{2}^{+}+v_{32}, & T_{3}^{+}=S_{3}^{+}-\mathbf{i} v_{32} \\
T_{1}^{-}=T_{1}, & T_{2}^{-}=S_{2}^{-}-\mathbf{i} v_{32}, & T_{3}^{-}=S_{3}^{-}+v_{32}
\end{array}
$$

See Figure 7.2 for an illustration. The key point is that $v_{32} \in \mathbb{Z}[\mathbf{i}]$, and thus that each $T_{i}^{ \pm}$is a prototile for $C_{i}$, allowing us to define

$$
T_{0}:=T_{1}^{+} \cup T_{1}^{-} \cup T_{2}^{+} \cup T_{2}^{-} \cup T_{3}^{+} \cup T_{3}^{-}
$$

By the lattice rules (Lemma 5.1), we compute that

$$
\mathrm{c}\left(T_{i}^{ \pm}\right)-p_{0}= \pm \frac{1}{2}\left(v_{k j}-\mathbf{i} v_{k j}\right),
$$

if $(i, j, k)$ is a rotation of $(1,2,3)$, and $p_{0}:=c\left(T_{1}^{-}\right)+\frac{1}{2}\left(v_{32}-\mathbf{i} v_{32}\right)$. Since each $T_{i}^{ \pm}$is $180^{\circ}$ symmetric, we must have $p_{0}=\mathrm{c}\left(T_{0}\right)$. Thus $T_{0}$ is a prototile for $C_{0}$ according to Definition 7.1. Note that the condition that $\mathrm{f}\left(T_{i}^{ \pm}\right) \subseteq T_{0}$ cannot fail, since $C_{i} \neq(0, \pm 1)$ for any $i$ implies that no $T_{i}^{ \pm}$is degenerate.

Lemma 7.3 is the heart of our inductive construction, showing that prototiles for circles in $\mathcal{B}$ are in fact $90^{\circ}$ symmetric tiles. Lemma 7.3 concerns a proper Descartes quadruple ( $C_{0}, C_{1}, C_{2}, C_{3}$ ), writing $C_{i}=\left(c_{i}, z_{i}\right)$ and $v_{i j}=$ $v\left(C_{i}, C_{j}\right)$, and where $T_{i}$ denotes some prototile for $C_{i}$. For simplicity of notation, we assume $c_{1} \geq c_{2} \geq c_{3}$, using the symmetries discussed in Section 5.4.


Figure 7.3. The double decomposition of a tile, for the proof of Lemma 7.3.

A major goal of the lemma is to understand the relationship between $T_{0}$ and the tiles $T_{0} \pm v_{0 i}$ for $i=1,2,3$. We do this by decomposing tiles, to deduce that tiles form a touching triple as a consequence of the fact that some smaller tiles form a touching triple. In particular, we write $T_{0}$ as a union of $T_{i}^{ \pm}$as in (7.1) and then define

$$
\begin{array}{lll}
R_{1}^{0}=T_{1}^{+}+v_{03}, & R_{2}^{0}=T_{2}^{+}+v_{01}, & R_{3}^{0}=T_{3}^{+}+v_{02}, \\
R_{1}^{1}=T_{1}^{+}-v_{02}, & R_{2}^{1}=T_{2}^{+}-v_{03}, & R_{3}^{1}=T_{3}^{+}-v_{01}, \\
R_{1}^{2}=T_{1}^{-}-v_{03}, & R_{2}^{2}=T_{2}^{-}-v_{01}, & R_{3}^{2}=T_{3}^{-}-v_{02},  \tag{7.3}\\
R_{1}^{3}=T_{1}^{-}+v_{02}, & R_{2}^{3}=T_{2}^{-}+v_{03}, & R_{3}^{3}=T_{3}^{-}+v_{01} .
\end{array}
$$

(See Figure 7.1.) When $c_{1}>1$, Definition 7.1 implies that the subtiles $T_{1}^{ \pm}$ have decompositions as in (7.1), and we write

$$
\begin{aligned}
& T_{1}^{+}=Q_{2}^{+} \cup Q_{2}^{-} \cup Q_{3}^{+} \cup Q_{3}^{-} \cup Q_{4}^{+} \cup Q_{4}^{-}, \\
& T_{1}^{-}=S_{2}^{+} \cup S_{2}^{-} \cup S_{3}^{+} \cup S_{3}^{-} \cup S_{4}^{+} \cup S_{4}^{-}
\end{aligned}
$$

where $Q_{i}^{ \pm}$and $S_{i}^{ \pm}$are each tiles for $C_{i}$, and $\left(C_{1}, C_{4}, C_{2}, C_{3}\right)$ is the proper Descartes quadruple given by letting $C_{4}=2\left(C_{1}+C_{2}+C_{3}\right)-C_{0} \in \mathcal{B}$ be the Soddy precursor of $C_{0}$. The double decomposition of $T_{0}$ is then the collection of tiles

$$
\left\{T_{2}^{+}, T_{2}^{-}, T_{3}^{+}, T_{3}^{-}\right\} \cup\left\{S_{i}^{ \pm} \mid i=2,3,4\right\} \cup\left\{Q_{i}^{ \pm} \mid i=2,3,4\right\}
$$

shown in Figure 7.3. Note that in the course of proving Lemma 7.3, we will show that $Q_{4}^{-}=S_{4}^{+}$.

Lemma 7.3. Let $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$ be a Descartes quadruple, let $(i, j, k)$ indicate any rotation of $(1,2,3)$, and let $v_{i j}=v\left(C_{i}, C_{j}\right)$. Denoting the tiles of the decomposition and double decomposition of $T_{0}$ as above, the following properties hold:
(T1) $T_{0}$ is a tile, and $\left|T_{0}\right|=c_{0}$. Moreover, $T_{0} \backslash T_{i}^{ \pm}$is a tile, that touches $T_{i}^{ \pm}$, for each $i=1,2,3$.
(T2) $T_{0}$ is $90^{\circ}$ symmetric.
(T3) $T_{0}, T_{0}+v_{i 0}$, and $T_{0}-v_{j 0}$ form a touching triple, provided $c_{0}>1$.
(T4) $T_{0}, T_{0}+v_{i 0}$, and $T_{i}$ form a touching triple whenever $\mathrm{c}\left(T_{i}\right)-\mathrm{c}\left(T_{0}\right)=$ $\frac{1}{2}\left(v_{i 0}+v_{0 i}\right)$ and $c_{0}>1$.
(T5) $T_{0}, T_{i}$, and $T_{j}$ form a touching triple whenever $\mathrm{c}\left(T_{i}\right)-\mathrm{c}\left(T_{0}\right)=\frac{1}{2}\left(v_{i 0}+v_{0 i}\right)$ and $\mathrm{c}\left(T_{j}\right)-\mathrm{c}\left(T_{0}\right)=\frac{1}{2}\left(v_{j 0}-v_{0 j}\right)$.
(T6) If $c_{1} \geq c_{2}>1$, then among other labeled tiles from Figure 7.1, the subtile $T_{i}^{ \pm}$intersects only those that are drawn adjacent to it or with overlap. Moreover, $\mathrm{f}\left(T_{1}^{+}\right) \cap \mathrm{f}\left(T_{2}^{+}\right) \subseteq \mathrm{f}\left(Q_{2}^{+}\right) \cap \mathrm{f}\left(T_{2}^{+}\right)$, and $\mathrm{f}\left(T_{1}^{+}\right) \cap \mathrm{f}\left(T_{1}^{-}\right) \subseteq$ $\left(\mathrm{f}\left(Q_{2}^{+}\right) \cap \mathrm{f}\left(S_{3}^{-}\right)\right) \cup \mathrm{f}\left(Q_{4}^{-}\right) \cup\left(\mathrm{f}\left(Q_{3}^{+}\right) \cap \mathrm{f}\left(S_{2}^{-}\right)\right)$.
Note that (T6) could be removed from Lemma 7.3 without compromising the induction, but this technical information will be necessary for our use of tiles in the construction of integer superharmonic representatives.

Before commencing with our inductive proof of Lemma 7.3 in the general case, we use it to prove the following version of the tiling theorem from the introduction. This will give a very simple example of an application of Lemma 6.1.

Theorem 7.4. For every circle $C \in \mathcal{B}$, there is a tile $T_{C} \subseteq \mathbb{Z}^{2}$ with $90^{\circ}$ rotational symmetry, such that $T_{C}+\Lambda_{C}$ is a tiling. Moreover, except when $C$ has radius 1, each tile in $T_{C}+\Lambda_{C}$ borders exactly six other tiles.

Note that Theorem 1.3 follows from this and our confirmation in Theorem 10.6 that $\Lambda_{C}=L_{C}$.

Proof of Theorem 7.4. In the case where $C \in \mathcal{B}$ has curvature 1, the lattice $\Lambda_{C}$ is $\mathbb{Z}^{2}$ and the corresponding tile $T$ is simply a single square $s_{x}$. Thus we let $T=T_{0}$ be the tile for a circle $C_{0}$ in a Descartes quadruple ( $C_{0}, C_{1}, C_{2}, C_{3}$ ) with $c_{0}>1$, and we write $v_{i j}$ for the vectors $v\left(C_{i}, C_{j}\right)$. We consider a $\mathbb{Z}^{2}$ periodic planar graph $G$ whose vertices correspond to the tiles $T_{0}+\Lambda_{C_{0}}$, where two vertices are adjacent if the corresponding tiles differ by a vector $\pm v_{0 i}$. (T3) and the area condition from (T1) now allow us to verify all the hypotheses of Lemma 6.1. Lemma 6.1 implies that each square of $\mathbb{Z}^{2}$ lies in exactly one tile among $T_{0}+\Lambda_{C_{0}}$ (indeed, $T_{0}+\Lambda_{C_{0}}$ is a tiling) and that tile intersections are in one-to-one correspondence to edges in $G$, which has degree 6 .

The rest of this section is devoted to the proof of Lemma 7.3. Since we have already addressed the case of Ford circles and diamond circles, we may assume that $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}$ is a proper Descartes quadruple, such that no $C_{i}$ is a line, and at most one $C_{i}$ has curvature $c_{i}=1$. In particular, $c_{0}>4$, and, rotating the parents $\left(C_{1}, C_{2}, C_{3}\right)$ and possibly conjugating the original
tuple using the symmetries in Section 5.4 , we may assume $C_{3}$ is a parent of $C_{2}$ and $C_{2}$ is a parent of $C_{1}$ so that $c_{1}>c_{2} \geq 4$. In particular, we also have that $\left(C_{1}, C_{4}, C_{2}, C_{3}\right)$ and $\left(C_{2}, C_{3}, C_{5}, C_{6}\right)$ are both proper Descartes quadruples for some $C_{5}, C_{6} \in \mathcal{B}$ where $C_{4}$ is the Soddy twin of $C_{0}$, and by induction, we may assume Lemma 7.3 holds for both of these quadruples.

For the sake of clarity and brevity, Claims 7.5 and 7.6 are stated with the aid of Figure 7.1 and Figure 7.3, respectively. For the purposes of the statements of the claims, tiles $S_{1}, S_{2}, S_{3}$ (from among the $T_{i}^{ \pm}$and $R_{i}^{j}$, or $T_{i}^{ \pm}$, $Q_{i}^{ \pm}$and $S_{i}^{ \pm}$labeled in the figure) are considered to be drawn adjacently if their corresponding regions in the figure overlap or share some portion of their boundaries as drawn. ( $T_{1}^{-}$and $T_{1}^{+}$in Figure 7.1 are the only tiles drawn with overlap.) In the course of our proof of Lemma 7.3, we will be verifying that the basic tile layout in the figure is correct.

CLAim 7.5. If $\left(S_{1}, S_{2}, S_{3}\right)$ is a triple of adjacent labelled tiles in Figure 7.1, not all of which are contained in $T_{0}$, then the $S_{i}$ form a touching triple, unless $\left|S_{1}\right|=\left|S_{2}\right|=\left|S_{3}\right|=1$, in which case the $S_{i}$ 's are pairwise intersecting squares.

