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# Algebraic aspects of Abelian sandpile models 

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#### Abstract

The Abelian sandpile models feature a finite Abelian group $G$ generated by the operators corresponding to particle addition at various sites. We stady the canonical decomposition of $G$ as a product of cyclic groups $G=Z_{d_{1}} \times Z_{d_{2}} \times Z_{d_{9}} \cdots \times Z_{d_{g}}$, where $g$ is the least number of generators of $G$, and $d_{i}$ is a multiple of $d_{i+1}$. The structure of $G$ is determined in terms of the toppling matrix $\Delta$. We construct scalar functions, linear in the height variables of the pile, that are invariant under toppling at any site. These invariants provide convenient coordinates to label the recurrent configurations of the sandpile. For an $L \times L$ square lattice, we show that $g=L$. In this case, we observe that the system has non-trivial symmetries, transcending the obvious symmetries of the square, namely those coming from the action of the cyclotomic Galois group $\mathrm{Gal}_{L}$ of the $2(L+1)$ th roots of unity (which operates on the set of eigenvalues of $\Delta$ ). These eigenvalues are algebraic integers, the product of which is the order $|G|$. With the help of Gal ${ }_{L}$ we are able to group the eigenvalues into certain subsets the products of which are separately integers, and thus obtain an explicit factorization of $|G|$. We also use $\mathrm{Gal}_{L}$ to define other simpler sets of toppling invariants.


## 1. Introduction

The concept of self-organized criticality was proposed by Bak, Tang and Wiesenfeld in 1987, and has given rise to growing interest in the study of self-organizing systems. Bak et al argued that, in many natural phenomena, the dissipative dynamics of the system is such that it drives the system to a critical state, thereby leading to ubiquitous power-law behaviours [1,2]. This mechanism has been invoked to understand the power laws observed in turbulent fluids, earthquakes, distribution of visible matter in the universe, solar flares and surface roughening of growing interfaces [3-5].

Sandpile automata are among the simplest theoretical models which show self-organized criticality. A specially nice subclass consists of the so-called Abelian sandpile models (ASMS) [6]. There have been many numerical as well as analytical studies of the ASM. The case when there is a preferred direction of particle transfer turns out to be equivalent to the voter model, and all the critical exponents can be determined in all dimensions [7]. When there is no preferred direction, the model turns out to be related to the $q \rightarrow 0$ limiting case of the Potts model [8]. The problem has been solved exactly in the mean field limit [9-11]. In two and three dimensions, only some of the critical exponents of the problem are known analytically $[6,8,12]$.

Most of these papers have been concerned with the critical properties of the ASM in the thermodynamic limit of large system sizes. In this paper, our interest is rather in the

[^0]properties of the finite automata. There have been only a handful of papers addressing these so far (and mainly in two dimensions). Creutz' paper [13] exhibits very interesting examples of geometrical patterns displayed by a special configuration (the so-called identity configuration) of the ASM on a square lattice. Liu et al [14] have studied the patterns obtained by relaxing some simple unstable configurations. Wiesenfeld et al [15] have studied the periods of deterministic ASM, again on a square lattice. These studies have been extended by Markosova and Markos [16] to lattices of size up to 19. This paper was inspired partly by the 'experimental' results of these authors. We will concentrate on such properties as the structures of the Abelian group and of the space of recurrent configurations.

The plan of this paper is as follows. In section 2, we give the definition of general ASM, and review their basic properties. In section 3, we construct, for a general ASM, a set of functions, defined on the space of recurrent configurations, which are invariant under topplings, and which can be used to label these configurations. With these functions, we show how to determine the structure of the Abelian group of the ASM in terms of the normal form decomposition of its toppling matrix. In section 4, we consider the problem of computing the rank of the Abelian group when the ASM is defined on an $L_{1} \times L_{2}$ rectangular portion of the square lattice, and show that the rank equals $L$ on an $L \times L$ square. In section 5 , we recapitulate briefly the results from general Galois theory needed by us, and use them to study the Galois group of the characteristic polynomial of the toppling matrix, and construct, in section 6, another set of algebraic functions, invariant under topplings, thereby extracting information about the Abelian group from Galois theory. This information is incomplete and we show that the method does not give the full structure of the group in general. In section 7, we remark on two other properties of the ASM on a square lattice. First, we derive a sharper upper bound on the time period of the deterministic ASM studied in [15]. The second one is an interesting connection between the structure of identity configurations on $2 L \times 2 L$ and $(2 L+1) \times(2 L+1)$ lattices. Some technical material involving detailed calculations is given in appendices A and B , while in appendix C we display some identity configurations for some small square lattices. In appendix $D$, we show how the toppling invariants discussed by Lee et al [17] are a special case of the invariants discussed in this paper.

## 2. Preliminaries and notation

The general ASM is defined on a set of $N$ sites, labelled say by integers 1 to $N$. Each site $i$ is assigned an integer variable $z_{i}$, called the height of the sandpile at site $i$. The time evolution of the sandpile is defined in terms of the following two processes:
(i) Addition of particles. We choose a site at random, and increase its height by 1 , while the heights at other sites remain unchanged. The probability of choosing the $k$ th site is denoted by $p_{k}$. For simplicity we take all $p_{k}$ 's to be non-zero. (This condition is needed to ensure the uniqueness of the steady state.)
(ii) Relaxation. If the height at some site $j$ equals or exceeds a prespecified value $z_{\text {crit }}(j)$, it is said to be unstable and topples, losing some grains of sand which either fall on other 'neighbouring' sites (with a consequent increase in height), or drop out of the system. The updating of heights is specified in terms of an $N \times N$ integer matrix $\Delta$, called the toppling matrix and satisfying $\Delta_{i i} \geqslant 0$ and $\Delta_{i j} \leqslant 0$ for $i \neq j$. If there is a toppling at some site $j$, all heights $\left\{z_{i}\right\}$ are updated according to the rule

$$
\begin{equation*}
\text { If } z_{j} \geqslant z_{\text {crit }}(j) \text { then } z_{i} \rightarrow z_{i}-\Delta_{i j} \text { for all } i . \tag{2.1}
\end{equation*}
$$

We may assume that

$$
\begin{equation*}
z_{\mathrm{crit}}(j)=\Delta_{j j} \tag{2.2}
\end{equation*}
$$

so that the values of $z_{j}$ in a stable configuration are between 0 and $\Delta_{j j}-1$.
A toppling at one site may make other sites unstable. A sequence of topplings caused by adding a single particle is called an avalanche. When all the unstable sites have toppled, we are left with a new stable configuration. This defines a single step of time evolution. At the next time, a new particle is added at a random site, and the system is allowed to relax, and so on.

It is convenient to define operators $a_{k}, k=1$ to $N$, corresponding to the process of adding a particle at site $k$, and allowing the system to relax. These operators $a_{k}$ map stable configurations into stable confgurations.

A general analysis of Abelian sandpile models was carried out in [6]. It was shown that the operators $\left\{a_{k}\right\}$ commute with each others. This allows a simple characterization of the critical steady state: only a small subset of all stable configurations occur with non-zero probability in the steady state. These are called recurrent configurations. Their number is equal to Det $\Delta$, and in the steady state they occur with equal probability. The $a_{k}$ 's map the space $\mathcal{R}$ of recurrent configurations onto itself, and are invertible on $\mathcal{R}$.

Let $G$ be the Abelian group generated by the operators $\left\{a_{i}\right\}, i=1$ to $N$. This is a finite group as these operators satisfy the closure relations [6]

$$
\begin{equation*}
\prod_{i=1}^{N} a_{i}^{\Delta_{i j}}=I \quad \text { for all } j \tag{2.3}
\end{equation*}
$$

The order of $G$, denoted by $|G|$, is equal to the number of recurrent configurations. This is a consequence of the fact that if $C$ and $C^{\prime}$ are any two recurrent configurations, then there is an element $g \in G$ such that $C^{\prime}=g C$. We thus have

$$
\begin{equation*}
|G|=|\mathcal{R}|=\operatorname{Det} \Delta . \tag{2.4}
\end{equation*}
$$

## 3. Toppling invariants and the group structure for general ASM

The space of all configurations $\left\{z_{i}\right\}$ (with non-negative heights $z_{i}$ ) constitutes a commutative semigroup over the given set of $N$ sites, with the operation given by sitewise addition of heights followed by relaxation if necessary. One can define an equivalence relation on this semigroup by declaring two configurations $\left\{z_{i}\right\}$ and $\left\{z_{i}^{\prime}\right\}$ equivalent (under toppling) if and only if there exist $N$ integers $n_{j}$ such that

$$
\begin{equation*}
z_{i}^{\prime}=z_{i}-\sum_{j} \Delta_{i j} n_{j} \quad \text { for all } i \tag{3.1}
\end{equation*}
$$

Each equivalence class contains one and only one recurrent configuration. Indeed with any configuration $\left\{z_{i}\right\}$, one can associate a recurrent configuration $C$ by letting

$$
\begin{equation*}
C\left[z_{i}\right]=\prod_{i} a_{i}^{z_{i}} C^{*} \tag{3.2}
\end{equation*}
$$

where $C^{*}$ is any fixed recurrent configuration. Then if $\left\{z_{i}\right\}$ and $\left\{z_{i}^{\prime}\right\}$ are related as in (3.1), one easily checks that $C\left[z_{i}\right]=C\left[z_{i}^{\prime}\right]$ on account of the relations (2.3). (Note that with the definition (3.1), two stable configurations may be equivalent under toppling.)

Toppling invariants are scalar functions defined on the space of all configurations of the sandpile, such that they take the same value for any two configurations which are equivalent under toppling. Toppling invariants which are linear in the height variables, and which are conserved modulo various integers were first introduced by Lee and Tzeng [17] in the context of specific deterministic one-dimensional ASM. Their results will be shown to be particular cases of our general construction (see appendix D).

Given the toppling matrix $\Delta$, we define $N$ rational functions $Q_{i}(i=1$ to $N)$ by setting

$$
\begin{equation*}
Q_{i}\left(\left\{z_{j}\right\}\right)=\sum_{j} \Delta_{i j}^{-1} z_{j} \bmod 1 \tag{3.3}
\end{equation*}
$$

It is straightforward to prove that the functions $Q_{i}$ are toppling invariants. Indeed, under toppling at site $k$, the configuration $C \equiv\left\{z_{j}\right\}$ changes to $C^{\prime} \equiv\left\{z_{j}^{\prime}=z_{j}-\Delta_{j k}\right\}$, and from the linearity of the functions $Q_{i}$ in the height variables, we have

$$
\begin{equation*}
Q_{i}\left(C^{\prime}\right)=Q_{i}(C)-\sum_{j} \Delta_{i j}^{-1} \Delta_{j k}=Q_{i}(C) \bmod 1 \tag{3.4}
\end{equation*}
$$

The functions $Q_{i}$ take rational values, but they are easily made integer-valued upon multiplication by some adequate integer. Being toppling invariants, the functions $Q_{i}$ can be used to label the recurrent configurations, and thus the space $\mathcal{R}$ can be replaced by the set of $N$-uples ( $Q_{1}, Q_{2}, \ldots, Q_{N}$ ). However, the labelling by the $Q_{i}$ 's is generally overcomplete, as they are not all independent. A simple example will readily establish that.

Throughout this paper, we shall consider in detail the special case of the ASM defined on an $L_{1} \times L_{2}$ rectangular portion of a two-dimensional square lattice. We choose the toppling matrix to be the discrete Laplacian, the diagonal entries of which are given by $\Delta_{i i}=4$, and the off-diagonal entries $\Delta_{i j}=-1$ or 0 according to whether the sites $i$ and $j$ are nearest neighbours or not. Note that this implies open boundary conditions on all four boundaries of the rectangle, and thus any toppling there involves a loss of sand.