Proof of claim. To prove this claim, we use the lattice rules (5.1) to simplify the differences between centers of adjacent tiles. One can check that adjacent tiles for $C_{i}$ are related by $\pm v_{j i}$ for some $j<i$. One can also check that, when $i \neq j$, adjacent tiles for $C_{i}$ and $C_{j}$ are related by $\frac{1}{2} \mathbf{i}^{s}\left(v_{i j}+v_{j i}\right)$ for $s=0,1,2,3$. From this, we see that all of the triples $\left(S_{1}, S_{2}, S_{3}\right)$ from the claim fall into one of the inductive versions of (T4) or (T5) for ( $C_{1}, C_{4}, C_{2}, C_{3}$ ) or $\left(C_{2}, C_{3}, C_{5}, C_{6}\right)$ so long as $\left|S_{i}\right| \neq 1$ for some $i$. The case where each $S_{i}$ is a tile for a circle of curvature 1 is easily checked by hand.

To analyze touching tiles within $T_{0}$, we need to make use of the double decomposition. The depiction in Figure 7.3 of the double decomposition is a bit deceptive, however, as it shows $S_{4}^{ \pm}$smaller than $S_{3}^{ \pm}$and $S_{2}^{ \pm}$, even though the size of $S_{4}^{ \pm}$relative to $S_{2}^{ \pm}, S_{3}^{ \pm}$is not constrainted by the hypotheses of the lemma. (It depends on the relative size of $C_{4}$ to $C_{2}, C_{3}$.) In particular, this is why the scope of Claim 7.5 is limited to triples of tiles not all lying inside $T_{0}$. To emphasize this point, examples of the six generally possible tangency structures (from the six possible relative size orders of the circles $C_{1}, C_{2}, C_{3}$ ) are illustrated in Figure 7.4.

The following claim is sufficient for our purposes, however, and does not depend on the relative size of $C_{4}$ to $C_{2}$ and $C_{3}$ :

CLAIm 7.6. $S_{4}^{+}=Q_{4}^{-}$, and if $\left(R_{1}, R_{2}, R_{3}\right)$ is a triple of adjacent labelled tiles in Figure 7.3, not all of which are contained in a single $T_{1}^{ \pm}$, then the $R_{i}$ form a touching triple of tiles.


Figure 7.4. The six possible internal tangency structures of the double decomposition. The relative sizes of the parents of the largest parent determine how the subtiles of the largest parents interact with the other parents. Boldfaced lines are boundaries of the first decomposition, while light lines are boundaries of subtiles in the double decomposition. Overlap in the decompositions makes it a little tricky to visually parse some of these cases; it is helpful to keep in mind that each subtile is centrally symmetric and to know that its label is drawn here at its center.

Proof of claim. Again we use the lattice rules (5.1) to simplify the differences between centers of adjacent tiles. In particular, we verify that $S_{4}^{+}=Q_{4}^{-}$ and that each triple covered by the claim is a case of (T4) or (T5) for the quadruple ( $C_{2}, C_{3}, C_{5}, C_{6}$ ).

The information from the double decomposition is not enough for us to claim yet that, e.g., $T_{1}^{+}$and $T_{2}^{+}$touch (in particular, that they intersect only on their boundaries), as we have not analyzed the topological relationships among all the tiles in the double decomposition from Figure 7.3, which would be quite cumbersome in light of the unknown relative size of the circle $C_{4}$. To proceed further at this point will require Lemma 6.1.

In principle, we would like to apply some topology to rule out the presence of extraneous tile relationships. However, Lemma 6.1 is designed to apply to tilings, so the overlap of tiles $T_{1}^{+}$and $T_{1}^{-}$precludes its direct use for this purpose.

To work around this, we will modify some of our tiles so that we actually have a tiling. In particular, we define $\widetilde{\mathcal{T}}_{L}$ by replacing each translate of $T_{1}^{+}$by the corresponding translate of $T_{1}^{+} \backslash Q_{4}^{-}$which, by induction, is a tile by (T1). We consider the graph $G$ whose vertex set is in correspondence with $\mathcal{T}_{L}$ and, also, therefore, $\widetilde{\mathcal{T}}_{L}$. For any vertex $v \in V(G)$, we write $T(v)$ for the corresponding tile in $\mathcal{T}_{L}$, and $\tilde{T}(v)$ for the corresponding tile in $\widetilde{\mathcal{T}}_{L}$. (Note that unless $T(v)$ is a $L$-translate of $T_{1}^{+}$, we have that $T(v)=\tilde{T}(v)$.) A pair of vertices $\{u, v\}$ is joined by an edge in $G$ if $T(u)$ and $T(v)$ are translates by a common vector in $L$ of a tile pair drawn adjacently or with overlap in Figure 7.1. Viewed as an abstract graph, $G$ is easily seen to have a drawing as a $\mathbb{Z}^{2}$-periodic planar triangulation.

We let $\mathcal{E}$ and $\tilde{\mathcal{E}}$ be the sets of pairs $\{T(u), T(v)\}$ and $\{\tilde{T}(u), \tilde{T}(v)\}$, respectively, for $\{u, v\} \in E(G)$. Lemma 6.1 will allow us to prove the following claim:

Claim 7.7. The tiles in $\widetilde{\mathcal{T}}$ form a tiling, the only intersecting pairs of tiles in $\tilde{\mathcal{T}}$ are those in $\mathcal{E}$, and all nonempty intersections of tile-pairs are paths.

Assuming the claim, we can now verify (T1)-(T5) for the quadruple $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$.
(T1) $T_{0}$ is a tile, and $\left|T_{0}\right|=c_{0}$. Moreover, $T_{0} \backslash T_{i}^{ \pm}$is a tile, that touches $T_{i}^{ \pm}$, for each $i=1,2,3$.

Proof. Since $T_{0}$ is a set of squares whose centers form a connected subgraph of the dual lattice $\mathbb{Z}^{2^{*}}$, the fact that $T_{0}$ is a tile follows from the fact that $T_{0}+L$ covers the plane without overlap (so $T_{0}$ has no "holes"), by Lemma 6.1. Similarly, since any five tiles from among

$$
T_{1}^{+} \backslash Q_{4}^{-}, T_{1}^{-}, T_{2}^{+}, T_{2}^{-}, T_{3}^{+}, T_{3}^{-}
$$

induce a connected subgraph of $(\widetilde{\mathcal{T}}, \tilde{\mathcal{E}})$ whose complement is also connected, we get that $T_{0} \backslash T_{i}^{ \pm}$is a tile. It also follows that $\left|T_{0}\right|=c_{0}$, since $c_{0}$ is the determinant of $L$.
(T2) $T_{0}$ is $90^{\circ}$ symmetric.
Proof. Claim 7.5 implies that the tiles $R_{i}^{k}$ for $i=1,2,3$ and $k=1,2,3,4$ surround $T_{0}$. By induction, each of the $R_{i}^{k}$ is $90^{\circ}$ symmetric. Moreover, using the lattice rules (5.1), we see that the union $S=\cup R_{i}^{k}$ is a $90^{\circ}$ symmetric union of squares. Since $T_{0}$ is the bounded component of the complement of $S \backslash \partial S$, it too must be $90^{\circ}$ symmetric.
(T3) $T_{0}, T_{0}+v_{i 0}$, and $T_{0}-v_{j 0}$ form a touching triple provided $c_{0}>1$.
Proof. Our use of Lemma 6.1 tells us that the only pairs of tiles in $\widetilde{\mathcal{T}}$ that intersect are those in $\tilde{\mathcal{E}}$. Thus the pairwise intersections of the tiles $T_{0}, T_{0}+v_{0 i}$, and $T_{0}-v_{0 j}$ can be written as the union of the subtiles from the respective tiles that are drawn adjacently in Figure 7.1. Each of these intersections is a simple path between points $\rho(F), \rho\left(F^{\prime}\right)$ that can be concatenated to a simple path. (The resulting path is simple since the subtiles are all nonoverlapping.) We are done by (T2).
(T4) $T_{0}, T_{0}+v_{i 0}$, and $T_{i}$ form a touching triple if $\mathrm{c}\left(T_{i}\right)-\mathrm{c}\left(T_{0}\right)=\frac{1}{2}\left(v_{i 0}+v_{0 i}\right)$ provided $c_{0}>1$.

Proof. Using the lattice rules (5.1), one can verify for each $i$ in $(1,2,3)$ that $T_{i}=R_{i}^{1}$, as defined in (7.3). In particular, as above, Lemma 6.1 implies that $T_{0}, T_{0}+v_{0 i}$, and $R_{i}^{1}$ form a touching triple. The statement follows now from (T2).
(T5) $T_{0}, T_{i}$, and $T_{j}$ form a touching triple if $\mathrm{c}\left(T_{i}\right)-\mathrm{c}\left(T_{0}\right)=\frac{1}{2}\left(v_{i 0}+v_{0 i}\right)$ and $\mathrm{c}\left(T_{j}\right)-\mathrm{c}\left(T_{0}\right)=\frac{1}{2}\left(v_{j 0}-v_{0 j}\right)$.

Proof. Using the lattice rules (5.1), we can verify that $T_{j}=R_{j}^{2}$, and Lemma 6.1 gives that $T_{0}, R_{i}^{1}$, and $R_{j}^{2}$ form a touching triple. Thus, since we also had that $T_{i}=R_{i}^{1}$, the statement follows now from (T2).
(T6) If $c_{1} \geq c_{2}>1$, then among other labeled tiles from Figure 7.1, the subtile $T_{i}^{ \pm}$intersects only those that are drawn adjacent to it or with overlap. Moreover, $\mathrm{f}\left(T_{1}^{+}\right) \cap \mathrm{f}\left(T_{2}^{+}\right) \subseteq \mathrm{f}\left(Q_{2}^{+}\right) \cap \mathrm{f}\left(T_{2}^{+}\right)$and $\mathrm{f}\left(T_{1}^{+}\right) \cap \mathrm{f}\left(T_{1}^{-}\right) \subseteq\left(\mathrm{f}\left(Q_{2}^{+}\right) \cap \mathrm{f}\left(S_{3}^{-}\right)\right) \cup$ $\mathrm{f}\left(Q_{4}^{-}\right) \cup\left(\mathrm{f}\left(Q_{3}^{+}\right) \cap \mathrm{f}\left(S_{2}^{-}\right)\right)$.

Proof. These are both consequences of the application of Lemma 6.1. The first follows from the fact that $\mathrm{f}\left(T_{1}\right) \cap \mathrm{f}\left(T_{2}\right)$ is a path between $\rho(F)$ and $\rho\left(F^{\prime}\right)$ for the faces $F, F^{\prime}$ that contain $T_{1}, T_{2}$. Similarly, the second follows from the fact that $\mathrm{f}\left(T_{1}^{+} \backslash Q_{4}^{-}\right) \cap \mathrm{f}\left(T_{1}^{-}\right)$is a path in $\left(\mathrm{f}\left(Q_{2}^{+}\right) \cap \mathrm{f}\left(S_{3}^{-}\right)\right) \cup \mathrm{f}\left(Q_{4}^{-}\right) \cup\left(\mathrm{f}\left(Q_{3}^{+}\right) \cap \mathrm{f}\left(S_{2}^{-}\right)\right)$.

Remark 7.8. An essential feature of (T3), (T4), (T5) is that their proofs make use of $90^{\circ}$ symmetry. In particular, each induction step essentially rotates known relationships by 90 degrees, which is then fixed by using (T2). The lack of $90^{\circ}$ symmetry will force a significant increase in the complexity of our odometer construction in comparison with this tile induction.