To be specific, consider the case $L_{1}=L_{2}=2$, with the configurations specified as $\left(\begin{array}{ll}z_{1} & z_{2} \\ z_{3} & z_{4}\end{array}\right)$. In this case $\Delta$ is a $4 \times 4$ matrix of determinant 192, and we find
$\Delta=\left(\begin{array}{cccc}4 & -1 & -1 & 0 \\ -1 & 4 & 0 & -1 \\ -1 & 0 & 4 & -1 \\ 0 & -1 & -1 & 4\end{array}\right) \quad$ and $\quad \Delta^{-1}=\frac{1}{24}\left(\begin{array}{llll}7 & 2 & 2 & 1 \\ 2 & 7 & 1 & 2 \\ 2 & 1 & 7 & 2 \\ 1 & 2 & 2 & 7\end{array}\right)$.
The definition (3.3) yields the four invariants

$$
\begin{align*}
Q_{1} & =\frac{1}{24}\left(7 z_{1}+2 z_{2}+2 z_{3}+z_{4}\right) \bmod 1  \tag{3.6a}\\
Q_{2} & =\frac{1}{24}\left(2 z_{1}+7 z_{2}+z_{3}+2 z_{4}\right) \bmod 1  \tag{3.6b}\\
Q_{3} & =\frac{1}{24}\left(2 z_{1}+z_{2}+7 z_{3}+2 z_{4}\right) \bmod 1  \tag{3.6c}\\
Q_{4} & =\frac{1}{24}\left(z_{1}+2 z_{2}+2 z_{3}+7 z_{4}\right) \bmod 1 \tag{3.6d}
\end{align*}
$$

They satisfy three linear relations

$$
\begin{align*}
& Q_{1}+Q_{2}=0 \bmod 1 / 8  \tag{3.7a}\\
& Q_{3}=4 Q_{1}-Q_{2} \bmod 1  \tag{3.7b}\\
& Q_{4}=-Q_{1}+4 Q_{2} \bmod 1 \tag{3.7c}
\end{align*}
$$

from which one sees that only two independent invariants remain, which we choose to be integer-valued for convenience:

$$
\begin{align*}
& I_{1} \equiv 24 Q_{1}=7 z_{1}+2 z_{2}+2 z_{3}+z_{4} \bmod 24  \tag{3.8a}\\
& I_{2} \equiv 8\left(Q_{1}+Q_{2}\right)=3 z_{1}+3 z_{2}+z_{3}+z_{4} \bmod 8 \tag{3.8b}
\end{align*}
$$

These two invariants provide a complete labelling for the space of recurrent configurations, of cardinality 192.

As this very simple example already shows, for an arbitrary $N \times N$ matrix $\Delta$, we construct $N$ toppling invariants $Q_{i}$, but they are generally not independent. It thus seems desirable to isolate a minimal set of independent invariants, as we did above to obtain the mininal set given by ( $I_{1}, I_{2}$ ). We now show how this can be done for an arbitrary ASM using the classical theory of the Smith normal form for integer matrices.

In current mathematics literature, as in Jacobson [18], this is often discussed for matrices with entries in a principal ideal domain (the ring of integers being a prime example). From theorems 3.8 and 3.9 of Jacobson, any non-singular $N \times N$ matrix $\Delta$ can be written in the form

$$
\begin{equation*}
\Delta=A D B \tag{3.9}
\end{equation*}
$$

where $A$ and $B$ are $N \times N$ integer matrices with determinants $\pm 1$, and $D$ is a diagonal matrix

$$
\begin{equation*}
D_{i j}=d_{i} \delta_{i j} \tag{3.10}
\end{equation*}
$$

such that
(i) $d_{i}$ is a multiple of $d_{l+1}$ for all $i=1$ to $N-1$, and
(ii) $d_{i}=e_{i-1} / e_{i}$, where $e_{i}$ stands for the greatest common divisor of the determinants of all the $(N-i) \times(N-i)$ submatrices of $\Delta$ (set $e_{N}=1$ ).
Thus the matrix $D$ is uniquely determined by $\Delta$, but the matrices $A$ and $B$ are far from unique. The integers $d_{i}$ are called the elementary divisors of $\Delta$. An efficient algorithm to compute them can be found in Cohen [19].

In terms of the decomposition (3.9), we define the set of scalar functions $I_{i}(C)$ by

$$
\begin{equation*}
I_{i}(C) \rightleftharpoons \sum_{j}\left(A^{-1}\right)_{i j} z_{j} \bmod d_{i} \tag{3.11}
\end{equation*}
$$

Due to the unimodularity of $A$, these functions are integer-valued. As argued for the $Q$ 's in equation (3.4), we see using (3.9) that the $I_{i}(C)$ so defined are invariant under the toppling of any site. Clearly, only those invariants $I_{i}$ with $d_{i} \neq 1$ are non-trivial. Note that they are written in terms of $A$, and hence not unique. It will be obvious from the discussion below that the set of non-trivial $I_{i}$ is always minimal and complete.

As suggested by the notation used in the above example, this second set $\left\{I_{i}\right\}$ is a minimal set of invariants chosen from the overcomplete set $\left\{Q_{i}\right\}$. The Smith decomposition precisely shows how to combine the overcomplete invariants $Q_{i}$ so as to obtain a complete set. Indeed, it is easily checked that the $I_{l}$ 's can be written in terms of the $Q_{i}$ 's as

$$
\begin{equation*}
I_{i}=\sum_{j} d_{i} B_{i j} Q_{j} \bmod d_{i} \tag{3.12}
\end{equation*}
$$

For instance, in the above $2 \times 2$ example, one finds the following Smith decomposition (written in the form $A^{-1} \Delta=D B$ ):

$$
\left(\begin{array}{cccc}
7 & 2 & 2 & 1  \tag{3.13}\\
3 & 3 & 1 & 1 \\
-1 & 0 & 0 & 0 \\
-2 & 0 & -1 & 0
\end{array}\right) \Delta=\operatorname{diag}(24,8,1,1)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
1 & 1 & 0 & 0 \\
-4 & 1 & 1 & 0 \\
-7 & 2 & -2 & 1
\end{array}\right) .
$$

From the first two rows of $A^{-1}$ and (3.11), one recovers the two invariants $I_{1}, I_{2}$ of (3:8). From the matrix $B$ and (3.12), one finds the relations (3.8) between ( $I_{1}, I_{2}$ ) and ( $Q_{1}, Q_{2}$ ), while the trivial character of $I_{3}, I_{4}$ leads to the linear relations ( $3.7 b, c$ ) among the $Q$ 's.

Let us now show that the set $\left\{I_{i}\right\}$ not only forms a complete set of toppling invariants, but that they also determine the structure of the Abelian group $G$.

Let $g$ be the number of $d_{i}>1$. With each recurrent configuration of the ASM, we associate a $g$-uple ( $I_{1}, I_{2}, \ldots, I_{g}$ ), where $0 \leqslant I_{i}<d_{i}$. The total number of distinct $g$-uples is $\prod_{i=1}^{g} d_{i}=|G|$ (since $A$ and $B$ in the Smith decomposition of $\Delta$ are unimodular).

We first show that this mapping from the set of recurrent configurations to $g$-uples is one-to-one. Let us define operators $e_{i}$ by the equation

$$
\begin{equation*}
e_{i}=\prod_{j=1}^{N} a_{j}^{A_{j i}} \quad 1 \leqslant i \leqslant g \tag{3.14}
\end{equation*}
$$

Acting on a fixed recurrent configuration $C^{*}=\left\{z_{j}\right\}, e_{i}$ yields a new recurrent configuration, equivalent under toppling to the configuration $\left\{z_{j}+A_{j i}\right\}$. If the $g$-uple corresponding to $C^{*}$ is ( $I_{1}^{*}, I_{2}^{*}, \ldots, I_{g}^{*}$ ), it is easy to see from (3.11) that $e_{i} C^{*}$ gives a configuration with toppling invariants $I_{k}=I_{k}^{*}+\delta_{i k}$. By operating with these operators $\left\{e_{i}\right\}$ sufficiently many times on $C^{*}$, all $|G|$ values for the $g$-uple $\left(I_{1}, I_{2}, \ldots, I_{g}\right)$ are obtainable. Thus, there is at least one recurrent configuration corresponding to any $g$-uple ( $I_{1}, I_{2}, \ldots, I_{g}$ ). As the total number of recurrent configurations exactly equals the total number of $g$-uples, we see that there is a one-to-one correspondence between the $g$-uples ( $I_{1}, I_{2}, \ldots, I_{g}$ ) and the recurrent configurations of the ASM.

To express the operators $a_{j}$ in terms of $e_{i}$, we need to invert the transformation (3.14). This is easily seen to be

$$
\begin{equation*}
a_{j}=\prod_{i=1}^{g} e_{i}^{\left(A^{-1}\right)_{i j}} \quad 1 \leqslant j \leqslant N . \tag{3.15}
\end{equation*}
$$

Thus the operators $e_{i}$ generate the whole of $G$. Since $e_{i}$ acting on a configuration increases $I_{i}$ by 1 , leaving the other invariants unchanged, and since $I_{i}$ is only defined modulo $d_{i}$, we see that

$$
\begin{equation*}
e_{i}^{d_{i}}=I \quad \text { for } i=1 \text { to } g . \tag{3.16}
\end{equation*}
$$

Note that the definition (3.14) for $e_{i}$ makes sense for $i$ between $g+1$ and $N$, and implies relations among the $a_{j}$ operators ( $e_{i}=I$ from (3.16)).

This shows that $G$ has the canonical decomposition as a product of cyclic groups

$$
\begin{equation*}
G=Z_{d_{1}} \times Z_{d_{2}} \times \ldots \times Z_{d_{s}} \tag{3.17}
\end{equation*}
$$

with the $d_{i}$ 's defined in (3.10). We thus have shown that the generators and the group structure of $G$ for an arbitrary ASM can be entirely determined from its toppling matrix $\Delta$, through its normal from decomposition (3.9).

The invariants $\left\{I_{i}\right\}$ also provide a simple additive representation of the group $G$. As mentioned at the beginning of this section, one can define a binary operation of 'addition' (denoted by $\oplus$ ) on the space $\mathcal{R}$ of recurrent configurations by adding heights sitewise, and then allowing the resultant configuration to relax. From the linearity of the $I_{i}$ 's in the height variables, and their invariance under toppling, it is clear that under this addition of configurations, the $I_{i}$ also simply add, i.e. for any recurrent configurations $C_{1}$ and $C_{2}$, one has.

$$
\begin{equation*}
I_{i}\left(C_{1} \oplus C_{2}\right)=I_{i}\left(C_{1}\right)+I_{i}\left(C_{2}\right) \bmod d_{i} . \tag{3.18}
\end{equation*}
$$

The $I_{i}$ 's provide a complete labelling of $\mathcal{R}$. There is a unique recurrent configuration, denoted by $C_{\mathrm{id}}$, for which all $I_{i}\left(C_{i d}\right)$ are zero. Also, each recurrent $C$ has a unique inverse $-C$, also recurrent, and determined by $I_{i}(-C)=-I_{i}(C) \bmod d_{i}$. Therefore the addition $\oplus$ is a group law on $\mathcal{R}$, with identity given by $C_{\mathrm{id}}$. A simple recursive algorithm to compute $C_{\text {id }}$ has been given by Creutz [13].

There exists a one-to-one correspondence between the recurrent configurations of ASM, and the elements of the group $G$ : we associate with the group element $g \in G$, the recurrent configuration $g C_{\mathrm{id}}$. It is then easy to see from (3.18) that for all $g, g^{\prime} \in G$

$$
\begin{equation*}
g C_{\mathrm{id}} \oplus g^{\prime} C_{\mathrm{id}}=\left(g g^{\prime}\right) C_{\mathrm{id}} . \tag{3.19}
\end{equation*}
$$

Thus the recurrent configurations with the operation $\oplus$ form a group which is isomorphic to the multiplicative group $G$, a result first proved in [13].

The invariants $\left\{I_{i}\right\}$ provide a simple labelling of all the recurrent configurations. Since a recurrent configuration can also be uniquely specified by the height variables $\left\{z_{i}\right\}$, the existence of forbidden subconfigurations in ASMS implies that these heights satisfy many inequality constraints. (Certain patterns cannot appear inside any recurrent configuration, and hence are called forbidden subconfigurations [6]; for instance, two neighbouring sites can never both have zero heights.)

## 4. The rank of $G$ for a rectangular lattice

For a general matrix $\Delta$, it is difficult to say much more about the group structure of G. In the rest of this paper, unless otherwise specified, we shall consider the special case when $\Delta$ is the toppling matrix corresponding to a finite $L_{1} \times L_{2}$ rectangle of the (two-dimensional) square lattice. Here it is more convenient to label the sites not by a single index $i$ going from 1 to $N=L_{1} L_{2}$, but by two Cartesian coordinates $(x, y)$, where $1 \leqslant x \leqslant L_{1}$ and $1 \leqslant y \leqslant L_{2}$. The toppling matrix is the discrete Laplacian, given by $\Delta(x, y ; x, y)=4, \Delta\left(x, y ; x^{\prime}, y^{\prime}\right)=-1$ if the sites $(x, y)$ and $\left(x^{\prime}, y^{\prime}\right)$ are nearest neighbours
(i.e. $\left|x-x^{\prime}\right|+\left|y-y^{\prime}\right|=1$ ), and zero otherwise. Without loss of generality, we can assume that $L_{1} \geqslant L_{2}$.

The relations satisfied by the particle addition operators $a(x, y)$ can be written in the form

$$
\begin{equation*}
a(x+1, y)=a^{4}(x, y) a^{-1}(x, y+1) a^{-1}(x, y-1) a^{-1}(x-1, y) \tag{4.1}
\end{equation*}
$$

where we adopt the convention that
$a(x, 0)=a\left(x, L_{2}+1\right)=a(0, y)=a\left(L_{1}+1, y\right)=I \quad$ for all $x, y$.
Equations (4.1) can be recursively solved to express any operator $a(x, y)$ as a product of powers of $a(1, y)$. Therefore the group $G$ can be generated by the $L_{2}$ operators $a(1, y)$. Denoting the rank of $G$ (minimal number of generators) by $g$, this implies that

$$
\begin{equation*}
g \leqslant L_{2} \tag{4.3}
\end{equation*}
$$

In the special case of the linear chain, $L_{2}=1$, we see that $g=1$, and thus $G$ is cyclic.
From (4.1), $a\left(L_{1}+1, y\right)$ can also be expressed as a product of powers of $a(1, y)$, say

$$
\begin{equation*}
a\left(L_{1}+1, y\right)=\prod_{y^{\prime}} a\left(1, y^{\prime}\right)^{n_{y}} \tag{4.4}
\end{equation*}
$$

where the $n_{y y^{\prime}}$ are integers which depend on $L_{1}$ and $L_{2}$, and which can be explicitly determined by solving the recurrence relation (4.1). The condition $a\left(L_{1}+1, y\right)=I$ then leads to the closure relations

$$
\begin{equation*}
\prod_{y^{\prime}=1}^{L_{2}} a\left(1, y^{\prime}\right)^{n_{y} y^{\prime}}=I \quad \text { for all } 1 \leqslant y \leqslant L_{2} \tag{4.5}
\end{equation*}
$$

Equations (4.5) give a presentation of $G$, the structure of which can be determined from the normal form decomposition of the $L_{2} \times L_{2}$ matrix $n_{y y^{\prime}}$. This is considerably easier to handle than the normal form decomposition of the much larger matrix $\Delta$ needed for an arbitrary ASM.

Even though this is a real computational improvement, the actual calculation, for arbitrary $L_{1}$, of the rank of $G$ is non-trivial. Even in the simplest case $L_{2}=2$, it depends in a complicated way on the number-theoretic properties of $L_{1}$. Details of this case are given in appendix A .

As to the square lattice where $L_{1}=L_{2}=L$, using the above algorithm we find the following structures of $G$ for the first values of $L$

$$
\begin{align*}
& L=2: G=Z_{24} \times Z_{8}  \tag{4.6a}\\
& L=3: G=Z_{224} \times Z_{112} \times Z_{4}  \tag{4.6b}\\
& L=4: G=Z_{6600} \times Z_{1320} \times Z_{8} \times Z_{8}  \tag{4.6c}\\
& L=5: G=Z_{102960} \times Z_{102960} \times Z_{48} \times Z_{16} \times Z_{4} \tag{4.6d}
\end{align*}
$$

This suggests, and it is, in fact, not difficult to prove, that for an $L \times L$ square

$$
\begin{equation*}
g=L \quad \text { for } L_{1}=L_{2}=L \tag{4.7}
\end{equation*}
$$

The idea of the proof of (4.7) is to use the fact that in this case, the matrix $\Delta$ has exactly $L$ linearly independent eigenvectors of eigenvalue 4 , of which all components can be chosen to be integers. We use them to construct $L$ independent operators (none can be expressed as product of powers of the others) $U_{k}, 0 \leqslant k \leqslant L-1$, such that

$$
\begin{equation*}
U_{k}^{4}=I \tag{4.8}
\end{equation*}
$$

This implies that the number of generators is at least equal to $L$. Combined with the inequality (4.3), we get the result.

To write the operators $U_{k}$ explicitly, we consider the eigenvectors $n_{k}(x, y)$ of $\Delta$ of eigenvalue 4:

$$
\begin{equation*}
\sum_{\left(x^{\prime}, y^{\prime}\right)} \Delta\left(x, y ; x^{\prime}, y^{\prime}\right) n_{k}\left(x^{\prime}, y^{\prime}\right)=4 n_{k}(x, \dot{y}) \tag{4.9}
\end{equation*}
$$

There are $L$ independent solutions to equation (4.9), and a possible choice is to set $(1 \leqslant x, y \leqslant L, 0 \leqslant k \leqslant L-1)$
$n_{k}(x, y)=(-1)^{x}[\delta(y-x-k)-\delta(y+x-k)-\delta(x+y-2 L-2+k)+\delta(y-x+k)]$.

It is then easy to verify by using (4.9) and the relations (2.3) satisfied by the $a(x, y)$, that the operators $U_{k}$ defined by

$$
\begin{equation*}
U_{k}=\prod_{x, y} a(x, y)^{n_{k}(x, y)} \tag{4.11}
\end{equation*}
$$

all satisfy $U_{k}^{4}=I$. Note that a strong form of independence of the eigenvectors $n_{k}(x, y)$ must hold for the operators $U_{k}$ to be multiplicatively independent: the eigenvectors must be linearly independent modulo the lattice $\left\{\sum_{x^{\prime}, y^{\prime}} m_{x^{\prime}, y^{\prime}} \Delta\left(x, y ; x^{\prime}, y^{\prime}\right): m_{x^{\prime}, y^{\prime}}\right.$ integers $\}$, on account of the relations (2.3). The vectors (4.10) are easily seen to satisfy this condition. This completes the proof.

When $L_{1} \neq L_{2}$, the above construction is easily generalized. If $f=\operatorname{gcd}\left[\left(L_{1}+1\right),\left(L_{2}+\right.\right.$ 1)], it is easy to check that the matrix $\Delta$ has $f-1$ independent eigenvectors of eigenvalue 4. Thus we obtain in this case that the rank $g_{L_{1} \times L_{2}}$ satisfies the two inequalities

$$
\begin{equation*}
f-1 \leqslant g_{L_{1} \times L_{2}} \leqslant L_{2} \tag{4.12}
\end{equation*}
$$

That the eigenvectors of integer eigenvalue play a particular role in the above proof should be clear. An arbitrary eigenvalue is generically an algebraic number, and the corresponding eigenvector has components which are also algebraic numbers, so that (4.11) becomes meaningless. However, we show in the next two sections that one can construct toppling invariants from the spectrum of $\Delta$, which yield the proper setting to generalize the above idea to any eigenvalue. The corresponding construction is applicable to any ASM which has a diagonalizable toppling matrix.

## 5. A reminder of Galois theory

In this section, we recall the basic ideas of Galois theory, which can be used in any ASM, whatever its type. For illustrating the technique, we shall consider the familiar ASM defined on an $L \times L$ square lattice. The toppling matrix $\Delta$ is the discrete Laplacian, and its eigenvalues are easily determined and read as

$$
\begin{equation*}
\lambda_{m, n}=4-2 \cos \frac{2 \pi m}{N}-2 \cos \frac{2 \pi n}{N} \quad 1 \leqslant m, n \leqslant L \tag{5.1}
\end{equation*}
$$

where we have defined

$$
\begin{equation*}
N=2(L+1) \tag{5.2}
\end{equation*}
$$

The order of the group $G$ is given by

$$
\begin{equation*}
|G|=\prod_{m=1}^{L} \prod_{n=1}^{L} \lambda_{m, n} \tag{5.3}
\end{equation*}
$$

We define the transformations

$$
\begin{equation*}
\sigma_{s}\left(\lambda_{m, n}\right)=\lambda_{m^{\prime}, n^{\prime}} \quad s \text { coprime with } N \tag{5.4}
\end{equation*}
$$

where $m^{\prime}, n^{\prime}$ are integers, $1 \leqslant m^{\prime}, n^{\prime} \leqslant L$, such that $\cos \left(2 \pi m^{\prime} / N\right)=\cos (2 \pi s m / N)$ and $\cos \left(2 \pi n^{\prime} / N\right)=\cos (2 \pi s n / N)$. We see that these transformations act on the set of eigenvalues by permutation. Moreover, they form a group, and as

$$
\begin{equation*}
\sigma_{s} \circ \sigma_{s^{\prime}}=\sigma_{s s^{\prime}}=\sigma_{s^{\prime}} \circ \sigma_{s} \tag{5.5}
\end{equation*}
$$

the group is commutative and isomorphic to the multiplication modulo $N$ of the numbers coprime with $N$. Finally from the explicit expression of $\lambda_{m, n}$, the actions of $\sigma_{s}$ and $\sigma_{-s}$ are the same. Thus the group of all $\sigma_{s}$ is isomorphic to the group $Z_{N}^{*}$ of invertible integers modulo $N$ (those which are coprime with $N$ ) quotiented by the subgroup $\{+1,-1\}$

$$
\begin{equation*}
\mathrm{GaI}_{L}=\left\{\sigma_{s}: 1 \leqslant s \leqslant L \text { and } s \text { coprime with } N\right\} \simeq Z_{N}^{*} /\{ \pm 1\} \tag{5.6}
\end{equation*}
$$

and is of order $\frac{1}{2} \varphi(N)$. We call group (5.6) the Galois group of this ASM.
It is instructive to group the set of eigenvalues $\lambda_{m, n}$ into orbits under the Galois group (5.6). As an example, let us take $L=3$ or $N=8$. The Galois group consists of the two transformations $s=1$ (the identity transformation) and $s=3$. We find $\sigma_{3}\left(\lambda_{1,1}\right)=\lambda_{3,3}$ and $\sigma_{3}\left(\lambda_{3,3}\right)=\lambda_{1,1}$, so that $\left\{\lambda_{1,1}, \lambda_{3,3}\right\}$ form an orbit under the Galois group. Computing the product of these two eigenvalues, we find $\lambda_{1,1} \lambda_{3,3}=8$. Likewise $\sigma_{3}\left(\lambda_{2,2}\right)=\lambda_{2,2}$ shows that $\lambda_{2,2}=4$ is an orbit on its own. (Note that because of the existence of some degeneracies in the eigenvalues, finding the orbit of $\lambda_{m, n}$ is not merely finding the orbit of ( $m, n$ ) under a diagonal multiplication by all $s$. In the example at hand, to say that $\left\{\lambda_{1,1}, \lambda_{3,3}\right\}$ is the orbit of $\lambda_{1,1}$ supposes that we have checked that the two eigenvalues are different complex numbers.) Doing the same calculation for the other eigenvalues, we can rearrange the product (5.3) giving the order of $G$ as a product over the orbits ( 6 in this case) to find

$$
\begin{align*}
L=3:|G| & =\left[\lambda_{11} \lambda_{33}\right]\left[\lambda_{22}\right]\left[\lambda_{12} \lambda_{32}\right]\left[\lambda_{21} \lambda_{23}\right]\left[\lambda_{13}\right]\left[\lambda_{31}\right] \\
& =8 \cdot 4 \cdot 14 \cdot 14 \cdot 4 \cdot 4=100352 . \tag{5.7}
\end{align*}
$$

We can do the same rearrangement for any value of $L$, writing $|G|$ as a product over orbits under the Galois group,

$$
\begin{equation*}
|G|=\prod_{\text {orbits } \mathcal{O}}\left[\prod_{\lambda_{m, n} \in \mathcal{O}} \lambda_{m, n}\right] \tag{5.8}
\end{equation*}
$$

and then each subproduct in square brackets is an integer. This follows from the fact that each square bracket is, by construction, invariant under the whole of the Galois group, in which case the Galois theorem then states that it must be a rational number. Because of the special form of the numbers $\lambda_{m, n}$ (they are algebraic integers), all square brackets must even be integers. Not only is each square bracket equal to an integer, but the Galois theorem asserts that none of them can be further split: the decomposition (5.8) is the maximal factorization of $|G|$ as a product of integers with the property that each integer is the product of a certain number of eigenvalues.