We now complete the proof of Lemma 7.3 with the following:
Proof of Claim 7.7. We claim that the graph $\left(\widetilde{\mathcal{T}}_{L}, \tilde{\mathcal{E}}\right)$ satisfies the hypotheses of Lemma 6.1. Hypothesis 1 is immediate, as is the first part of Hypothesis 2. For the second part of Hypothesis 2, observe that (using (1.7)) we have

$$
\sum_{T \in \widetilde{\mathcal{T}}_{L} / L}|T|=2 c_{1}+2 c_{2}+2 c_{3}-c_{4}=c_{0}=\left|\operatorname{det} \Lambda_{C_{0}}\right|
$$

It remains to verify Hypotheses 3 and 4 . We will need the following:
Claim 7.9. If $s_{x} \notin T_{1}^{-} \cup T_{1}^{+}$, then $s_{x} \cap \mathrm{f}\left(Q_{4}^{-}\right)=\varnothing$, unless $\left|T_{3}^{-}\right|=\left|T_{2}^{+}\right|=1$.
Proof. By (T1), $T_{1}^{+} \backslash Q_{4}^{-}$and $T_{1}^{-} \backslash S_{4}^{+}$are both tiles that each touch $T_{4}:=Q_{4}^{-}=S_{4}^{+}$. Thus neither of these tiles shares any squares with $S_{4}^{+}$, and unless $\left|T_{3}^{-}\right|=\left|T_{2}^{+}\right|=1$, Claim 7.6 or inductively, by Lemma 7.3 , give that the intersection of $T_{1}^{+} \backslash Q_{4}^{-}$and $T_{1}^{-} \backslash S_{4}^{+}$contains edges of $\mathbb{Z}^{2}$. Thus, the union $\mathcal{S} \subseteq \mathbb{Z}^{2}$ of these tiles satisfies that its dual $\mathcal{S}^{*}$ is connected; since $\mathcal{S}$ is centrally symmetric about the center of $S_{4}^{+}$, we have as a consequence that the squares of $\mathcal{S}$ surround $T_{4}$. In particular, no square of $S_{4}^{+}$can intersect any square outside of $\mathcal{S}$.

Consider now hypothesis (3) of Lemma 6.1. For any pair $\tilde{T}(u), \tilde{T}(v)$ in $\tilde{\mathcal{E}}$, we have from the definition of $G$ that the pair $T(u), T(v)$ is drawn adjacently or with overlap in Figure 7.1. Claims 7.5 and 7.6 now imply that $f(T(u)) \cap$ $\mathrm{f}(T(v))$ contains at least two vertices. Finally, Claim 7.9 implies that $\mathrm{f}(\tilde{T}(u)) \cap$ $\mathrm{f}(\tilde{T}(v))=\mathrm{f}(T(u)) \cap \mathrm{f}(T(v))$, unless, up to symmetry, we have $T(u)=T_{1}^{+}$, $T(v)=T_{1}^{-}$. In this case, however, Claim 7.6 implies directly that $\mathrm{f}(\tilde{T}(u)) \cap$ $\mathrm{f}(\tilde{T}(v))$ contains at least two vertices.

To check Hypothesis 4 for the graph $(\widetilde{\mathcal{T}}, \tilde{\mathcal{E}})$, we warm up by examining this for the graph $(\mathcal{T}, \mathcal{E})$, where it also holds. We choose an assignment of the points $\rho(F)$ for faces $F$ of $G$. For any face $F=\{u, v, w\}$ whose three corresponding tiles form a triple covered by Claim 7.5, we must choose $\rho(F) \in$ $\mathbb{Z}[\mathbf{i}]$ to be the unique point in the three-way intersection of the footprints of the tiles. Remaining faces $F$ are those whose corresponding tile triple lies entirely within $T_{0}$ : i.e., the triple is either $\left(T_{1}^{ \pm}, T_{2}^{ \pm}, T_{3}^{\mp}\right)$ or $\left(T_{1}^{ \pm}, T_{1}^{\mp}, T_{2}^{ \pm}\right)$. In the first case, we use Claim 7.6 to choose (without loss of generality) $\rho\left(\left\{T_{1}^{+}, T_{2}^{+}, T_{3}^{-}\right\}\right.$) to be the unique point in the intersection of the touching triple $Q_{2}^{+}, T_{2}^{+}, T_{3}^{-}$.

In the second case, we use Claim 7.6 to choose (without loss of generality) $\rho\left(\left\{T_{1}^{+}, T_{1}^{-}, T_{2}^{+}\right\}\right)$to be the unique point in the intersection of the touching triple $Q_{2}^{+}, T_{2}^{+}, S_{3}^{-}$.

With this choice for the $\rho(F)$ 's, we see that for any adjacent $F, F^{\prime}$ corresponding to tiles from Figure 7.1 other than the pair of faces corresponding to the tile triples

$$
\begin{equation*}
\left\{T_{1}^{+}, T_{1}^{-}, T_{2}^{+}\right\}, \quad\left\{T_{1}^{+}, T_{1}^{-}, T_{2}^{-}\right\} \tag{7.4}
\end{equation*}
$$

the points $\rho(F), \rho\left(F^{\prime}\right)$ are joined by a simple path given as the intersection of a single pair of tiles $\mathrm{f}\left(U^{1}\right) \cap \mathrm{f}\left(U^{2}\right)$ known to touch either by Claim 7.5 or Claim 7.6. This would verify Hypothesis 4 of Lemma 6.1 in the graph $(\mathcal{T}, \mathcal{E})$ for these pairs of faces. Note now that Claim 7.9 implies that the points $\rho(F)$ selected above also lie in the three way intersections $\mathrm{f}(\tilde{T}(u)) \cap \mathrm{f}(\tilde{T}(v)) \cap \mathrm{f}(\tilde{T}(w))$ for $\{u, v, w\}=F$ and, moreover, that the simple paths in the intersections of the footprints of pairs of tiles $T(u), T(v)$ used above remain in the intersection $\mathrm{f}(\tilde{T}(u)) \cap \mathrm{f}(\tilde{T}(v))$. In particular, we have verified that with this assignment of $\rho(F)$ 's, Hypothesis 4 holds also in the graph $(\widetilde{\mathcal{T}}, \tilde{\mathcal{E}})$ for all faces other than those corresponding to translates of triples

$$
\left\{T_{1}^{+} \backslash Q_{4}^{-}, T_{1}^{-}, T_{2}^{+}\right\}, \quad\left\{T_{1}^{+} \backslash Q_{4}^{-}, T_{1}^{-}, T_{2}^{-}\right\}
$$

and we deal with this final case separately now.
Let $L^{+}$be the set of edges of $\partial Q_{4}^{-}$whose two incident squares both lie in $T_{1}^{+}$, and define $L^{-}$similarly. Note that $L^{-}$is the central reflection through $\mathrm{c}\left(Q_{4}^{-}\right)$of $L^{+}$and that Claim 7.9 implies that $L^{+} \cup L^{-}$equals the edge-set of $\partial Q_{4}^{-}$. We claim that $L^{+}$is the edge-set of a path, which suffices for us since then, letting $V\left(L^{+}\right)$denote the set of vertices incident with edges in $L^{+}$, we have that $\mathrm{f}\left(S_{3}^{-}\right) \cap \mathrm{f}\left(S_{4}^{+}\right) \cap \mathrm{f}\left(Q_{2}^{+}\right) \subseteq V\left(L^{+}\right) \cap V\left(L^{-}\right)$and $\mathrm{f}\left(S_{2}^{-}\right) \cap \mathrm{f}\left(S_{4}^{+}\right) \cap \mathrm{f}\left(Q_{3}^{+}\right) \subseteq$ $V\left(L^{+}\right) \cap V\left(L^{-}\right)$, so that

$$
\left(S_{3}^{-} \cap Q_{2}^{+}\right) \cup V\left(L^{+}\right) \cup\left(Q_{3}^{+} \cap S_{2}^{-}\right)
$$

lies in the intersection $\left(T_{1}^{+} \backslash Q_{4}^{-}\right) \cap T_{1}^{-}$and contains a path joining the points assigned to the triples from (7.4).

To see that this $L^{+}$is indeed the edge-set of a path, suppose it is false. In particular, we can, in cyclic order in $\partial Q_{4}^{-}$, find edges $e_{1}, e_{2}, e_{3}, e_{4}$ such that $e_{1}, e_{3} \in L^{+}$and $e_{2}, e_{4} \in L^{-} \backslash L^{+}$. But then since $T_{1}^{+} \backslash Q_{4}^{-}$is connected, there is a path of squares in $T_{1}^{+} \backslash Q_{4}^{-}$from $e_{1}$ to $e_{3}$. But, together with the squares of $Q_{4}^{-}$, these squares must enclose either $s_{2}$ or $s_{4}$, where $s_{i}$ is the square incident with $e_{i}$ lying outside of $Q_{4}^{-}$. In particular, $T_{1}^{+} \backslash Q_{4}^{-}$cannot be simply connected, contradicting the inductive hypothesis (T1).

Thus the hypotheses of Lemma 6.1 are satisfied for $(\widetilde{\mathcal{T}}, \tilde{\mathcal{E}})$.

Remark 7.10. We emphasize that our proof of Claim 7.7 did not work by showing that the intersections of adjacent tiles in $\widetilde{\mathcal{T}}_{L}$ are all simple paths. In particular, for the tile pair $\left(T_{1}^{+} \backslash Q_{4}^{-}, T_{1}^{-}\right)$, the proof only establishes that the intersection contains a path. The problem with using induction to establish that the intersection equals a path for this pair is that the tile relationships in the double decomposition (e.g., $S_{3}^{-}, S_{4}^{+}, Q_{2}^{+}$, etc.) are not fixed and depend on the relative order of the sizes of the circles $C_{1}, C_{2}, C_{3}, C_{4}$. Instead, the proof given above uses an ad hoc argument for this pair based on Claim 7.9. In principle, one could prove a stronger version of Claim 7.7 by considering several cases according to the relative sizes of the circles and accounting for the various subtile-relationships that arise in each case. This would allow one to apply a more straigthforward topological lemma at the expense of a blow-up in the tile analysis.

## 8. Boundary strings

Having constructed tiles, we now turn our attention towards constructing odometers. The main idea is to mimic the tile construction, using the duality between the rules for the $v_{i j}$ and $a_{i j}$ in (5.1) to attach superharmonic data to the tiles. However, there is a problem in directly lifting the tile construction: the odometers are only $180^{\circ}$ symmetric in general, and we used $90^{\circ}$ symmetry in constructing the tiles.

An examination of the argument in Section 7 suggests that we made essential use of $90^{\circ}$ symmetry only to conclude that the lattice $\Lambda_{C}$ generated by $\left\{v_{10}, v_{20}, v_{30}\right\}$ is a tiling lattice for $T_{0}$ from the fact that $\mathbf{i} \Lambda_{C}$ is a tiling lattice for $T_{0}$ (Remark 7.8). To give a proof using only $180^{\circ}$ symmetry, we must therefore find a way to express the interface between $T_{0}$ and $T_{0}+v_{i 0}$ directly in terms of relationships between subtiles.