Let us now put the above considerations in a more general perspective. A clear account of Galois theory can be found, for example, in [20].

Let $P(x)$ be a finite degree polynomial with integer coefficients, and let us define the algebraic extension $F[P]$ by adjoining to the field of rational numbers $Q$ all the roots of $P(x)$. (This kind of algebraic extension is technically known as a separable normal extension.) Thus, $F[P]$ is constructed in two steps: first we form the set of all rational linear combinations of the roots of $P$ and, second, we promote this set to a field by adding whatever is needed to make it closed under multiplication. Hence $F[P]$ contains all roots of $P$ and all products of roots, but it is always possible to choose a finite number of base elements (linearly independent over $Q$ ) of which $F[P]$ is the linear $Q$-span. The number of base elements is called the degree of the extension $F[P]$, which has, in general, no relation with the degree of $P$. For instance the degree of $F\left[x^{4}-2\right]$ is 8 , while that of $F\left[x^{4}-1\right]$ is 2.

Clearly, $Q$ is a subfield of $F[P]$, and it is thus a well defined problem to look for automorphisms of $F[P]$ which leave invariant every element of $Q \subset F[P]$. The group of all such automorphisms is called the Galois group of $F[P]$ with respect to $Q$, and denoted $\operatorname{Gal}(F[P] / Q)$. The effect of the Galois group on $F[P]$ is to permute the roots of $P$, although not all permutations are allowed since they must be automorphisms. An equivalent definition of $\operatorname{Gal}(F[P] / Q)$ is the set of permutations of the roots of $P(x)$ which preserve all algebraic relations among them. It can be shown that the order of the Galois group is equal to the degree of $F[P]$. We can now formulate the Galois correspondence theorem: there is a one-to-one correspondence between the subfields of $F[P]$ and the subgroups of $\operatorname{Gal}(F[P] / Q)$. The bijection goes by associating a subfield $M \subset F[P]$ with the maximal subgroup $H \subset \operatorname{GaI}(F[P] / Q)$ which leaves the elements of $M$ invariant. Hence the bigger $H$ is, the smaller $M$ is. For instance $M=Q$ is associated with $H I=\operatorname{Gal}(F[P] / Q)$ by definition of the Galois group (hence an element of $F[P]$ is in $Q$ iff it is invariant under the whole Galois group), and $M=F[P]$ is associated with the trivial subgroup $H=\{1\}$.

Before closing this digression, we make a final comment about algebraic numbers and algebraic integers. All the elements of an algebraic extension of $Q$ are algebraic numbers, which means that they all satisfy polynomial equations with integer coefficients. However, some of them may as well satisfy a polynomial equation with integer coefficients and with coefficient of the highest power equal to 1 . These elements are distinguished in the extension and are called algebraic integers. For instance, in $F\left[x^{2}-2\right]=Q[\sqrt{2}]$, the number $\sqrt{2}$ is an algebraic integer whereas $1 / \sqrt{2}$ is not. Also the eigenvalues of a matrix with integer entries are all algebraic integers. While the extension itself is a field, the subset of its algebraic
integers is only a ring. In $Q$, this amounts to the usual distinction between the rationals and the integers. An immediate consequence is that an algebraic integer of $F[P]$ which is invariant under the whole Galois group is actually an algebraic integer of $Q$, i.e. an integer of $Z$.

Let us now see how these general ideas work if we take the extension $F[\operatorname{det}(x-\Delta)]$ obtained by adjoining to $Q$ the roots of the characteristic polynomial of the toppling matrix of an ASM. We choose the ASM defined in section 3 for an $L \times L$ lattice, with $\Delta$ the discrete Laplacian, but clearly the same analysis can be done in any model.

We thus consider the algebraic extension $F_{L} \equiv F\left[\operatorname{det}\left(x-\Delta_{L}\right)\right]$ obtained by adding to $Q$ the roots $\lambda_{m, n}$ given in (5.1),

$$
\begin{equation*}
\lambda_{m, n}=4-2 \cos \frac{2 \pi m}{N}-2 \cos \frac{2 \pi n}{N} \quad 1 \leqslant m, n \leqslant L \tag{5.9}
\end{equation*}
$$

Since $F_{L}$ contains all rational combinations of the $\lambda_{m, n}$, it contains $\cos 2 \pi m / N$ for all $1 \leqslant m \leqslant L$ since

$$
\begin{equation*}
\cos \frac{2 \pi m}{N}=2-\frac{1}{2 L} \sum_{n=1}^{L} \lambda_{m, n} \tag{5.10}
\end{equation*}
$$

We also see from

$$
\begin{equation*}
\cos \frac{2 \pi m}{N} \cos \frac{2 \pi n}{N}=\frac{1}{2} \cos \frac{2 \pi(m+n)}{N}+\frac{1}{2} \cos \frac{2 \pi(m-n)}{N} \tag{5.11}
\end{equation*}
$$

that any product of $\lambda_{m, n}$ can be expressed as a rational combination of 1 and the numbers (5.10), which shows that any element of $F_{L}$ can be written as a rational combination of $1, \cos (2 \pi / N), \cos (4 \pi / N), \ldots, \cos (2 \pi L / N)$. However, this writing is not unique as these $L+1$ cosines are not independent over $Q$, but satisfy

$$
\begin{equation*}
\sum_{i=0}^{p-1} \cos \frac{2 \pi\left(a p^{k}+b N / p^{k}+\mathrm{i} N / p\right)}{N}=0 \quad 0 \leqslant a \leqslant \frac{N}{p^{k}}-1,0 \leqslant b \leqslant p^{k-1}-1 \tag{5.12}
\end{equation*}
$$

for every prime power $p^{k}$ entering the prime decomposition of $N$. Using these relations, it can be shown that only $\frac{1}{2} \varphi(N)$ among the $L+1$ cosines are independent over $Q$. Thus the extension $F_{L}$, usually denoted by $Q(\cos (2 \pi / N))$, has degree $\frac{1}{2} \varphi(N)$. A basis of $F_{L}$ over $Q$ can be taken to be $\left\{\cos (2 \pi m / N): 0 \leqslant m \leqslant \frac{1}{2} \varphi(N)-1\right\}$.

To find the Galois group of $F_{L}$, we first note the following identities

$$
\begin{equation*}
\cos \frac{2 \pi m}{N}=T_{m}\left(\cos \frac{2 \pi}{N}\right) \tag{5.13}
\end{equation*}
$$

where $T_{m}$ is the $m$ th Chebyshev polynomial of the first kind. Equation (5.13) shows that the Galois transformation of $\cos (2 \pi m / N), m \geqslant 2$, is determined by that of $\cos (2 \pi / N)$. The next step is to notice that, since the Galois group acts on the roots $\lambda_{m, n}$ by permutations, the Galois transformations of $\cos (2 \pi / N)$ must be in the set of $\cos (2 \pi s / N), 1 \leqslant s \leqslant L$. Hence suppose that $\sigma_{s}(\cos (2 \pi / N))=\cos (2 \pi s / N)$ is a Galois transformation for $s$ between 1 and $L$. Being a group element, $\sigma_{s}$ must have an inverse which we denote by $\sigma_{s^{\prime}}=\sigma_{s}^{-1}$. It satisfies $\sigma_{s^{\prime}}(\cos (2 \pi s / N))=\cos \left(2 \pi s s^{\prime} / N\right)=\cos (2 \pi / N)$, hence $s s^{\prime}= \pm 1 \bmod N$. This is impossible if $s$ and $N$ have a common factor, while if $s$ is coprime with $N$, there is a
unique $s^{\prime}$ between 1 and $L$ satisfying either $s^{\prime}=s^{-1} \bmod N$ or $s^{\prime}=-s^{-1} \bmod N$. Thus the Galois group contains at least the transformations $\sigma_{s}$ for $s$ between 1 and $L$ and coprime with $N$. The number of such $s$ is equal to $\frac{1}{2} \varphi(N)$, and since this must also be the order of the Galois group (equal to the degree of $F_{L}$ ), one recovers the symmetry group of (5.6)

$$
\begin{equation*}
\mathrm{Gal}_{L}=\left\{\sigma_{s}: 1 \leqslant s \leqslant L \text { and }(s, N)=1\right\} \tag{5.14}
\end{equation*}
$$

We have completely identified the extension $F_{L}$ and determined its Galois group. The last point concerns the ring of algebraic integers of $F_{L}$. This is a far more difficult question, and we just quote the result: the algebraic number $w=a_{0}+\sum_{m=1}^{[\varphi(N)-1] / 2} 2 a_{m} \cos (2 \pi m / N)$ of $F_{L}$ is an algebraic integer if and only if all coefficients $a_{m}$ are integers of $Z$. This ring of integers is usually denoted by $Z(\cos (2 \pi / N))$. In particular, as noted above, the eigenvalues $\lambda_{m, n}$ are all algebraic integers.

Looking back at the decomposition (5.8) of $|G|$ as a product of orbits of eigenvalues under the Galois group, all the properties we mentioned easily follow from the above facts. The product of eigenvalues belonging to an orbit is, by definition, Galois invariant and thus must be an rational number. Since the eigenvalues are algebraic integers, this rational number must, in fact, be an integer. Finally the decomposition is maximal by definition of the orbits as the smallest subsets invariant under the Galois group.

Similarly to the decomposition of $G=\operatorname{Det} \Delta$, we may consider the factorization of the characteristic polynomial of $\Delta$ in factors which are all invariant under the Galois group

$$
\begin{equation*}
P_{L}(x)=\prod_{m, n=1}^{L}\left(x-\lambda_{m, n}\right)=\prod_{\text {orbits } \mathcal{O}}\left[\prod_{\lambda_{m, n} \in \mathcal{O}}\left(x-\lambda_{m, n}\right)\right] \tag{5.15}
\end{equation*}
$$

In this case, the Galois theory states that the polynomial within each square bracket, call it $P_{\mathcal{O}}(x)$, has integer coefficients, and that the decomposition (5.15) is maximal, namely one cannot perform a further splitting into polynomials with rational coefficients (or, in other words, all $P_{\mathcal{O}}(x)$ are irreducible over the field of rationals).

Examples of decomposition (5.15) for small values of $L$ are:

$$
\begin{align*}
& P_{2}(4-x)=x^{2}\left(x^{2}-4\right)  \tag{5.16a}\\
& P_{3}(4-x)=-x^{3}\left(x^{2}-2\right)^{2}\left(x^{2}-8\right)  \tag{5.16b}\\
& P_{4}(4-x)=x^{4}(x-1)^{2}(x+1)^{2}\left(x^{2}-5\right)^{2}\left(x^{2}-2 x-4\right)\left(x^{2}+2 x-4\right) \tag{5.16c}
\end{align*}
$$

$$
\begin{align*}
P_{7}(4-x)=- & x^{7}\left(x^{2}-2\right)^{2}\left(x^{2}-8\right)\left(x^{4}-4 x^{2}+2\right)^{2}\left(x^{4}-8 x^{2}+8\right)^{2}\left(x^{4}-16 x^{2}+32\right) \\
& \times\left(x^{4}-8 x^{2}-8 x-2\right)^{2}\left(x^{4}-8 x^{2}+8 x-2\right)^{2} \tag{5.16d}
\end{align*}
$$

The number $F(L)$ of factors in the decompositions (5.8) or (5.15) can be explicitly computed. The behaviour of $F(L)$ is rather chaotic, as can be checked from table 1 , where we give some of its values. It strongly depends on the prime decomposition of $(L+1)$, and for large $L, F(L)$ increases linearly with $L$. An analytic formula for $F(L)$ is given in appendix B.

Finally, let us note that the Galois group (5.14) also acts on the $L \times L$ lattice on which the ASM is defined. This action is defined by $\sigma_{s}(x, y)=\left(x^{\prime}, y^{\prime}\right)$ where $\cos \left(2 \pi x^{\prime} / N\right)=$

Table 1. Number $F(L)$ of irreducible factors in $\operatorname{det}(x-\Delta)$.

| $L$ | $F(L)$ | $L$ | $F(L)$ | $L$ | $F(L)$ | $L$ | $F(L)$ |
| :--- | :---: | ---: | :--- | :--- | :---: | :--- | :--- |
| 1 | 1 | 6 | 16 | 21 | 74 | 26 | 112 |
| 2 | 4 | 7 | 19 | 22 | 64 | 27 | 112 |
| 3 | 6 | 8 | 28 | 23 | 122 | 28 | 84 |
| 4 | 12 | 9 | 34 | 24 | 88 | 29 | 192 |
| 5 | 18 | 10 | 28 | 25 | 90 | 30 | 88 |

$\cos (2 \pi s x / N)$ and $\cos \left(2 \pi y^{\prime} / N\right)=\cos (2 \pi y s / N)$. The operator $\sigma_{s}$ permutes the lattice sites. In turn, this induces an action on the set of all configurations of the ASM by setting

$$
\begin{equation*}
\sigma_{s} C: \sigma_{s}\left(z_{i}\right)=z_{\sigma_{s}(i)} \tag{5.17}
\end{equation*}
$$

So we can also look at the Galois transformations either as automorphisms of the lattice or as automorphisms of the set of all configurations of the ASM. In fact, they even define automorphisms of $\mathcal{R}$, because, although $\sigma_{s} C$ is not necessarily recurrent even if $C$ is, there is a unique recurrent configuration equivalent to $\sigma_{s} C$.

## 6. Toppling invariants from eigenvectors

Using the basic results of Galois theory briefly recalled in section 5 , we study here toppling invariants constructed from the (left) eigenvectors of $\Delta$. This can be done for any ASM with diagonalizable toppling matrix, but for definiteness, we proceed with the ASM defined in section 4, for a square $L \times L$ lattice.

Thus $\Delta$ is the discrete Laplacian, with the eigenvalues $\lambda_{m, n}$ given in (5.1), and for which the eigenvectors read as

$$
\begin{equation*}
v_{m, n}(x, y)=\frac{\sin 2 \pi m x / N}{\sin 2 \pi m / N} \frac{\sin 2 \pi n y / N}{\sin 2 \pi n / N} \quad 1 \leqslant m, n \leqslant L . \tag{6.1}
\end{equation*}
$$

The normalization is purely conventional, and will be explained below.
For each eigenvector $v_{m, n}$, we define an algebraic toppling invariant by taking its inner product with the $z$ vector of an arbitrary stable configuration, and by evaluating the result modulo the eigenvalue $\lambda_{m, n}$ of $v_{m, n}$

$$
\begin{equation*}
A_{m, n}(C)=\sum_{x, y} v_{m, n}(x, y) z_{x, y} \bmod \lambda_{m, n} \tag{6.2}
\end{equation*}
$$

Since the eigenvalues and the eigenvector components are, in general, algebraic numbers in the algebraic extension $F[\operatorname{det}(x-\Delta)]$, equal to $Q(\cos (2 \pi / N))$ in this case, the congruence (6.2) must be understood in that extension. Moreover, modulo $\lambda_{m, n}$ means modulo all multiples $w \lambda_{m, n}$ with $w$ in the ring of algebraic integers of $Q(\cos (2 \pi / N))$.

Clearly under the toppling at site $\left(x_{0}, y_{0}\right)$, we have $A_{m, n} \longrightarrow A_{m, n}+\lambda_{m, n} v_{m, n}\left(x_{0}, y_{0}\right)$, so that $A_{m, n}(C)$ is invariant provided $v_{m, n}\left(x_{0}, y_{0}\right)$ is an algebraic integer for all $x_{0}, y_{0}$. This was precisely the reason for choosing the normalization (6.1). On account of the fact that, for integer $x, \sin \alpha x / \sin \alpha$ is a polynomial of $2 \cos \alpha$ with integer coefficients, the components of $v_{m, n}$ are indeed algebraic integers of $Q(\cos (2 \pi / N)$ ) (see section 5).

Thus the invariants $A_{m, n}$ are all valued in $Z(\cos (2 \pi / N))$ where the congruence (6.2) is to be taken. To give an explicit example, consider $L=3, N=8$. The relevant extension is
$Q(\cos (2 \pi / 8))=Q(\sqrt{2})$. Choosing for instance $m=n=1$, one obtains $\lambda_{1,1}=4-2 \sqrt{2}$, and
$A_{1,1}(C)=z_{1,1}+\sqrt{2} z_{1,2}+z_{1,3}+\sqrt{2} z_{2,1}+2 z_{2,2}+\sqrt{2} z_{2,3}+z_{3,1}+\sqrt{2} z_{3,2}+z_{3,3} \bmod 4-2 \sqrt{2}$.

It can be seen that this invariant takes eight different values by setting $z_{1,1}=k, z_{1,2}=l$ for $k=0,1,2,3$ and $l=0,1$, and all other $z_{x, y}=0$. Thus the set of values taken by $A_{1,1}$ can be presented as $\{k+l \sqrt{2}: k=0,1,2,3$ and $l=0,1\}$, which is an additive group isomorphic to $Z_{4} \times Z_{2}$. Also recall from section 5 that $\lambda_{1,1}$ belongs to the orbit $\mathcal{O}=\left\{\lambda_{1,1}, \lambda_{3,3}\right\}$ under the Galois group, and that the associated irreducible polynomial is $P_{O}(x)=\left(x-\lambda_{1,1}\right)\left(x-\lambda_{3,3}\right)=x^{2}-8 x+8$. So we have in this case that $P_{O}(0)=8$ is the number of values taken by the invariant $A_{1,1}$.

It seems that we have defined $L^{2}$ invariants $A_{m, n}$, one for each eigenvalue. However, some of them are related by Galois transformations and are thus not independent. Indeed if $\lambda_{m, n}$ and $\lambda_{m^{\prime}, n^{\prime}}$ are in the same Galois orbit, the corresponding eigenvectors are related by a Galois transformation $\sigma$, and so are the corresponding invariants, $A_{m^{\prime}, n^{\prime}}=\sigma\left(A_{m, n}\right)$. So in fact we associate an invariant not to each eigenvalue, but rather to each orbit $\mathcal{O}$ under the Galois group or, equivalently, to each irreducible factor $P_{\mathcal{O}}(x)$ of the characteristic polynomial of $\Delta$. We thus have $F(L)$ (the number of orbits) algebraic invariants. Moreover, because of the normalization we chose in (6.1) for the eigenvectors, each invariant is linear in the height of the left top corner $(x=y=1)$ with coefficient 1 :

$$
\begin{equation*}
A_{m, n}(C)=z_{1,1}+\cdots \bmod \lambda_{m, n} . \tag{6.4}
\end{equation*}
$$

Since 1 is not an eigenvalue, we see that none of the $A_{m, n}$ is trivial. However, the exact number of values they take is, in general, difficult to estimate. In most cases, the number of values taken by $A_{m, n}$ is bounded by $\left|P_{\mathcal{O}}(0)\right|$, where $\mathcal{O}$ is the orbit comprising $\lambda_{m, n}$, but can also be larger (when the extension containing the eigenvalue $\lambda_{m, n}$ is a proper subfield of that containing the eigenvector components $v_{m, n}$ ).

The next question is whether the invariants $A_{m, n}$ are linearly independent. This is easily answered negatively. There are $F(L)>L$ invariants and we have shown that none of them is trivial, so that if they were independent, the group $G$ would be the product of at least $F(L)$ cyclic groups. (Remember from section 3 that toppling invariants provide additive representations of $G$.) This is of course impossible since, from section 4 , we know that $G$ is the product of exactly $L$ cyclic factors.

Although not linearly independent, we can still consider the problem of their completeness. We do not have general results regarding this question, but generically, and this is confirmed by concrete examples, we do not expect them to provide a complete labelling, because the elementary divisors of $\Delta$, determining the structure of $G$, are not related to its spectral decomposition. It is however possible to find, by the following algorithm, a complete set of algebraic invariants, analogous to the complete set of integer invariants $I_{i}$ defined in (3.11) in terms of the Smith normal form.

We note that since the invariants involve congruences, there is no need to use the exact eigenvectors. We can define invariants by setting

$$
\begin{equation*}
A_{\lambda}(C)=\sum_{x, y} v_{\lambda}(x, y) z_{x, y} \bmod \lambda \tag{6.5}
\end{equation*}
$$

where we now take for $v_{\lambda}$ a vector with the property that

$$
\begin{equation*}
v_{\lambda} \cdot \Delta=0 \bmod \lambda \tag{6.6}
\end{equation*}
$$

With no loss of generality, we can assume that the components of $v_{\lambda}$ are in the ring of integers of some algebraic extension (otherwise $\lambda$ is simply rescaled), so that we can see the vectors $v_{\lambda}$ as left eigenvectors of zero eigenvalue over a finite ring. We will call them modular eigenvectors (with zero eigenvalue). We look for a set of such vectors, as small as possible and with as large values of $\lambda$ as possible, in order to generate a complete set of invariants with the fewest invariants. To avoid irrelevant overall factors (if $v_{\lambda}$ is a modular eigenvector modulo $\lambda, 2 v_{\lambda}$ is a modular eigenvector modulo $2 \lambda$ ), we may assume that at least one component of $v_{\lambda}$ is coprime with $\lambda$, which we can then choose equal to 1 . The system (6.6) is overdetermined since it yields $L^{2}$ equations for only $L^{2}-1$ components of $v_{\lambda}$ (one of them was set equal to 1). The extra equation can be seen as a constraint on the values of $\lambda$. By the Chinese remainder theorem, it is sufficient to look for solutions of (6.6) for $\lambda$ a prime power. For if $v_{\lambda_{1}}$ and $v_{\lambda_{2}}$ are modular eigenvectors modulo $\lambda_{1}$ and $\lambda_{2}$ respectively, with $\lambda_{1}$ and $\lambda_{2}$ coprime, $v=v_{\lambda_{1}} \lambda_{2}+v_{\lambda_{2}} \lambda_{1}$ is a modular eigenvector modulo $\lambda_{1} \lambda_{2}$. Determining all independent solutions modulo prime powers and systematically using the Chinese remainder theorem to obtain the smallest number of modular eigenvectors, we get a set of eigenvectors $\left\{v_{\lambda_{1}}\right\}$ where $\lambda_{i}$ is a multiple of $\lambda_{i+1}$. ((6.6) implies that $\lambda_{1}$, hence all $\lambda_{i}$, divides the determinant of $\Delta$.) The corresponding set of invariants $\left\{A_{\lambda}\right\}$ defines a complete set of toppling invariants.