We need a new type of tile decomposition, which we call boundary strings. If $\mathcal{T}$ is a regular tiling, then we call a sequence $T^{0}, T^{1}, \ldots, T^{n}$ of tiles in $\mathcal{T}$ the $\mathcal{T}$-string (or simply: string) from $T^{0}$ to $T^{n}$ if
(1) $T^{i}$ touches $T^{i+1}$ for $0 \leq i<n$;
(2) each $\mathrm{c}\left(T^{i}\right)(0<i<n)$ lies in the closed half-plane to the left of the ray from $\mathrm{c}\left(T^{0}\right)$ to $\mathrm{c}\left(T^{n}\right)$ - i.e.,

$$
\operatorname{Im}\left(\frac{\mathrm{c}\left(T^{i}\right)-\mathrm{c}\left(T^{0}\right)}{\mathrm{c}\left(T^{n}\right)-\mathrm{c}\left(T^{0}\right)}\right) \geq 0
$$

(3) each $T^{i}(0<i<n)$ touches some $S^{i} \in \mathcal{T}$ whose centroid lies outside of the closed half-plane to the left of the ray from $\mathrm{c}\left(T^{0}\right)$ to $\mathrm{c}\left(T^{n}\right)$.
A $\mathcal{T}$-string is the left-handed approximation of the line segment from $\mathrm{c}\left(T^{0}\right)$ to $\mathrm{c}\left(T^{n}\right)$ in $\mathcal{T}$; see Figure 8.1. It is not hard to show from the definition that there is a unique $\mathcal{T}$-string between any two tiles in a regular tiling (indeed, the


Figure 8.1. A $\mathcal{T}$-string $T^{0}, \ldots, T^{7}$ in the regular tiling $\mathcal{T}=$ $T_{C}+\Lambda_{C}$ for the circle $C=(28,7+20 \mathbf{i})$.
partial sequences $T^{0}, T^{1}, \ldots, T^{j}$ are uniquely determined, inductively, given $T_{0}$ and the ray from $\mathrm{c}\left(T^{0}\right)$ towards $\left.\mathrm{c}\left(T^{n}\right)\right)$. The interior tiles of the $\mathcal{T}$-string $\mathcal{S}$ from $T^{0}$ to $T^{n}$ are the tiles in the $\mathcal{T}$-string other than the endpoints $T^{0}, T^{n}$, and the interior of the $\mathcal{T}$-string is the union of the footprints of all interior tiles.

Given a tile $T_{0}$ for the circle $C_{0}$ in the proper Descartes quadruple ( $C_{0}, C_{1}$, $\left.C_{2}, C_{3}\right)$, the $C_{i}$ boundary-string $(i=1,2,3)$ for the tile $T_{0}$ is the string from the tiles $R_{i}^{-}$to $R_{i}^{+}$for the tiling associated to $C_{i}$, where

$$
\mathrm{c}\left(R_{i}^{ \pm}\right)=\mathrm{c}\left(T_{0}\right)+\frac{1}{2}\left(v_{i 0} \pm v_{0 i}\right) .
$$

The following lemma shows that the boundary strings for a tile can be constructed by concatenating certain smaller strings together.

Lemma 8.1. Let $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ be a proper Descartes quadruple, write $v_{i j}=v\left(C_{i}, C_{j}\right)$, suppose $T_{0}$ is a tile for $C_{0}$ with the tile decomposition $\left\{T_{i}^{ \pm}\right\}$, and suppose $i \in\{1,2,3\}, c_{i}>0$, and that $R_{i}^{ \pm}$is a tile for $C_{i}$ satisfying

$$
\mathrm{c}\left(R_{i}^{ \pm}\right)=\mathrm{c}\left(T_{0}\right)+\frac{1}{2}\left(v_{i 0} \pm v_{0 i}\right) .
$$

Then we have that

$$
\begin{equation*}
\mathrm{c}\left(R_{i}^{+}\right)-\mathrm{c}\left(T_{i}^{-}\right)=v_{k i}, \quad \mathrm{c}\left(T_{i}^{-}\right)-\mathrm{c}\left(R_{i}^{-}\right)=-v_{j i}, \tag{8.1}
\end{equation*}
$$

and the $\left(R_{i}^{-}+\Lambda_{C_{i}}\right)$-string from $R_{i}^{-}$to $R_{i}^{+}$is the concatenation of the $\left(R_{i}^{-}+\Lambda_{C_{i}}\right)$ strings from $R_{i}^{-}$to $T_{i}^{-}$and from $T_{i}^{-}$to $R_{i}^{+}$, respectively.

Note that (8.1) implies via Lemma 5.3 that the tiles $T_{i}^{-}$and $R_{i}^{+}$both lie in tiling $R_{i}^{-}+\Lambda_{C_{i}}$, ensuring that the referenced strings are well defined.

It is not hard at this point to use the lattice rules to strengthen Lemma 8.1, to show inductively that the strings from $R_{i}^{-}$to $T_{i}^{-}$and from $T_{i}^{-}$to $R_{i}^{+}$are


Figure 8.2. The three boundary strings of a tile.
themselves boundary strings of subtiles (when they have more than two tiles). This induction also gives, for example, that the interior tiles of a boundary string for $T_{0}$ lie in $T_{0}$, as seen in Figure 8.2. We postpone this calculation until the next section, however, when we are prepared to simultaneously show that the strings have important compatibility properties with respect to our odometer construction.

Proof of Lemma 8.1. The offsets in (8.1) result from a straightforward calculation using (5.1); for convenience, note that referring to the same tile collection (7.3) used in the proof of Lemma 7.3 (see Figure 7.1), one can check that $R_{i}^{+}=R_{i}^{1}, R_{i}^{-}=R_{i}^{0}$ in that decomposition.

Examining the definition of a string, we see that the statement regarding the concatenation fails only if the interior of the triangle $\Delta \mathrm{c}\left(T_{i}^{-}\right) \mathrm{c}\left(R_{i}^{-}\right) \mathrm{c}\left(R_{i}^{+}\right)$ contains the center of some tile in the tiling $T_{i}^{-}+\Lambda_{C_{i}}$. However, the lattice generated by $-v_{j i}$ and $v_{k i}$ has determinant

$$
\frac{1}{2}\left(\overline{v_{i k}} v_{j i}+v_{i k} \overline{v_{j i}}\right)=c_{i}
$$

by the lattice rules (5.1). In particular, the triangle $\triangle \mathrm{c}\left(T_{i}^{-}\right) \mathrm{c}\left(R_{i}^{-}\right) \mathrm{c}\left(R_{i}^{+}\right)$has area $\frac{1}{2} c_{i}$. Thus the lemma follows from the fact that any triangle of area half the determinant of a lattice containing its three vertices can contain no other points of the lattice. (This is a special case of Pick's theorem, for example.)

The following topological lemma allows us to analyze tile interfaces and, in particular, odometer interfaces - using only $180^{\circ}$ symmetry.

Lemma 8.2. If $R$ and $S$ are tiles in the tiling $\mathcal{T}=T_{0}+\Lambda_{C_{0}}$ corresponding to the circle $C_{0} \in \mathcal{B}$ with $c_{0} \geq 1$, then the intersection of the interiors of the
$\mathcal{T}$-string $\mathcal{R}$ from $R$ to $S$ and the $\mathcal{T}$-string $\mathcal{S}$ from $S$ to $R$ contains a path in $\mathbb{Z}[\mathbf{i}]$ from $R$ to $S$.

Proof. We consider the graph $G$ on $\mathcal{T}$ where $T, T^{\prime}$ are adjacent if they touch. $G$ is isomorphic to the graph of the triangular lattice unless $c_{0}=1$, in which case the lemma is easy to verify directly.

For the former case, we draw $G$ in the plane by placing each vertex at the center of the corresponding tile and drawing straight line segments between adjacent vertices of $G$; the result is a planar embedding of $G$ that is affineequivalent to the equilateral triangle embedding of the triangular lattice.

We now consider the sequences $F^{0}, F^{1}, \ldots, F^{k}$ of faces of $G$ through which the line segment from $\mathrm{c}(R)$ to $\mathrm{c}(S)$ passes, in order. If we associate to each $F$ the point $\rho(F)$ in $\mathbb{Z}[\mathbf{i}]$ lying in the intersection of the three tiles of $F$, then the consecutive faces $F^{i}$ and $F^{i+1}$ either share an edge corresponding to touching tiles $T \in \mathcal{R}, T^{\prime} \in \mathcal{S}$, or they share a vertex corresponding to a tile in both $\mathcal{R}$ and $\mathcal{S}$; in either case, there is a path in $\mathbb{Z}[\mathbf{i}]$ from $\rho(F)$ to $\rho\left(F^{\prime}\right)$ lying in the intersection of the interiors of $\mathcal{S}$ and $\mathcal{R}$. Concatenating these paths consecutively, we get a walk in $\mathbb{Z}[\mathbf{i}]$ from $\rho\left(F^{0}\right)$ to $\rho\left(F^{k}\right)$, lying in the intersection of the interiors of $\mathcal{R}$ and $\mathcal{S}$.

## 9. Tile odometers

In this section we attach function data to our tiles. Since we are no longer concerned with topological issues, our definition of a tile as a set of squares $s_{x}$ is no longer useful, and from here on we identify a tile $T$ with the vertex set of its footprint graph, $\bigcup_{s_{x} \in T} s_{x} \subset \mathbb{Z}[\mathbf{i}]$. In particular, $T \cap T^{\prime}$ now denotes a subset of $\mathbb{Z}[\mathbf{i}]$ rather than a set of squares.
9.1. Basic definitions. A partial odometer is a function $h: T \rightarrow \mathbb{Z}$ with a finite domain $T \subseteq \mathbb{Z}^{2}$. We write $T(h)$ for the domain of $h$ and $\mathbf{s}(h) \in \mathbb{C}$ for the slope of $h$, which is the average of

$$
\begin{align*}
& \frac{1}{2}(h(x+1)-h(x)+h(x+1+\mathbf{i})-h(x+\mathbf{i})) \\
& \quad+\frac{\mathbf{i}}{2}(h(x+\mathbf{i})-h(x)+h(x+1+\mathbf{i})-h(x+1)) \tag{9.1}
\end{align*}
$$

over squares $\{x, x+1, x+\mathbf{i}, x+1+\mathbf{i}\} \subseteq T$; this is a measure of an average gradient for $h$. Note that the slope is not defined when $T$ is a singleton. We say that two partial odometers $h_{1}$ and $h_{2}$ are translations of one another if

$$
\begin{equation*}
T\left(h_{1}\right)=T\left(h_{2}\right)+v \quad \text { and } \quad h_{1}(x)=h_{2}(x+v)+a \cdot x+b \tag{9.2}
\end{equation*}
$$

for some $v, a \in \mathbb{Z}^{2}$ and $b \in \mathbb{Z}$.

Definition 9.1. We say that two partial odometers $h_{1}$ and $h_{2}$ are compatible if $h_{2}-h_{1}=c$ on $T\left(h_{1}\right) \cap T\left(h_{2}\right)$ for some constant $c$, which we call the offset constant for the pair ( $h_{1}, h_{2}$ ), or if $T\left(h_{1}\right) \cap T\left(h_{2}\right)=\varnothing$.

When the offset constant is 0 , or in the second case, we write $h_{1} \cup h_{2}$ for the common extension of the $h_{i}$ to the union of their domains. The next lemma allows us to glue together pairwise compatible partial odometers. Recall we have defined a tiling as a collection of tiles $\mathcal{T}$ such that every square $s_{x}=$ $\{x, x+1, x+i, x+1+i\}$ of $\mathbb{Z}^{2}$ is contained in exactly one element of $\mathcal{T}$.

Lemma 9.2. If $\mathcal{H}$ is a collection of pairwise compatible partial odometers such that $\mathcal{T}=\{T(h): h \in \mathcal{H}\}$ is a hexagonal tiling, then there is a function $g: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$, unique up to adding a constant, that is compatible with every $h \in \mathcal{H}$.

Proof. Since $\mathcal{T}$ is a hexagonal tiling, its intersection graph $G$ is a planar triangulation, and each face $\left\{h_{0}, h_{1}, h_{2}\right\}$ of $G$ corresponds to a touching triple of tiles $T\left(h_{i}\right)$ with $T\left(h_{0}\right) \cap T\left(h_{1}\right) \cap T\left(h_{2}\right) \neq \varnothing$. In particular, writing $d\left(h_{i}, h_{j}\right)$ for the constant value of $h_{i}-h_{j}$ on $T\left(h_{i}\right) \cap T\left(h_{j}\right)$, we have $d\left(h_{0}, h_{1}\right)+d\left(h_{1}, h_{2}\right)+$ $d\left(h_{2}, h_{1}\right)=0$. It follows that the sum of $d$ over any cycle is zero and therefore that $d$ can be written as the gradient of a vertex function $f: \mathcal{T} \rightarrow \mathbb{Z}$ that is unique up to additive constant. Now fix any $h_{0} \in \mathcal{H}$ and set $g(x)=h_{0}(x-y)$ $+f(T(h))$ on $x \in T(h)$, where $y \in \mathbb{Z}^{2}$ is the translation such that $T(h)=$ $T\left(h_{0}\right)+y$.

Our goal is to associate a partial odometer $h$, unique up to odometer translation (9.2), to every circle $C \in \mathcal{B}$.

Definition 9.3. A partial odometer $h_{0}: T_{0} \rightarrow \mathbb{Z}$ is a tile odometer for $C_{0} \in \mathcal{B}$ if $T_{0}$ is a tile for $C_{0}$ and either

- $C_{0}=(0, \pm 1)$ (so $T_{0}$ is a singleton);
- $C_{0}=(1,1+2 z)$ for some $z \in \mathbb{Z}[\mathbf{i}]$, and $h$ is any translation of the partial odometer $h^{\prime}:\{0,1, \mathbf{i}, 1+\mathbf{i}\} \rightarrow \mathbb{Z}$ given by $h(0)=h(1)=h(\mathbf{i})=0, h(1+\mathbf{i})=$ $\operatorname{Im}(z) / 2$; or
- $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ is a proper Descartes quadruple, $C_{i}=\left(c_{i}, z_{i}\right), a_{i j}=$ $a\left(C_{i}, C_{j}\right), v_{i j}=v\left(C_{i}, C_{j}\right)$, and

$$
\begin{equation*}
h_{0}=h_{1}^{+} \cup h_{1}^{-} \cup h_{2}^{+} \cup h_{2}^{-} \cup h_{3}^{+} \cup h_{3}^{-}, \tag{9.3}
\end{equation*}
$$

where $h_{i}^{ \pm}$is a tile odometer for $C_{i}$ such that

$$
\mathrm{c}\left(T\left(h_{i}^{ \pm}\right)\right)-\mathrm{c}\left(T\left(h_{0}\right)\right)= \pm \frac{1}{2}\left(v_{k j}-\mathbf{i} v_{k j}\right)
$$

and

$$
c_{i}=0 \quad \text { or } \quad \mathbf{s}\left(h_{i}^{ \pm}\right)-\mathbf{s}\left(h_{0}\right)= \pm \frac{1}{2}\left(a_{k j}+\mathbf{i} a_{k j}\right)
$$

for all rotations $(i, j, k)$ of $(1,2,3)$.

We call the $h_{i}^{ \pm}$'s the subodometers of $h_{0}$.
When $T_{1}, \ldots, T_{n}$ is a sequence of tiles such $T_{k}$ is a subtile of $T_{k+1}$, we call $T_{1}$ an ancestor tile of $T_{n}$. Restrictions $h \mid T$ of odometers to ancestor tiles $T$ of $T(h)$ are called ancestor odometers of $h_{0}$. The following lemma asserts that this makes sense:

Lemma 9.4. If $h_{0}: T_{0} \rightarrow \mathbb{Z}$ is a tile odometer for $C_{0}$ and the tile $T$ for the circle $C \in \mathcal{B}$ is an ancestor of $T_{0}$, then $h_{0} \mid T$ is a tile odometer for $C$.

By induction, tile odometers are easily seen to be centrally symmetric:
Lemma 9.5. If $h$ is a tile odometer with domain $T$, then $x \mapsto h(-x)$ with domain $-T$ is a translation of $h$.

Our main goal in this section is to prove inductively that each circle $C \in \mathcal{B}$ has an associated tile odometer. (Note that, inductively, the definition immediately gives that tile odometers for a given circle are unique up to odometer translation.) However, before proving that circles in $\mathcal{B}$ do have tile odometers, we will prove that tile odometers must have certain compatibility properties when they do exist.

Given a proper Descartes quadruple $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$ with $c_{0}>0$, we say that tile odometers $h_{0}$ and $h_{0}^{\prime}$ for $C_{0}$ are left-lattice adjacent if

$$
\begin{aligned}
\mathrm{c}\left(T\left(h_{0}^{\prime}\right)\right)-\mathrm{c}\left(T\left(h_{0}\right)\right) & = \pm v_{i 0}, \\
\mathrm{~s}\left(h_{0}^{\prime}\right)-\mathbf{s}\left(h_{0}\right) & = \pm a_{i 0}
\end{aligned}
$$

for some $i \in\{1,2,3\}$ (with matching signs, as usual) and similarly right-lattice adjacent if

$$
\begin{aligned}
\mathrm{c}\left(T\left(h_{0}^{\prime}\right)\right)-\mathrm{c}\left(T\left(h_{0}\right)\right) & = \pm v_{0 i}, \\
\mathrm{~s}\left(h_{0}^{\prime}\right)-\mathrm{s}\left(h_{0}\right) & = \pm a_{0 i}
\end{aligned}
$$

for some $i \in\{1,2,3\}$. Given tile odometers $h_{0}$ and $h_{i}$ for $C_{0}$ and $C_{i}(i=1,2,3)$, we say that $h_{i}$ is subtile-lattice adjacent to $h_{0}$ if

$$
\begin{aligned}
\mathrm{c}\left(T\left(h_{i}\right)\right)-\mathrm{c}\left(T\left(h_{0}\right)\right) & =\mathbf{i}^{s} \frac{1}{2}\left(v_{i 0}+v_{0 i}\right), \\
\mathbf{s}\left(h_{i}\right)-\mathbf{s}\left(h_{0}\right) & =(-\mathbf{i})^{s} \frac{1}{2}\left(a_{i 0}+a_{0 i}\right)
\end{aligned}
$$

for some $s \in\{0,1,2,3\}$. Note that the subtile-lattice adjacency relationship is not symmetric.
9.2. Outline of construction. The essential difference between the induction in Section 7 and the argument we are forced to carry out in this section is that the odometer analog of (T2) in Lemma 7.3 is false; tile odometers are not $90^{\circ}$ symmetric in a straightforward way. In particular, we cannot build a $\Lambda_{C}$-periodic global odometer using that inductive argument. (Instead, that
argument would give an $\mathbf{i} \Lambda_{C}$ periodic function, which does not have the correct growth to be an odometer for $C$. These $\mathbf{i} \Lambda_{C}$ functions are interesting in their own right, however.)

Note that Definition 9.3 is analogous to Definition 7.1 and that left-lattice adjacency and subtile-lattice adjacency as defined here correspond for the domains of $h_{0}$ and $h_{t}^{\prime}(t=0,1,2,3)$ to the two types of pairwise tile relationships that were seen in (T3) and (T4) of Lemma 7.3. However, as the proof of (T4) ends with an application of $90^{\circ}$ symmetry, the inductive argument of Section 7 can actually be adapted to the cases of right-lattice adjacency and subtile-lattice adjacency of tile-odometers, but not left-lattice adjacency. This adaptation is carried out in Lemma 9.6 below, which shows that rightlattice adjacency and subtile-lattice adjacency do indeed reduce, inductively, to (subtile- or left-lattice) adjacencies among smaller tile odometers.

The main innovation in this section is then to make use of boundary strings to establish the compatibility of left-lattice adjacent tile odometers (Lemmas 9.7 and 9.8). Once left-lattice adjacency is handled, the reduction from the adapatation of the tile argument described above implies that right-lattice and subtile-lattice adjacent tile odometers are also compatible (Lemma 9.9). Finally, this allows us to prove that tile odometers exist for each circle (Lemma 9.10).
9.3. Adapting the tile argument. Here we adapt the tile argument to reduce right- and subtile-lattice adjacency to subtile- and left-lattice adjacency among smaller tile odoemters.

Lemma 9.6. If the tile odometers $h$ and $h_{0}$ are right-lattice adjacent, then any subodometer of $h$ intersecting $h_{0}$ is subtile-lattice adjacent to $h_{0}$. Similarly, if the tile odometer $h$ is subtile-lattice adjacent to $h_{0}$, then for any subodometer $h_{i}^{ \pm}$of $h_{0}$ that intersects $h$, either $h, h_{i}^{ \pm}$are left-lattice adjacent, or $h$ is subtilelattice adjacent to $h_{i}^{ \pm}$, or $h_{i}^{ \pm}$is subtile-lattice adjacent to $h$.

Proof. Recall the tiles $R_{i}^{t}$ defined in (7.3) from the proof of Lemma 7.3, as shown in Figure 7.1; each $R_{i}^{t}$ is defined by

$$
\begin{aligned}
& R_{i}^{0}=T_{i}^{+}+v_{0 k}, \\
& R_{i}^{1}=T_{i}^{+}-v_{0 j}, \\
& R_{i}^{2}=T_{i}^{-}-v_{0 k}, \\
& R_{i}^{3}=T_{i}^{-}-v_{0 j},
\end{aligned}
$$

where $(i, j, k)$ is a rotation of $(1,2,3)$. We define for each $R_{i}^{t}$ a translation $h_{i}^{t}$ of the tile odometer for the circle $C_{i}$ whose domain is $R_{i}^{t}$, by

$$
\mathrm{c}\left(h_{i}^{0}\right)=\mathrm{c}\left(h_{i}^{+}\right)+v_{0 k}, \quad \mathrm{~s}\left(h_{i}^{0}\right)=\mathbf{s}\left(h_{i}^{+}\right)+a_{0 k},
$$

$$
\begin{array}{ll}
\mathrm{c}\left(h_{i}^{1}\right)=\mathrm{c}\left(h_{i}^{+}\right)-v_{0 j}, & \mathrm{~s}\left(h_{i}^{1}\right)=\mathrm{s}\left(h_{i}^{+}\right)-a_{0 j} \\
\mathrm{c}\left(h_{i}^{2}\right)=\mathrm{c}\left(h_{i}^{-}\right)-v_{0 k}, & \mathrm{~s}\left(h_{i}^{2}\right)=\mathrm{s}\left(h_{i}^{-}\right)-a_{0 k} \\
\mathrm{c}\left(h_{i}^{3}\right)=\mathrm{c}\left(h_{i}^{-}\right)+v_{0 j}, & \mathrm{~s}\left(h_{i}^{3}\right)=\mathrm{s}\left(h_{i}^{-}\right)+a_{0 j}
\end{array}
$$

We know from (T4) and the application in Section 7 of Lemma 6.1 that the boundary of $T_{0}$ is covered by the $R_{i}^{t}$ 's. Thus, to prove the first part of the lemma, it is sufficient to show compatibility of $h_{0}$ with the $h_{i}^{t}$. And by part (T6) of Lemma 7.3, it is sufficient to show compatibility of tile odometers $h_{i}^{ \pm}$and $h_{s}^{t}$ whose domains are drawn as touching in Figure 7.1. It can be checked by hand using the lattice rules (5.1) that such a pair $h_{i}^{ \pm}$and $h_{s}^{t}$ are either left-lattice adjacent or subtile-lattice adjacent, proving the full statement of the lemma, since the $h_{s}^{t}$ include all odometers that are subtile-lattice adjacent to $h_{0}$.

Looking ahead, if we knew left-lattice adjacent tile odometers to be compatible, then by induction, Lemma 9.6 would give compatibility of right-lattice adjacent and subtile-lattice adjacent odometers as well. Indeed, we will give this as Lemma 9.9, below.
9.4. Using boundary strings for left-lattice adjacency. Unlike the argument for tiles, we cannot apply $90^{\circ}$ symmetry and $v_{0 i}=\mathbf{i} v_{i 0}$ to add the case of left-lattice adjacent tile odometers to Lemma 9.6 , since odometers are only $180^{\circ}$ symmetric in general. Instead, we will express the shared boundary of the touching tiles in terms of boundary strings and inductively use the compatibility of the restrictions of the odometers to the tiles making up the boundary strings.

To do this, we first need to strengthen our notion of boundary string: We say a partial odometer respects a string when its domain includes all tiles of the string, and its restrictions to those tiles are tile odometers that are consecutively left-lattice adjacent.