As an example, consider once more the $2 \times 2$ case. Let us choose the $v_{\lambda}$ 's and the $\lambda$ 's in $Z$, and assume

$$
\begin{equation*}
A_{\lambda}(C)=z_{1}+b z_{2}+c z_{3}+d z_{4} \bmod \lambda \tag{6.7}
\end{equation*}
$$

where the configuration $C$ is represented by $\left(\begin{array}{ll}z_{1} & z_{2} \\ z_{3} & z_{4}\end{array}\right)$. Equation (6.6) implies $c=$ $4-b \bmod \lambda$ and $d=4 b-1 \bmod \lambda$, as well as

$$
\begin{equation*}
8 b-16=16 b-8=0 \bmod \lambda \tag{6.8}
\end{equation*}
$$

For $\lambda=p^{m}$ with $p$ an odd prime, (6.8) is equivalent to $b=2=2^{-1} \bmod p^{m}$, which has no solution unless $p^{m}=3$, in which case one finds one invariant,

$$
\begin{equation*}
A_{3}(C)=z_{1}+2 z_{2}+2 z_{3}+z_{4} \bmod 3 \tag{6.9}
\end{equation*}
$$

For $\lambda=2^{m}$, the two equations (6.8) are trivially satisfied for $2^{m}=8$, but cannot be satisfied for $2^{m} \geqslant 16(16 b-8=-8 \neq 0 \bmod 16)$. Thus the maximal value of $\lambda$ is 8 , and one finds two independent invariants modulo $8(b=1,-2)$

$$
\begin{align*}
& A_{8}(C)=z_{1}+z_{2}+3 z_{3}+3 z_{4} \bmod 8  \tag{6.10a}\\
& A_{8}^{\prime}(C)=z_{1}-2 z_{2}-2 z_{3}-z_{4} \bmod 8 \tag{6.10b}
\end{align*}
$$

Using the Chinese remainder theorem, we obtain two invariants, modulo 24 and 8 respectively:

$$
\begin{align*}
& A_{24}(C)=8 A_{3}(C)+3 A_{8}^{\prime}(C)=11 z_{1}+10 z_{2}+10 z_{3}+5 z_{4} \bmod 24  \tag{6.11a}\\
& A_{8}(C)=z_{1}+z_{2}+3 z_{3}+3 z_{4} \bmod 8 \tag{6.11b}
\end{align*}
$$

As expected, one recovers the invariants (3.8) constructed out of the Smith normal form of $\Delta, I_{1}=5 A_{24}$ and $I_{2}=3 A_{8}$. Clearly, this is a general fact: choosing to work on $Z$, the above algorithm allows us to compute the elementary divisors of $\Delta$ (the $\lambda_{i}$ 's), as well as the relevant lines of the matrix $A^{-1}$ needed to obtain the invariants (see (3.11)), since from (3.9) and (3.10), the rows of $A^{-1}$ are precisely modular eigenvectors modulo the elementary divisors of $\Delta$. However, the algorithm is completely general regarding the ring of integers in which we want to work.

As a conclusion to this section, we found it rather attractive to use Galois theory in order to associate toppling invariants with irreducible factors of the characteristic polynomial of $\Delta$, to the end of extracting information about $G$ from the Galois decomposition of this polynomial. Although this method can be useful, the usefulness being dependent on the spectrum of $\Delta$, it generally provides only partial information on $G$.

## 7. Some related problems

In this section, we comment briefly on two intriguing questions related to sandpile automata on finite square lattices.

Question 1. Wiesenfeld et al [15] have studied a deterministic version of the sandpile automaton considered here, in which the particle addition is always done at the central site of a $(2 L+1) \times(2 L+1)$ square lattice. They observed, and it is easy to prove, that starting from an initial arbitrary configuration, once the transient configurations are gone, the system shows a cyclic behaviour, and the period of the cycle $T_{L}$ is independent of the initial configuration. The exact dependence of $T_{L}$ on $L$ is still quite puzzling.

Wiesenfeld et al have numerically determined $T_{L}=4,16,104,544,146248,7889840$ for $L=0$ to 5 , and found that $T_{L}$ increases at a much slower rate than Det $\Delta$, $T_{L} \sim \exp 0.44 L^{2}$ against $\operatorname{Det} \Delta \sim \exp 4.67 L^{2}$. The period for $L$ up to 9 has been determined by Markosova and Markos [16]. It is easy to see that $T_{L}$ is the order of the operator $a(L+1, L+1)$ on the space $\mathcal{R}$ of recurrent configurations and so is the smallest positive integer such that

$$
\begin{equation*}
a(L+1, L+1)^{T_{L}}=I . \tag{7.1}
\end{equation*}
$$

If $a(L+1, L+1)^{T_{L}}=I$ on $\mathcal{R}$, we obtain that all the toppling invariants (3.3) must be equal to zero modulo 1 at the configuration which is zero everywhere except at the central site, where it is $T_{L}$. Multiplying the invariants by Det $\Delta$ and defining the integer matrix $E=(\operatorname{Det} \Delta) \Delta^{-1}$, we obtain that $T_{L}$ is the smallest positive integer satisfying

$$
\begin{equation*}
E_{i,(L+1, L+1)} T_{L}=0 \bmod (\operatorname{Det}, \Delta) \quad \text { for all } i . \tag{7.2}
\end{equation*}
$$

Setting $M=\operatorname{gcd}\left\{E_{i,(L+1, L+1)}: 1 \leqslant i \leqslant 2 L+1\right\}$, it is easy to see that

$$
\begin{equation*}
T_{L}=(\operatorname{Det} \Delta) / M \tag{7.3}
\end{equation*}
$$

Since $T_{L}$ is independent of the initial configuration, we may choose a configuration which has the symmetry of the square. This symmetry is preserved under symmetrical topplings, and under addition of sand at the central site, hence we can restrict ourselves to $\mathcal{R}_{\text {sym }}$, the subspace of recurrent configurations having the symmetry of a square. On
a $(2 L+1) \times(2 L+1)$ lattice, a symmetrical configuration is completely specified by the heights in an octant having $\frac{1}{2}(L+1)(L+2)$ sites. Instead of studying the symmetrical configurations on a square lattice, we can just as well study all recurrent configurations on an octant $O_{L}$ with a new toppling matrix $\Delta_{\text {sym }}$ given by

$$
\begin{equation*}
\left(\Delta_{\mathrm{sym}}\right)_{i j}=\sum_{j^{\prime}} \Delta_{i, j^{\prime}} \quad . \quad \text { for } i, j \in O_{L} \tag{7.4}
\end{equation*}
$$

where the sum over $j^{\prime}$ is over sites which are related to $j$ by the symmetries of the square (dihedral group of order 8 ). Even though it does not satisfy $\sum_{i}\left(\Delta_{\text {sym }}\right)_{i j} \geqslant 0$ for all $i$, it can be shown that the general results of [6] hold for this case, so that for example

$$
\begin{equation*}
\left|\mathcal{R}_{\mathrm{sym}}\right|=\operatorname{Det} \Delta_{\mathrm{sym}}=\left|G_{\mathrm{sym}}\right| \tag{7.5}
\end{equation*}
$$

where $G_{\text {sym }}$ is the Abelian group generated by the operators $a_{i}, i \in O_{L}$, subjected to the relations (2.3) with $\Delta$ replaced by $\Delta_{\text {sym }}$. A simple calculation shows that

$$
\begin{equation*}
\text { Det } \Delta_{\mathrm{sym}}=\prod_{m \leqslant n} \lambda_{m . n} \tag{7.6}
\end{equation*}
$$

where the product over $m$ and $n$ is over odd values of $m \leqslant n$ between 1 and ( $2 L+1$ ), and with $\lambda_{m, n}$ given in (5.1). For large $L,\left|G_{\text {sym }}\right|$ varies as $|G|^{1 / 8} \sim \exp 0.58 L^{2}$, so that $T_{L}$ increases at a substantially lower rate than $\left|G_{\text {sym }}\right|$. As in section 3 , the structure of $G_{\text {sym }}$ can be determined by computing the elementary divisors of $\Delta_{\text {sym }}$. We find for the first four values of $L$

$$
\begin{align*}
& L=1: G_{\text {sym }}=Z_{16} \times Z_{2}  \tag{7.7a}\\
& L=2: G_{\text {sym }}=Z_{104} \times Z_{4} \times Z_{2}  \tag{7.7b}\\
& L=3: G_{\text {sym }}=Z_{544} \times Z_{32} \times Z_{2} \times Z_{2}  \tag{7.7c}\\
& L=4: G_{\text {sym }}=Z_{146248} \times Z_{8} \times Z_{4} \times Z_{2} \times Z_{2} . \tag{7.7d}
\end{align*}
$$

We see that the order of the largest cyclic group in $G_{\text {sym }}$ is apparently nothing but the period $T_{L}$. In terms of the invariants (3.11) constructed from the Smith normal form of $\Delta_{\text {sym }}$, this would mean that the element of the first row of $A^{-1}$ corresponding to the central site is coprime with the largest elementary divisor $d_{1}^{\text {sym }}$, implying $T_{L}=d_{1}^{\text {sym }}$. What is certainly true is that $T_{L}$ must divide $d_{1}^{\text {sym }}$ since $a(L+1, L+1)$ generate a cyclic subgroup of $G_{\text {sym }}$, the largest one being of order $d_{1}^{\text {sym }}$. Although there are plausibility arguments that $T_{L}$ is actually equal to $d_{1}^{\text {sym }}$, the proof (or disproof) of this simple fact is still lacking.

The question of the rank of $G_{\text {sym }}$ can also be addressed. However, here an argument based on the eigenvectors of eigenvalue 4 , similar to the one we used in section 4 , does not work since the degeneracy of the eigenvalue 4 is only $\left[\frac{1}{2} L\right]+1$, with $[x]$ the integral part of $x$. (It nevertheless implies that so many cyclic factors of $G_{\text {sym }}$ have order multiple of 4.) But on the other hand, it is easy to see that there are exactly $L+1$ independent modular eigenvectors modulo 2 (of $\Delta_{\text {sym }}$ ), which shows that the rank of $G_{\text {sym }}$ is $L+1$. ( $L+1$ is also an upper bound since all operators $a(x, y)$ for $(x, y) \in O_{L}$ can be expressed in terms of $a(1, y)$ for $y=1,2, \ldots, L+1$.) The independent modular eigenvectors have only one non-zero component, at any of the $L+1$ positions on the principal diagonal of
$O_{L}: v_{\lambda=2}^{(k)}(x, y)=\delta_{x, k} \delta_{y, k}$ for $1 \leqslant k \leqslant L+1$. This proves that each cyclic subgroup of $G_{\text {sym }}$ has order multiple of 2.

The fact that $T_{L}$ must be a divisor of $\left|G_{\text {sym }}\right|$ already gives us a fairly good upper-bound estimate of it. It is instructive to look at the values of $\left[G_{\text {sym }} \mid / T_{L}\right.$ for small $L$. From the results of Markosova et al [16], we see that for $L=0$ up to 8 , these values are $1,2,2^{3}, 2^{7}$, $2^{7}, 2^{11}, 2^{9}, 2^{22}, 2^{15} \times 2701$. Note that they show a rather irregular, non-monotonic behaviour with $L$. They have a large power of 2 as divisor, but sometimes have other factors as well.

The fluctuations in $T_{L}$ appear to be related to the appearance of degeneracies in the eigenvalue spectrum of $\Delta_{\text {sym }}$. For example, the fact that the eigenvalue 4 is $\left(\left[\frac{1}{2} L\right]+1\right)$-fold degenerate and the existence of $L+1$ modular eigenvectors modulo 2 imply that $T_{L}$ is a divisior of $2^{-(L+[L / 2])}\left|G_{\text {sym }}\right|$. Another accidental degeneracy in the spectrum of $G_{\text {sym }}$ occurs when $2 L+2$ is a multiple of 6 . Writing $2 L+2=6 \ell$, we can easily verify that

$$
\begin{align*}
& \lambda_{2 \ell-m, 2 \ell+m}=\lambda_{3 \ell, m}  \tag{7.8}\\
& \lambda_{4 \ell-m, 4 \ell+m}=\lambda_{3 \ell, 6 \ell-m} .
\end{align*}
$$

In the Galois factorization of $\left|G_{\text {sym }}\right|$ this implies that there are some factors which are repeated. For example, for $L=8$, the factor 2701 occurs twice in $\left|G_{\text {sym }}\right|$. While for a given value of $L$, the degeneracies, and the group structure of $G_{\text {sym }}$ can be explicitly determined, it seems difficult at this stage to say much more about this 'irregular' variation of $T_{L}$ with $L$ for general $L$.