Lemma 9.7. Suppose $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$ is a proper Descartes quadruple, $h_{0}$ is a tile odometer for the circle $C_{0}$ with domain $T_{0}$, and write $v_{i j}=v\left(C_{i}, C_{j}\right)$. For each $i=1,2,3$ for which $c_{i}>0$, we have that if $R_{i}^{ \pm}$are the endpoints of the $C_{i}$ boundary string $\mathcal{R}$ for $T_{0}$ and $h_{R_{i}^{ \pm}}$is a tile odometer for $C_{i}$ with domain $R_{i}^{ \pm}$, which satisfies

$$
\mathbf{s}\left(h_{R_{i}^{ \pm}}\right)-\mathbf{s}\left(h_{0}\right)=\frac{1}{2}\left(v_{i 0} \pm v_{0 i}\right)
$$

then $h_{R_{i}^{+}}$and $h_{R_{i}^{-}}$are each subtile-lattice adjacent to $h_{0}$, and $h_{0} \cup h_{R_{i}^{+}} \cup h_{R_{i}^{-}}$ respects the $C_{i}$ boundary string of $T_{0}$.

Proof. Recall from the definition of boundary strings that

$$
\mathrm{c}\left(R_{i}^{ \pm}\right)-\mathrm{c}\left(T_{0}\right)=\frac{1}{2}\left(v_{i 0} \pm v_{0 i}\right)
$$

Therefore each $h_{R_{i}^{ \pm}}$is subtile-lattice adjacent to $h_{0}$ as claimed. Moreover, Lemma 8.1 gives


Figure 9.1. Verifying a boundary string.

$$
\begin{equation*}
\mathrm{c}\left(R_{i}^{+}\right)-\mathrm{c}\left(T_{i}^{-}\right)=v_{k i} \quad \text { and } \quad \mathrm{c}\left(T_{i}^{-}\right)-\mathrm{c}\left(R_{i}^{-}\right)=-v_{j i} . \tag{9.4}
\end{equation*}
$$

In particular, for $i=1$, then our assumption that $c_{1}>c_{2}, c_{3}$ implies that $R_{1}^{-}, T_{i}^{-}, R_{1}^{+}$are a triple of consecutively touching tiles in a tiling of $T_{1}^{-}$under the lattice generated by $\left\{v_{41}, v_{21}, v_{31}\right\}$; this is already sufficient to imply that they form a string (either from $R_{1}^{-}$to $R_{1}^{+}$or vice versa), and the sign in (5.1e) implies that they are the string from $R_{1}^{-}$to $R_{1}^{+}$. Moreover, since

$$
\mathbf{s}\left(h_{1}^{-}\right)-\mathbf{s}\left(h_{0}\right)=-\frac{1}{2}\left(a_{32}+\mathbf{i} a_{32}\right),
$$

we can calculate using the lattice rules (5.1) that

$$
\begin{equation*}
\mathbf{s}\left(h_{R_{1}^{+}}\right)-\mathbf{s}\left(h_{i}^{-}\right)=a_{31}, \quad \mathbf{s}\left(h_{1}^{-}\right)-\mathbf{s}\left(h_{R_{1}^{-}}\right)=-a_{21}, \tag{9.5}
\end{equation*}
$$

and (9.4) and (9.5) together give that $h_{R_{i}^{-}}, h_{T_{1}^{-}}, h_{R_{i}^{+}}$are consecutively leftlattice adjacent; thus $h_{0} \cup h_{R_{i}^{+}} \cup h_{R_{i}^{-}}$respects the string.

The cases $i=2,3$ require the induction hypothesis and, since they are similar, we handle only the case $i=3$. Decomposing $h_{0}$ according to tile odometers $h_{i}^{ \pm}$on tiles $T_{i}^{ \pm}$according to Definition 9.3 and using the lattice rules (5.1), we check that

$$
\begin{array}{r}
\mathrm{c}\left(R_{3}^{+}\right)-\mathrm{c}\left(T_{2}^{+}\right)=\frac{1}{2}\left(v_{32}-v_{23}\right) \quad \text { and } \quad \mathrm{c}\left(T_{3}^{-}\right)-\mathrm{c}\left(T_{2}^{+}\right)=\frac{1}{2}\left(v_{32}+v_{23}\right), \\
\mathrm{s}\left(h_{R_{3}^{+}}\right)-\mathrm{s}\left(h_{2}^{+}\right)=\frac{1}{2}\left(a_{32}-a_{23}\right) \quad \text { and } \quad \mathrm{s}\left(h_{3}^{-}\right)-\mathrm{s}\left(h_{2}^{+}\right)=\frac{1}{2}\left(a_{32}+a_{23}\right) ;
\end{array}
$$

see Figure 9.1. In particular, the string $\mathcal{S}^{\prime}$ from $T_{3}^{-}$to $R_{3}^{+}$is the $C_{3}$ boundary string of $T_{2}^{+}$, and by induction, $h_{0} \cup h_{R_{i}^{+}} \cup h_{R_{i}^{-}}$respects this string. In exactly the same way, we can check that the string $\mathcal{S}$ from $R_{3}^{-}$to $T_{3}^{-}$is the $C_{3}$ boundary string of $T_{1}^{+}$and is respected by $h_{0} \cup h_{R_{i}^{+}} \cup h_{R_{i}^{-}}$. By Lemma 8.1, the $C_{3}$-string for $T_{0}$ is the concatenation of $\mathcal{S}$ and $\mathcal{S}^{i}$, and this string is respected by $h_{0} \cup$ $h_{R_{i}^{+}} \cup h_{R_{i}^{-}}$since all pairs of consecutive tiles in the concatenation are already consecutive tiles in either $\mathcal{S}$ or $\mathcal{S}^{\prime}$.

We now prove the compatibility of left-lattice adjacent odometers:


Figure 9.2. Verifying tiling $T_{C_{0}}$ by $\Lambda_{C_{0}}$ using boundary strings.

Lemma 9.8. If the odometer $h_{0}^{\prime}$ for $C_{0}$ is left-lattice adjacent to $h_{0}$, then $h_{0}^{\prime}$ and $h_{0}$ are compatible. Moreover, if $c_{i}>1$, where $\mathrm{c}\left(T\left(h_{0}^{\prime}\right)\right)-\mathrm{c}\left(T\left(h_{0}\right)\right)= \pm v_{i 0}$, then any vertex in $T\left(h_{0}\right) \cap T\left(h_{0}^{\prime}\right)$ lies, together with all of its lattice neighbors, in the union of the domains of some pair of proper ancestor odometers of $h_{0}$ and $h_{0}^{\prime}$ that are pairwise left-lattice adjacent.

Proof. If $c_{i}=0$, then $C_{0}$ is a Ford circle and this lemma is easy to verify from the construction in Section 4.2, so we may assume $c_{i}>0$.

From the definition of left-lattice adjacency, we have without loss of generality that $h_{0}^{\prime}(x)=h_{0}^{i}(x)=h_{0}\left(x-v_{i 0}\right)+a_{i 0} \cdot x$ of $h_{0}$ with domain $T_{0}^{i}=T_{0}+v_{i 0}$, and we let $R_{i}$ and $S_{i}$ be tiles for $C_{i}$ satisfying

$$
\begin{aligned}
& \mathrm{c}\left(R_{i}\right)-\mathrm{c}\left(T_{0}\right)=\frac{1}{2}\left(v_{i 0}-v_{0 i}\right), \\
& \mathrm{c}\left(S_{i}\right)-\mathrm{c}\left(T_{0}\right)=\frac{1}{2}\left(v_{i 0}+v_{0 i}\right) .
\end{aligned}
$$

From Lemma 7.3, we know that $R_{i}^{-}, T_{0}, T_{0}^{i}$ and $R_{i}^{+}, T_{0}, T_{0}^{i}$ are both touching triples; in particular, the intersection of $T_{0}$ and $T_{0}^{i}$ is a path from $S_{i}$ to $R_{i}$. Therefore, let $\mathcal{R}=R_{i}^{0}, R_{i}^{1}, \ldots, R_{i}^{t}$ and $\mathcal{S}=S_{i}^{0}, S_{i}^{1}, \ldots, S_{i}^{t}$ be the string from $R_{i}^{-}=R_{i}^{0}$ to $R_{i}^{+}=R_{i}^{t}$ and the string from $R_{i}^{+}=S_{i}^{0}$ to $R_{i}^{-}=S_{i}^{t}$, respectively; note that they are equivalent under a central reflection (Figure 9.2). Lemma 9.7 guarantees that the interiors of $\mathcal{R}$ and $\mathcal{S}$ lie in $T_{0}$ and $T_{0}+v_{i 0}$, respectively. By Lemma 8.2, the intersection of the interiors of $\mathcal{R}$ and $\mathcal{S}$ contains a simple path from $S_{i}$ to $R_{i}$, so $T_{0} \cap T_{0}^{i}$ equals the intersection of the interiors of $\mathcal{R}$ and $\mathcal{S}$. So it suffices to show for each pair of touching $R_{i}^{\ell_{1}}, S_{i}^{\ell_{2}}$ from $\mathcal{R}$ and $\mathcal{S}$ that the restrictions $h_{0} \mid R_{i}^{\ell_{1}}$ and $h_{0} \mid S_{i}^{\ell_{2}}$ are left-lattice adjacent (thus compatible by induction). Note that $c_{i}>1$ implies that there is no vertex in $\mathbb{Z}^{2}$ that lies in
the intersection of four tiles in a $T_{i}+\Lambda_{C_{i}}$ tiling of $\mathbb{Z}^{2}$, justifying the Moreover clause.

We let $h_{R_{i}^{ \pm}}$be defined as in Lemma 9.7 and let $f_{0}^{i}=h_{0} \cup h_{R_{i}^{-}} \cup h_{R_{i}^{+}}$ and $f_{1}^{i}=h_{0}^{i} \cup h_{R_{i}^{-}} \cup h_{R_{i}^{+}}$which, by Lemma 9.7 and inductive application of Lemma 9.9 to the sublattice adjaceny of $h_{R_{i}^{ \pm}}$to $h_{0}$, are well-defined partial odometers that respect the strings $\mathcal{R}$ and $\mathcal{S}$, respectively. Considering now touching tiles $R_{i}^{\ell_{1}}$ and $S_{i}^{\ell_{2}}$ from $\mathcal{R}$ and $\mathcal{S}$, respectively, we see that

$$
R_{i}^{\ell_{1}}, R_{i}^{\ell_{1}-1}, \ldots, R_{i}^{0}=S_{i}^{t}, S_{i}^{t-1}, \ldots, S_{i}^{\ell_{2}}
$$

is a sequence of $\left(\ell_{1}+\ell_{2}+1-t\right)$ tiles where each consecutive pair $U^{m}, U^{m+1}$ in the sequence forms a two-tile string that is respected either by $f_{0}^{i}$ or $f_{1}^{i}$. In particular, let $C_{i}^{1}, C_{i}^{2}, C_{i}^{3}$ be the parents of $C_{i}$ in clockwise order. For each pair $U^{m}, U^{m+1}$ that is respected by $f_{q_{m}}^{i}\left(q_{m} \in\{0,1\}\right)$, we have for some $n$ (which may vary with $m$ ) that

$$
\begin{equation*}
\mathrm{s}\left(f_{q_{m}}^{i} \mid U^{m+1}\right)-\mathrm{s}\left(f_{q_{m}}^{i} \mid U^{m}\right)= \pm a\left(C_{i}, C_{i}^{n}\right) \tag{9.6}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{c}\left(U^{m+1}\right)-\mathrm{c}\left(U^{m}\right)= \pm v\left(C_{i}, C_{i}^{n}\right) \tag{9.7}
\end{equation*}
$$

We now have that

$$
\pm v\left(C_{i}, C_{i}^{n}\right)=\mathrm{c}\left(S_{i}^{\ell_{2}}\right)-\mathrm{c}\left(R_{i}^{\ell_{1}}\right)=\sum_{m}\left(\mathrm{c}\left(U^{m+1}\right)-\mathrm{c}\left(U^{m}\right)\right)
$$

implies that

$$
\mathrm{s}\left(f_{1}^{i} \mid S_{i}^{\ell_{2}}\right)-\mathrm{s}\left(f_{0}^{i} \mid R_{i}^{\ell_{1}}\right)=\sum_{m}\left(\mathrm{~s}\left(f_{q_{m}}^{i} \mid U^{m+1}\right)-\mathrm{s}\left(f_{q_{m}}^{i} \mid U^{m}\right)\right)= \pm a\left(C_{i}, C_{i}^{n}\right)
$$

In particular, $f_{1}^{i} \mid S_{i}^{\ell_{2}}$ and $f_{0}^{i} \mid R_{i}^{\ell_{1}}$ are left-lattice adjacent.
As noted earlier, Lemma 9.6 now gives us the following:
Lemma 9.9. If $h$ is a tile odometer which is right-lattice or subtile-lattice adjacent to the tile odometer $h_{0}$, then $h$ and $h_{0}$ are compatible.
9.5. Existence of tile odometers. We conclude with the following.