Question 2. This concerns the structure of the identity configuration. Even the $L \times L$ square case shows non-trivial features. The identity configuration $\dot{C}_{i d}$ is the unique recurrent configuration which is equivalent under toppling to the configuration with all heights zero. We have already noted that all toppling invariants $I_{i}\left(C_{\mathrm{id}}\right)$ must be zero. It also implies that $C_{i d}$ is the unique recurrent configuration such that $T_{i} \equiv \sum_{j}\left(\Delta^{-1}\right)_{i j} z_{j}\left(C_{i d}\right)$ are integers for all $i$ (by using the invariants (3.3)). This integer vector has the interesting interpretation that $T_{i}$ is the number of topplings occurring at $i$ when $C_{i d}$ is added to itself.

The identity configuration shows complicated fractal structures. Some colour-coded pictures of these may be found in [21]. We mention here two remarkable properties of the identity $\dot{\dagger}$. First, in the central area of a $2 L \times 2 L$ lattice, there is a whole square where all sites have maximal height, equal to 3 . The linear size of this central square was numerically measured for $L$ up to 100 , and goes like $4 L / 5$, so that the central square covers approximately $4 / 25$ of the whole lattice.

The second property, even more remarkable, is that the identity configuration on the $(2 L+1) \times(2 L+1)$ lattice appears to be related in a very simple way to that on the $2 L \times 2 L$ lattice. Indeed, if we divide the configuration $C_{\mathrm{id}}^{(2 L)}$ into four equal squares and pull them apart by one lattice spacing so as to leave a cross in the middle, we get $C_{\mathrm{id}}^{(2 L+1)}$ provided the heights on the cross are properly assigned. In obvious notation, if

$$
C_{\mathrm{id}}^{(2 L)}=\left(\begin{array}{ll}
B_{1} & B_{2}  \tag{7.9}\\
B_{3} & B_{4}
\end{array}\right) \quad\left(B_{i} \text { are } L \times L \text { blocks }\right)
$$

[^1]where the four blocks $B_{i}$ are related by the symmetry transformations of the square (since $C_{i d}$ has that symmetry), then
\[

C_{\mathrm{id}}^{(2 L+1)}=\left($$
\begin{array}{ccc}
B_{1} & R_{1} & B_{2}  \tag{7.10}\\
R_{2} & z_{\text {mid }} & R_{3} \\
B_{3} & R_{4} & B_{4}
\end{array}
$$\right) \quad\left(R_{i} are length L rows or columns\right) .
\]

In addition the height at the centre $z_{\text {mid }}$ is always 0 , and the branches of the cross given by the $R_{i}$ (also related by symmetry transformations) have a simple structure. The first instance of this phenomenon occurs when going from the $2 \times 2$ to the $3 \times 3$ lattices:

$$
C_{\mathrm{id}}^{(2)}=\left(\begin{array}{ll}
2 & 2  \tag{7.11}\\
2 & 2
\end{array}\right) \rightarrow C_{\mathrm{id}}^{(3)}=\left(\begin{array}{lll}
2 & 1 & 2 \\
1 & 0 & 1 \\
2 & 1 & 2
\end{array}\right)
$$

More identity configurations are listed in appendix C , for which this property may be checked (as well as the first one mentioned above).

In addition, the array $T_{i}$ defined in a previous paragraph also has the 'cross' property (7.9-10). For instance,
$T^{(6)}=\left(\begin{array}{llllll}2 & 3 & 4 & 4 & 3 & 2 \\ 3 & 5 & 6 & 6 & 5 & 3 \\ 4 & 6 & 7 & 7 & 6 & 4 \\ 4 & 6 & 7 & 7 & 6 & 4 \\ 3 & 5 & 6 & 6 & 5 & 3 \\ 2 & 3 & 4 & 4 & 3 & 2\end{array}\right) \rightarrow T^{(7)}=\left(\begin{array}{lllllll}2 & 3 & 4 & 4 & 4 & 3 & 2 \\ 3 & 5 & 6 & 6 & 6 & 5 & 3 \\ 4 & 6 & 7 & 7 & 7 & 6 & 4 \\ 4 & 6 & 7 & 7 & 7 & 6 & 4 \\ 4 & 6 & 7 & 7 & 7 & 6 & 4 \\ 3 & 5 & 6 & 6 & 6 & 5 & 3 \\ 2 & 3 & 4 & 4 & 4 & 3 & 2\end{array}\right)$.
Because the $C_{\mathrm{id}}$ and $T$ arrays are related by $C_{\mathrm{id}}=\Delta T$, the fact that both have the 'cross' property fixes their values on the branches of the cross (for odd lattices). For instance, the rows in $T$ corresponding to the $R_{i}$ (see (7.10)) must be equal to the rows bordering them, as shown in (7.12) for $2 L+1=7$. The first property mentioned for $C_{i d}$, namely the existence of a big central square with constant values, does not hold for $T$. There is, in fact, in $T$ a central square of constant values, but its linear size cannot exceed 3 if $C_{\text {id }}$ is to be recurrent at all (see (7.12)). As a final remark, we observed that the total number of topplings occurring when the identity is added to itself, i.e. $\sum_{i} T_{i}^{(L)}$, equal to 160 and 235 for $L=6$ and 7 , grows like a power of $L$, with an exponent close to $4(\sim 3.9)$.

Similar geometric structures are found when the ASM are allowed to relax from special unstable states, say all heights equal to 4 [14]. These fractal structures are not well understood yet.

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## Appendix A

In this appendix, we determine the structure of the group $G$ for the ASM defined on the $L \times 2$ lattice. In this case, the particle addition operators $a(x, y)$ satisfy the relations

$$
\begin{equation*}
a(x+1,1)=a^{4}(x, 1) a^{-1}(x-1,1) a^{-1}(x, 2) \tag{A1}
\end{equation*}
$$

and

$$
\begin{equation*}
a(x+1,2)=a^{4}(x, 2) a^{-1}(x-1,2) a^{-1}(x, 1) \tag{A2}
\end{equation*}
$$

Here $1 \leqslant x \leqslant L$, and by convention we assume that

$$
\begin{equation*}
a(0, y)=a(L+1, y)=I \quad \text { for } y=1,2 . \tag{A3}
\end{equation*}
$$

Using the recursion relations (A1) and (A2) we can express $a(x, 1)$ and $a(x, 2)$ as product of powers of $a \equiv a(1,1)$ and $b \equiv a(1,2)$. Explicitly we have

$$
\begin{align*}
& a(x, 1)=a^{\mu_{x}} b^{-v_{x}}  \tag{A4}\\
& a(x, 2)=a^{-v_{x}} b^{\mu_{x}} \tag{A5}
\end{align*}
$$

where $\mu_{x}$ and $\nu_{x}$ are positive integers satisfying the following recurrence relations

$$
\begin{align*}
& \mu_{x+1}=4 \mu_{x}+v_{x}-\mu_{x-1}  \tag{A6}\\
& v_{x+1}=4 v_{x}+\mu_{x}-v_{x-1} \tag{A7}
\end{align*}
$$

with the initial conditions $\mu_{0}=\nu_{0}=v_{1}=0$, and $\mu_{1}=1$. This can easily be solved by defining the new functions

$$
\begin{equation*}
m_{x}=\mu_{x}-v_{x} \quad n_{x}=\mu_{x}+v_{x} \tag{A8}
\end{equation*}
$$

They satisfy uncoupled binary recurrence relations

$$
\begin{align*}
& m_{x+1}=3 m_{x}-m_{x-1}  \tag{A9}\\
& n_{x+1}=5 n_{x}-n_{x-1} \tag{A10}
\end{align*}
$$

with $m_{0}=n_{0}=0, m_{1}=n_{1}=1$. The solutions read

$$
\begin{align*}
& m_{x}=\frac{1}{\sqrt{5}}\left[\left(\frac{3+\sqrt{5}}{2}\right)^{x}-\left(\frac{3-\sqrt{5}}{2}\right)^{x}\right]  \tag{Al1}\\
& n_{x}=\frac{1}{\sqrt{21}}\left[\left(\frac{5+\sqrt{21}}{2}\right)^{x}-\left(\frac{5-\sqrt{21}}{2}\right)^{x}\right] \tag{A12}
\end{align*}
$$

Table A1.

| $x$ | 0 | 1 | 2 | 3 | 4 | 5 | 6 | 7 | 8 |
| :--- | :--- | :--- | :--- | ---: | ---: | ---: | ---: | ---: | ---: |
| $m_{x}$ | 0 | 1 | 3 | 8 | 21 | 55 | 144 | 377 | 987 |
| $n_{x}$ | 0 | 1 | 5 | 24 | 115 | 551 | 2640 | 12649 | 60605 |
| $\mu_{x}$ | 0 | 1 | 4 | 16 | 68 | 303 | 1392 | 6513 | 30796 |
| $\nu_{x}$ | 0 | 0 | 1 | 8 | 47 | 248 | 1248 | 6136 | 29809 |

The first values of $m_{x}, n_{x}, \mu_{x}, \nu_{x}$ are given in table A1.
One may note that $m_{x}$ are precisely the even terms in the standard Fibonacci sequence.
To determine the structure of $G$ in terms of the generators $a$ and $b$, we note from (A3)-(A5) that the only relations among them are

$$
\begin{align*}
& a^{\mu_{L+1}} b^{-\nu_{L+1}}=I  \tag{A13}\\
& a^{-\nu_{L+1}} b^{\mu_{L+1}}=I . \tag{A14}
\end{align*}
$$

According to the discussions in sections 3 and 4, the structure of $G$ is given in terms of the elementary divisors $d_{1}, d_{2}$ of the matrix $\left(\begin{array}{cc}\mu_{L+1} & -v_{L+1} \\ -v_{L+1} & \mu_{L+1}\end{array}\right)$. We find that

$$
\begin{equation*}
G_{L \times 2}=Z_{d_{1}} \times Z_{d_{2}} \quad \text { with } d_{1}=\left(\mu_{L+1}^{2}-v_{L+1}^{2}\right) / d_{2} \text { and } d_{2}=\operatorname{gcd}\left(\mu_{L+1}, v_{L+1}\right) \tag{A15}
\end{equation*}
$$

As a function of $L$, the number $d_{2}$ has an irregular behaviour, with very large sudden jumps. The only result we could firmly establish follows from the double inequality (4.12): $d_{2}>1$ if $L+1$ is divisible by 3 , so in that case, $G$ is not cyclic.

Rather surprisingly, the elementary divisors can equally be expressed in terms of $m_{L+1}$ and $n_{L+1}$ by the neater formula

$$
\begin{equation*}
d_{1}=m_{L+1} n_{L+1} / d_{2} \quad d_{2}=\operatorname{gcd}\left(m_{L+1}, n_{L+1}\right) \tag{A16}
\end{equation*}
$$

From the relation (A8), equations (A15) and (A16) are clearly equivalent when $m_{x}$ and $n_{x}$ are both odd, which is the case if $x$ is not divisible by 3 . That they are equivalent for $x$ a multiple of 3 lies in the delicate fact that the largest power of 2 that divides $m_{x}$ is the same as that dividing $n_{x}$. To see this, we consider the subsequences $M(r)=m_{3 r}$ and $N(r)=n_{3 r}$. We find from (A11)-(A12) that they satisfy the recursions

$$
\begin{align*}
& M(r+1)=18 M(r)-M(r-1)  \tag{A17}\\
& N(r+1)=110 N(r)-N(r-1) \tag{A18}
\end{align*}
$$

subjected to the initial conditions $M(0)=N(0)=0, M(1)=8$, and $N(1)=24$. Now the desired result that for each $r, M(r)$ and $N(r)$ have the same 2-potency (i.e. the largest powers of 2 dividing them are the same) follows from a more general result by the last author [22], which gives a formula for the 2-potency for any 'binary recursive sequence' $\theta(r)$ defined by $\theta(0)=0, \theta(1)=\theta$, and $\theta(r+1)=2 R \theta(r)+S \theta(r-1)$, where $\theta$ is an arbitrary integer, and $R, S$ are odd integers.