Lemma 9.10. Every circle $C_{0} \in \mathcal{B}$ has a tile odometer $h_{0}: T_{0} \rightarrow \mathbb{Z}$.
Proof. In light of Section 4, we may assume that $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$ is a proper Descartes quadruple with $c_{1} \geq c_{2}>1$. Copying the proof of Lemma 7.2, we easily obtain tile odometers $h_{i}^{ \pm}$for each circle $C_{i}$ such that

$$
\mathrm{c}\left(T\left(h_{i}^{ \pm}\right)\right)-\mathrm{c}\left(T_{0}\right)= \pm \frac{1}{2}\left(v_{k j}-\mathbf{i} v_{k j}\right)
$$

and

$$
c_{i}=0 \quad \text { or } \quad \mathbf{s}\left(h_{i}^{ \pm}\right)-\mathbf{s}\left(h_{0}\right)= \pm \frac{1}{2}\left(a_{k j}+\mathbf{i} a_{k j}\right)
$$

for all rotations $(i, j, k)$ of $(1,2,3)$. The difficulty lies in checking that there are height offsets so that the $h_{i}^{ \pm}$have a common extension to $T_{0}$.

We know from part (T6) of Lemma 7.3 that the only pairs of subodometers $h_{i}^{ \pm}$whose domains intersect are the pairs

$$
\begin{array}{ll}
\left(h_{i}^{ \pm}, h_{j}^{\mp}\right) & \text { for } i \neq j, \\
\left(h_{1}^{ \pm}, h_{2}^{ \pm}\right), & \text {and } \\
\left(h_{1}^{+}, h_{1}^{-}\right) . &
\end{array}
$$

Thus, as in the proof of Lemma 9.2, we only need to check compatibility of these pairs. The tile odometers $h_{i}^{ \pm}, h_{j}^{\mp}$ for $i \neq j$ are easily verified to be subtilelattice adjacent, and so they are compatible by Lemma 9.9. It remains to check compatibility for tile odometer pairs corresponding to the tile pairs ( $T_{1}^{-}, T_{2}^{-}$), $\left(T_{1}^{+}, T_{2}^{+}\right)$, and ( $T_{1}^{+}, T_{1}^{-}$). To accomplish this, we decompose $T_{1}^{+}$and $T_{1}^{-}$into unions of $Q_{i}^{ \pm}$and $S_{i}^{ \pm}$for $i=2,3,4$ according to (7.1), as shown in Figure 7.3. The second part of (T6) implies that it suffices to check compatibility of the restrictions of the odometers $h_{i}^{ \pm}$to the pairs $\left(S_{2}^{-}, T_{2}^{-}\right),\left(Q_{2}^{+}, T_{2}^{+}\right),\left(Q_{2}^{+}, S_{3}^{-}\right)$, and $\left(Q_{3}^{+}, S_{2}^{-}\right)$. Using the lattice rules we verify that these pairs are rightlattice, right-lattice, subtile-lattice, and subtile-lattice adjacent, respectively, so we are done by induction.

## 10. Global odometers

Having constructed tile odometers $h_{C}$ for each circle $C \in \mathcal{B}$ by Lemma 9.10, we now check that the $h_{C}$ extend to global odometers $g_{C}$ satisfying $\Delta g_{C} \leq 1$ and prove our main theorem. We first observe that the partial odometers glue together to form global odometers with the correct periodicity and growth at infinity. (As in the previous section, a tile now is just a subset of $\mathbb{Z}[\mathbf{i}]$.)

Lemma 10.1. For every circle $C \in \mathcal{B}$, there is a function $g_{C}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ that has a restriction to a tile odometer for $C$ for which the periodicity condition (1.5) holds for $v \in \Lambda_{C}$ and for which

$$
\begin{equation*}
x \mapsto g_{C}(x)-\frac{1}{2} x^{t} A_{C} x-b \cdot x \tag{10.1}
\end{equation*}
$$

is $\Lambda_{C}$-periodic for some $b \in \mathbb{R}^{2}$.
Proof. We may suppose $C=C_{0}$ and $\left(C_{0}, C_{1}, C_{2}, C_{3}\right)$ is a proper Descartes quadruple, $h_{i}$ is a tile odometer for $C_{i}$ whose domain is the tile $T_{i}$, and write $v_{i j}=v\left(C_{i}, C_{j}\right), a_{i j}=a\left(C_{i}, C_{j}\right)$. We write $h=h_{0}$ and let $h(k, \ell)$ be the tile odometer on the domain $T(h)+k v_{10}+\ell v_{20}$ given by

$$
x \mapsto h\left(x-k v_{10}-\ell v_{20}\right)+\left(k a_{10}+\ell a_{20}\right) \cdot x .
$$

Now any pair of overlapping tile odometers in

$$
\mathcal{H}=\{h(k, \ell): k, \ell \in \mathbb{Z}\}
$$

are left-lattice adjacent and thus compatible. By Lemma 9.2, there is a function $g: \mathbb{Z}^{2} \rightarrow \mathbb{Z}$ with $g(0)=0$ that is compatible with every $h \in \mathcal{H}$. Since $\mathcal{H}$ is invariant under

$$
h \mapsto\left(x \mapsto h\left(x-v_{i 0}\right)+a_{i 0} \cdot x\right)
$$

for each $i \in\{1,2,3\}$, we see from Lemma 9.2 that

$$
x \mapsto g\left(x-v_{i 0}\right)+a_{i 0} \cdot x
$$

differs from $g$ by some constant:

$$
g\left(x+v_{i 0}\right)=\beta_{i}+a_{i 0} \cdot x+g(x)
$$

But $g(0)=0$ implies that $g\left(v_{i 0}\right)=\beta_{i}$, so that

$$
g\left(x+v_{i 0}\right)=g\left(v_{i 0}\right)+a_{i 0} \cdot x+g(x)
$$

for all $i=1,2,3$ and $x \in \mathbb{Z}^{2}$. Together with $a_{i 0}=A_{C_{0}} v_{i 0}$ from (5.2), this implies the periodicity condition (1.5) for $v \in \Lambda_{C}$ and that (10.1) is $\Lambda_{C}$-periodic for some $b \in \mathbb{R}^{2}$.

To prove this criterion for general odometers $g_{C}$, we follow the outline of Proposition 4.3. To begin, we need to understand the Laplacian $\Delta g_{C}$ for all $C \in \mathcal{B}$. Let $N(x)$ and $\bar{N}(x)$ denote the set $\{x \pm 1, x \pm \mathbf{i}\}$ of lattice neighbors of $x \in \mathbb{Z}[\mathbf{i}]$ and $\{x\} \cup N(x)$, respectively.

LEMMA 10.2. Let $h^{1}$ and $h^{2}$ be compatible tile odometers for (tangent or identical) circles in $\mathcal{B}$, let $h=h^{1} \cup h^{2}$, and let $x$ such that $x \in T\left(h^{1}\right) \cap T\left(h^{2}\right)$ and $\bar{N}(x) \subseteq T\left(h^{1}\right) \cup T\left(h^{2}\right)$. If $h^{1}, h^{2}$ are left-lattice adjacent, then $\Delta h(x)=1$. If $h^{1}, h^{2}$ are right-lattice adjacent or subtile-lattice adjacent, then $\Delta h(x)=0$ if $x \in s_{y} \nsubseteq T\left(h^{1}\right) \cup T\left(h^{2}\right)$ for some $y$, and $\Delta h(x)=1$ otherwise.

Proof. The proof is by induction on the areas of $T\left(h^{1}\right)$ and $T\left(h^{2}\right)$. For the cases of right-lattice and subtile-lattice adjacency, the base case occurs when $h^{1}$ and $h^{2}$ are both tile odometers for a circle of curvature 1 , in which case the statement can be checked by hand; it is sufficient to check for the circle $(1,1)$. Thus for the inductive step in this case, we assume (without loss of generality) $h^{1}$ is a tile odometer for a circle of curvature $>1$ : in particular, it can be decomposed into subodometers according to Definition 9.3.

For this case we let $x$ be a point such that $\bar{N}(x) \subseteq T\left(h^{1}\right) \cup T\left(h^{2}\right)$. We claim that there is a subodometer $h^{\prime}$ of $h^{1}$ such that $\bar{N}(x) \subseteq T\left(h^{\prime}\right) \cup T\left(h^{2}\right)$. Indeed, the definition of a tile ensures that if $h^{i}(i=1,2)$ covers $s_{x}$ and $s_{x-1-\mathbf{i}}$, then it must also cover either $s_{x-1}$ or $s_{x-\mathbf{i}}$. (In particular, at least three of the four squares containing $x$ as a vertex are covered by $h$.) Thus, without loss of generality, we have that $s_{x}$ is covered by $h^{2}$ and $s_{x-1-\mathbf{i}}$ is covered by $h^{1}$. In this case we let $h^{\prime}$ be the subodometer of $h^{1}$ whose domain covers $s_{x-1-\mathbf{i}}$, and we have that $\bar{N}(x) \subseteq T\left(h^{\prime}\right) \cup T\left(h^{2}\right)$. Lemma 9.6 now implies that $h^{\prime}$ and $h^{2}$ are
subtile-lattice adjacent or left-lattice adjacent, so we are done by induction. (In particular, note that the condition on $s_{y}$ that determines whether $\Delta h(x)$ is 0 or 1 is unchanged with the inductive step.)

For the case of left-lattice adjacency, Lemma 9.8 gives the statement by induction. The base case is when $h^{1}$ and $h^{2}$ are left-lattice adjacent along a $v_{i 0}$ for which the curvature $c_{i}$ of the corresponding parent circle $C_{i}$ is 1 . In this case, however, the proof of Lemma 9.8 gives that $x$ lies, together with all of its neighbors, in the union of four tile odometers for $C_{i}=(1,1+2 z)(z \in \mathbb{Z}[\mathbf{i}])$ that are cyclically left-lattice adjacent, and this case can be checked by hand.

Note that from the case of left-lattice adjacency, Lemma 10.2 has the crucial consequence that $\Delta g_{C} \equiv 1$ on the "web" of its tile boundaries, as seen in Figure 1.2. Next, we analyze $\Delta g_{C}$ on the interior of $T_{C}$.

Lemma 10.3. Suppose $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}^{4}$ is a proper Descartes quadruple, write $C_{i}=\left(c_{i}, z_{i}\right)$ and $T_{i}$ is a tile for $C_{i}$, and decompose $T_{0}$ as a union of $T_{i}^{ \pm}$as in (7.1). Let $x \in T_{0} \backslash \partial T_{0}$, and let $k$ denote the number of boundaries $\partial T_{i}^{ \pm}$of subtiles that contain $x$. Then we have
(1) if $k=2$ and $x \neq \mathrm{c}\left(T_{0}\right)$, then $\Delta g_{C_{0}}(x)=1$;
(2) if $k=2$ and $x=\mathrm{c}\left(T_{0}\right)$, then $\Delta g_{C_{0}}(x)=0$;
(3) if $k=3$, then $\Delta g_{C_{0}}(x)=0$;
(4) if $k=4$ and $x \neq \mathrm{c}\left(T_{0}\right)$, then $\Delta g_{C_{0}}(x)=-1$;
(5) if $k=4$ and $x=\mathrm{c}\left(T_{0}\right)$, then $\Delta g_{C_{0}}(x)=-2$.

In particular, as $180^{\circ}$ symmetry precludes $k=3$ when $x=\mathrm{c}\left(T_{0}\right)$, we have

$$
\Delta g_{C_{0}}(x)=3-k-\mathbf{1}_{\{\mathrm{c}(T)\}}(x)
$$

whenever $k(x) \geq 2$. Note that cases (4) and (5) arise only when $C_{0}$ is a Ford or Diamond circle, so these cases have been proved already in Section 4. These rules can be witnessed in "general" tile decomposition in Figure 1.4. (When looking at this figure, keep in mind that the curve surrounding the Soddy twin tile is not the meeting of two borders.)

Remark 10.4. The above lemma provides a fast algorithm for recursively generating tiles with their associated Laplacian patterns; the base case for the recursion is given by the base cases of Definition 9.3. The appendix lists all such patterns associated to circles in $\mathcal{B}$ of curvature $1 \leq c \leq 100$.

One may be tempted to define the odometers using this lemma in place of our complicated recursive construction. Indeed, by Liouville theorem, the Laplacian of the odometer determines the odometer up to a harmonic polynomial. Every periodic Laplacian pattern can be integrated to a function in this way. However, there is no guarantee in general that the function so constructed is integer valued. Thus in some sense, an important point of the construction in
this manuscript is that these particular patterns do integrate to integer-valued functions.

Proof of Lemma 10.3. For the Ford and Diamond circles, this lemma has already been verified in Section 4.2. Thus, we may assume by induction that the $c_{i}$ are distinct and positive and that the lemma holds for the proper Descartes quadruples $\left(C_{1}, C_{4}, C_{2}, C_{3}\right)$ and $\left(C_{2}, C_{3}, C_{5}, C_{6}\right)$ for $C_{4}, C_{5}, C_{6} \in \mathcal{B}$. Since $c_{1}>1$, we can decompose $T_{1}^{+}$and $T_{1}^{-}$into $Q_{i}^{ \pm}$and $S_{i}^{ \pm}$as in Claim 7.6 and Figure 7.3.

First consider the case where $k=3$. Definition 9.3 (and part (T6) from Lemma 7.3) give that the two subtiles containing $x$ are subtile-lattice adjacent. Since $k=3$, there exists the square $s_{y}$ in the final hypothesis of Lemma 10.2, and thus Lemma 10.2 gives that $\Delta g_{C_{0}}(x)=0$.

Next consider the case where $x \neq \mathrm{c}\left(T_{0}\right)$ and $k=2$. If the two subtiles whose boundaries contain $x$ are not the pair $\left\{T_{1}^{+}, T_{1}^{-}\right\}$, then by Definition 9.3 (and part (T6) from Lemma 7.3), the two subtiles containing $x$ are subtilelattice adjacent, and Lemma 10.2 now gives the result.

If on the other hand $k=2$ and $x \in \partial T_{1}^{+} \cap \partial T_{1}^{-}$, then we use the double decomposition of $T_{1}^{+}, T_{1}^{-}$. As in the proof of Lemma 10.2, we are guaranteed that two tiles from among $S_{3}^{-}, Q_{3}^{+}, S_{2}^{-}, Q_{2}^{+}, S_{4}^{+}$cover the neighborhood $\bar{N}(x)$. As in the proof of Lemma 9.10, all pairs among these tiles are known to have induced tile odometers that are subtile-lattice adjacent or right-lattice adjacent except for the pairs $\left\{Q_{2}^{+}, S_{2}^{-}\right\}$and $\left\{S_{3}^{-}, Q_{3}^{+}\right\}$. But (T6) implies that this case cannot occur unless $c_{i}=0$ for some $i \in\{3,4\}$, in which case $C_{0}$ is a Ford circle.

Finally, if $x=\mathrm{c}(T)$, then $x \in \partial T_{1}^{+} \cap \partial T_{1}^{-}$implies that $c_{4}=0$. In particular, $C_{1}=C_{p q}$ is a Ford circle with Ford circle parents $C_{2}=C_{p_{1} q_{1}}$ and $C_{3}=C_{p_{2} q_{2}}$. We have that $\bar{N}(x)$ is covered by $T_{1}^{+} \cup T_{1}^{-}$. In this case, the tile odometers $g_{p q}^{ \pm}$on $T_{1}^{ \pm}$are related by $q_{p q}^{+}\left(x-\left(q_{2}-q_{1}, q\right)\right)=q_{p q}^{-}(x)+\left(p, p_{2}-p_{1}\right) \cdot x+k$ for some constant $k \in \mathbb{Z}$. Thus, the explicit formula from Section 4.2 can be used to verify this case.

We now generalize the inductive argument in Proposition 4.3 to obtain maximality of general odometers.

Lemma 10.5. For each $C \in \mathcal{B}, g_{C}$ is maximal.
Proof. Suppose $X \subseteq \mathbb{Z}^{2}$ is connected and infinite and $\Delta\left(g_{C}+\mathbf{1}_{X}\right) \leq 1$. Let $Y$ be a connected component of $\mathbb{Z}^{2} \backslash X$, and observe that $\Delta\left(g_{C}-\mathbf{1}_{Y}\right) \leq 1$. In particular, we may assume that $Y=\mathbb{Z}^{2} \backslash X$. It is enough to show that $Y$ must be empty. Now, if $Y$ is not contained in $T \backslash \partial T$ for some $T \in T_{C}+\Lambda_{C}$, then by Lemma 10.2 , there is a point $x \in Y$ such that $\Delta g_{C}(x)=1$ and $\Delta \mathbf{1}_{X}(x)>0$, since the tile odometers of which $g_{C}$ consist are left-lattice adjacent. Thus, we may assume $Y \subseteq T_{C} \backslash \partial T_{C}$. The lemma is now immediate from the following claim.

Claim. Suppose $Y \subseteq T_{C}$ is simply connected, $Y \backslash \partial T_{C}$ is nonempty, $Y \cap \partial T$ is connected, and $Y \cap \partial T_{C} \cap T_{i}^{ \pm}$is nonempty for at most one subtile $T_{i}^{ \pm}$. Then there is a vertex $x \in Y \backslash \partial T_{C}$ such that $\Delta\left(g_{C}-\mathbf{1}_{Y}\right)(x)>1$.

We prove this by induction on the curvature of $C$. Note that, by the proofs of Propositions 4.3 and 4.7, we may assume that $C$ is neither a Ford nor a Diamond circle. In particular, each $T_{i}^{ \pm}$contains at least one square and the pairwise intersections are exactly what we expect from the picture. We may assume that $C_{1}$ is the largest parent. Thus $T_{i}^{ \pm}$and $T_{j}^{\mp}$ have simply connected intersection when $i \neq j$ and $T_{i}^{ \pm}$and $T_{i}^{\mp}$ are disjoint except when $i=1$, in which case the intersection is simply connected.

Case 1. $Y$ is contained in the interior of some $T_{i}^{ \pm}$. The claim follows either by the induction hypothesis or by the corresponding results for Ford and Diamond circles in Propositions 4.3 and 4.7.

Case 2. Some $x \in \partial Y \backslash \partial T$ lies in the boundary of exactly $T_{i}^{+}$and $T_{j}^{-}$ with $i \neq j$. Observe that $\Delta \mathbf{1}_{Y}(x)<0$ and $x \neq \mathrm{c}(T)$. Thus Lemma 10.3 gives the claim.

Case 3. Some $x \in \partial Y \backslash \partial T$ lies in the boundary of exactly three $T_{i}^{ \pm}$. Since Case 2 is excluded, we must have $\Delta \mathbf{1}_{Y}(x)<-1$, and thus Lemma 10.3 again gives the claim.

Case 4. In the exclusion of the above three cases, the topology of the tile decomposition implies that $Y$ lies in the union of the interior of $T_{1}^{+}$, the interior of $T_{1}^{-}$, and the intersection of one $T_{1}^{ \pm} \cap \partial T$. If $Y \cap T_{1}^{+}$is nonempty, then we can inductively apply the claim to $Y \cap T_{1}^{+}$. Otherwise, we can inductively apply the claim to $Y \subseteq T_{1}^{-}$.

From Lemmas 10.1 and 10.5, we immediately obtain Theorem 1.2, modulo checking that $\Lambda_{C_{0}}$ and $L_{C_{0}}$ are in fact the same lattice.

Theorem 10.6. $L_{C}=\Lambda_{C}$ for all $C \in \mathcal{B}$.
Proof. Lemma 5.2 verified $\Lambda_{C} \subseteq L_{C}$, thus it remains to verify $L_{C} \subseteq \Lambda_{C}$. For $C \in \mathcal{B}$, let $g_{C}$ be the odometer for $C$, as constructed in Section 10. Recall from (1.5) that $g_{C}$ satisfies

$$
g_{C}(x+v)=g_{C}(x)+x^{t} A_{C} v+g_{C}(v)
$$

for $v \in \Lambda_{C}$. We will now modify $g_{C}$ to produce an odometer for $C$ that satisfies the periodicity condition (1.5) for the lattice $L_{C}$. In particular, for $v \in L_{C}$, we let

$$
g_{C}^{v}(x)=g_{C}(x+v)-x^{t} A_{C} v-g_{C}(v),
$$

which is an integer since $A_{C} v \in \mathbb{Z}^{2}$ by the definition of $L_{C}$. Note that for $v \in \Lambda_{C}$, we have $g^{v}(x)=g(x)$.

We now define

$$
g^{\prime}(x):=\min _{v \in L_{C}} g^{v}(x) .
$$

Note that the periodicity $g^{v}(x)=g(x)$ for $v \in \Lambda_{C}$ implies that this can be interpreted as a finite minimum over the quotient $L_{C} / \Lambda_{C}$.

In particular, up to an additive constant, $g^{\prime}(x)$ is still an odometer for $A_{C}$ and now satisfies the periodicity condition (1.5) for the lattice $L_{C}$. In particular, we have that the average Laplacian $\bar{\Delta} g^{\prime}$ of $g^{\prime}$ over one period of $L_{C}$ must satisfy $\frac{1}{d}=\operatorname{trace}\left(A_{C}\right) \leq \bar{\Delta} g^{\prime} \leq \bar{\Delta} g=\frac{1}{d}$. But then we must have $\operatorname{det}\left(L_{C}\right) \geq \frac{1}{d}$, and thus $L_{C}=\Lambda_{C}$.

## Appendix A. Table of odometer patterns

Here we display proper Descartes quadruples $\left(C_{0}, C_{1}, C_{2}, C_{3}\right) \in \mathcal{B}$, along with the Soddy precursor $C_{4}=2\left(C_{1}+C_{2}+C_{3}\right)-C_{0}$ of $C_{0}$, the vectors $v\left(C_{i}, C_{0}\right)$ and $a\left(C_{i}, C_{0}\right)(i=1,2,3)$, the tile odometer for $C_{0}$, and a tiling neighborhood in $T_{C}+L_{C}$. We display quadruples up to symmetry for $1 \leq c_{0} \leq 156$. An extended appendix with more circles is available as an ancillary file for this manuscript at arXiv.org.

| $C_{0}$ |  |  |
| ---: | ---: | ---: | ---: |
| $C_{1}$ | $v\left(C_{1}, C_{0}\right)$ | $a\left(C_{1}, C_{0}\right)$ |
| $C_{2}$ | $v\left(C_{2}, C_{0}\right)$ | $a\left(C_{2}, C_{0}\right)$ |
| $C_{3}$ | $v\left(C_{3}, C_{0}\right)$ | $a\left(C_{3}, C_{0}\right)$ |
| $C_{4}$ |  |  |$\quad$ tile odometer $\quad$| One |
| :--- |
| of $T_{C}$ in $T_{C}+L_{C}$. |









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