## Appendix $B$

In this appendix, we give the explicit formula giving the number $F(L)$ of irreducible factors (over $Q$ ) of the polynomial $P_{L}(x)$ for arbitrary $L$. In section $5, P_{L}(x)$ was the characteristic equation of the discrete Laplacian on a $L \times L$ lattice with open boundary conditions and was defined by ( $N=2(L+1)$ )

$$
\begin{equation*}
P_{L}(x)=\prod_{m, n=1}^{L}\left(x-4+2 \cos \frac{2 \pi m}{N}+2 \cos \frac{2 \pi n}{N}\right) . \tag{B1}
\end{equation*}
$$

The roots $\lambda_{m, n}=4-2 \cos (2 \pi m / N)-2 \cos (2 \pi n / N)$ of $P_{L}$ belong to the extension $Q(\cos (2 \pi / N))$ and are permuted by its Galois group $\mathrm{Gal}_{L}$, which acts on them by (see section 5)

$$
\begin{equation*}
\sigma_{s}\left(\lambda_{m, n}\right) \equiv \lambda_{s m, s n} \quad \text { for } s \in Z_{N}^{*} /\{ \pm 1\} \tag{B2}
\end{equation*}
$$

As discussed in section 5, finding the irreducible factors of $P_{L}$ amounts to split the set of roots $\left\{\lambda_{m, n}\right\}$ into orbits under $\mathrm{Gal}_{L}$, so that

$$
\begin{equation*}
F(L)=\left|\left\{\lambda_{m, n}: 1 \leqslant m, n \leqslant L\right\} / \operatorname{Gal}_{L}\right| . \tag{B3}
\end{equation*}
$$

By a classical theorem (see for instance [23]), the number of orbits of a set $X$ under the action of a group $G$ is equal to the average number of fixed points of $G$ in $X$, that is

$$
\begin{equation*}
|X / G|=\frac{1}{|G|} \sum_{g \in G}\left|\left\{x \in X^{-}: g x=x\right\}\right| . \tag{B4}
\end{equation*}
$$

This allows us to recast (B3) as

$$
\begin{equation*}
F(L)=\frac{2}{\varphi(N)} \sum_{\left.s \in Z_{N}^{*} / \neq 1\right\}}\left\{\left\{1 \leqslant m, n \leqslant L: \lambda_{s m, s n}=\lambda_{m, n}\right\}\right\} . \tag{B5}
\end{equation*}
$$

From the expression of $\lambda_{m, n}$, we have the obvious relations $\lambda_{m, n}=\lambda_{ \pm m, \pm n}=\lambda_{ \pm n, \pm m}$ where the signs are uncorrelated. The only subtle point of the analysis lies in the possible existence of other relations. For example, if $L=4$ (i.e. $N=10$ ), $s$ takes the two values 1 and 3, we have that $\sigma_{3}\left(\lambda_{1,4}\right) \equiv \lambda_{3,2}$ is actually equal to $\lambda_{1,4}=4$ even though $(3,2) \neq( \pm 1, \pm 4)$ nor $( \pm 4, \pm 1) \bmod 10$. This situation is, however, not generic, and can be proved to only happen for the pairs ( $m, n$ ) with $m+n=L+1$ (related to the eigenvalue 4 of the Laplacian). Thus if we leave them aside, we obtain the result that $\lambda_{s m, s n}=\lambda_{m, n}$ if and only if $(s m, s n)=( \pm m, \pm n)$ or $( \pm n, \pm m) \bmod N$. Each of the $L$ eigenvalues $\lambda_{m, L+1-m}=4$ clearly forms an orbit on its own under the Galois group and we obtain

$$
\begin{gather*}
F(L)=L+\frac{2}{\varphi(N)} \sum_{n \in Z_{N}^{*} /\{ \pm 1]}\{\{1 \leqslant m, n \leqslant L \text { and } m+n \neq L+1:(s m, s n) \\
= \pm(m, \pm n) \text { or } \pm(n, \pm m) \bmod N\}\} . \tag{B6}
\end{gather*}
$$

The actual calculation of (B6) being straightforward but somewhat lengthy, we only quote the final result. Let $N=2(L+1)$ have the prime factorization given by

$$
\begin{equation*}
N=\prod_{p} p^{a_{r}} . \tag{B7}
\end{equation*}
$$

Then the number of irreducible polynomial factors of $P_{L}(x)$ equals

$$
\begin{align*}
F(L)=\frac{N}{2}+6 & +\frac{1}{2} \prod_{p} \frac{p^{a_{p}+1}+p^{a_{p}}-2}{p-1}+2\left(2 a_{2}-1\right) \prod_{p \neq 2}\left(2 a_{p}+1\right) \\
& +\prod_{p=1 \mathrm{mod} 4}\left(2 a_{p}+1\right)-\left(7 a_{2}+5\right) \prod_{p \neq 2}\left(a_{p}+1\right)+\left[\prod_{p \neq 2}\left(2 a_{p}+1\right)-1\right] \delta\left(a_{2}=1\right) \tag{B8}
\end{align*}
$$

It can be shown that the average growth of $F(L)$ is linear in $L$. Some particular values are given in table 1 , section 5 .

## Appendix $\mathbf{C}^{*}$

We list here the identity configurations $C_{\mathrm{id}}^{(L)}$, discussed in sections 3 and 7 , for $L=4$ and 5 and $L=10$ and 11 .

$$
\begin{aligned}
C_{\mathrm{id}}^{(4)} & =\left(\begin{array}{llll}
2 & 3 & 3 & 2 \\
3 & 2 & 2 & 3 \\
3 & 2 & 2 & 3 \\
2 & 3 & 3 & 2
\end{array}\right) \\
C_{\mathrm{id}}^{(5)} & =\left(\begin{array}{lllll}
2 & 3 & 2 & 3 & 2 \\
3 & 2 & 1 & 2 & 3 \\
2 & 1 & 0 & 1 & 2 \\
3 & 2 & 1 & 2 & 3 \\
2 & 3 & 2 & 3 & 2
\end{array}\right) \\
C_{\mathrm{id}}^{(10)} & =\left(\begin{array}{llllllllll}
2 & 3 & 3 & 0 & 3 & 3 & 0 & 3 & 3 & 2 \\
3 & 2 & 2 & 1 & 2 & 2 & 1 & 2 & 2 & 3 \\
3 & 2 & 2 & 3 & 3 & 3 & 3 & 2 & 2 & 3 \\
0 & 1 & 3 & 2 & 2 & 2 & 2 & 3 & 1 & 0 \\
3 & 2 & 3 & 2 & 2 & 2 & 2 & 3 & 2 & 3 \\
3 & 2 & 3 & 2 & 2 & 2 & 2 & 3 & 2 & 3 \\
0 & 1 & 3 & 2 & 2 & 2 & 2 & 3 & 1 & 0 \\
3 & 2 & 2 & 3 & 3 & 3 & 3 & 2 & 2 & 3 \\
3 & 2 & 2 & 1 & 2 & 2 & 1 & 2 & 2 & 3 \\
2 & 3 & 3 & 0 & 3 & 3 & 0 & 3 & 3 & 2
\end{array}\right)
\end{aligned}
$$

$$
C_{\text {id }}^{(11)}=\left(\begin{array}{ccccccccccc}
2 & 3 & 3 & 0 & 3 & 2 & 3 & 0 & 3 & 3 & 2 \\
3 & 2 & 2 & 1 & 2 & 1 & 2 & 1 & 2 & 2 & 3 \\
3 & 2 & 2 & 3 & 3 & 2 & 3 & 3 & 2 & 2 & 3 \\
0 & 1 & 3 & 2 & 2 & 1 & 2 & 2 & 3 & 1 & 0 \\
3 & 2 & 3 & 2 & 2 & 1 & 2 & 2 & 3 & 2 & 3 \\
2 & 1 & 2 & 1 & 1 & 0 & 1 & 1 & 2 & 1 & 2 \\
3 & 2 & 3 & 2 & 2 & 1 & 2 & 2 & 3 & 2 & 3 \\
0 & 1 & 3 & 2 & 2 & 1 & 2 & 2 & 3 & 1 & 0 \\
3 & 2 & 2 & 3 & 3 & 2 & 3 & 3 & 2 & 2 & 3 \\
3 & 2 & 2 & 1 & 2 & 1 & 2 & 1 & 2 & 2 & 3 \\
2 & 3 & 3 & 0 & 3 & 2 & 3 & 0 & 3 & 3 & 2
\end{array}\right) .
$$

## Appendix D

In [17], Lee and Tzeng have considered the following one-dimensional deterministic sandpile model. The $L$ sites of a linear chain, labelled 0 to $L-1$, are assigned a height variable $z_{i}$ taking positive integer values. A site $j$ becomes unstable if the difference $z_{j}-z_{j+1}$ exceeds some integer $N$, considered as a parameter. In this case, the site $j$ loses $N$ grains of sand which fall on its $N$ right nearest neighbours. They take open boundary conditions $z_{i}=0$ for all $i \geqslant L$, so that, under the toppling at site $j$, the sand is redistributed according to $z_{i} \rightarrow z_{i}-\Delta_{i j}$, where the toppling matrix is the lower triangular matrix

$$
\Delta_{i j}= \begin{cases}N & \text { if } i=j  \tag{D1}\\ -1 & \text { if } 1 \leqslant i-j \leqslant N .\end{cases}
$$

The deterministic dynamics is defined by dropping, at each time step, one grain of sand on the leftmost site $i=0$ and by letting the system relax. The number of recurrent configurations is equal to $N^{L}$, also equal to the determinant of $\Delta$. Note that this model is a non-Abelian sandpile model as the toppling condition is not local.

It was shown in [17] that for $N=2$ and 3 , the recurrent configurations can be labelled by the values of a toppling invariant, an integer function taken modulo $N^{L}$, (see also [24]). Lee and Tzeng wrote down the invariant when $N=2$ for all values of $L$, as well as when $N=3$ for some particular values of $L$. We show here that our general method allows to prove that one invariant is sufficient to label all recurrent configurations, for all $N$, and we will compute this invariant for $N=2,3$, all values of $L$.

That only one invariant gives a complete labelling of the recurrent configurations amounts to prove that the elementary divisors of $\Delta$ are $d_{1}=N^{L}$ and $d_{i}=1$ for $i \geqslant 2$. Equivalently, it amounts to showing that among all the invariants (3.3) constructed from $\Delta^{-1}$, only one is independent. The matrix $\Delta^{-1}$ is easily obtained and reads as

$$
\begin{equation*}
\Delta_{i j}^{-1}=\sum_{k=0}^{L-1} c_{k} \delta_{i-j, k} \tag{D2}
\end{equation*}
$$

where the entries $c_{k}$ satisfy the recurrence relations

$$
\begin{align*}
& N c_{k}=c_{k-1}+c_{k-2}+\cdots+c_{k-N} \\
& c_{0}=\frac{1}{N} \quad c_{-k}=0 \quad \text { for all } k \geqslant 1 . \tag{D3}
\end{align*}
$$

We thus obtain the $L$ invariants

$$
\begin{equation*}
Q_{i}(C)=\sum_{k=0}^{i} c_{i-k} z_{k} \bmod 1 \quad i=0,1, \ldots L-1 \tag{D4}
\end{equation*}
$$

The first invariants read as

$$
\begin{align*}
& Q_{0}=\frac{1}{N} z_{0} \bmod 1 \quad Q_{1}=\frac{1}{N^{2}} z_{0}+\frac{1}{N} z_{1} \bmod 1  \tag{D5}\\
& Q_{2}=\frac{N+1}{N^{3}} z_{0}+\frac{1}{N^{2}} z_{1}+\frac{1}{N} z_{2} \bmod 1
\end{align*}
$$

so that we have

$$
\begin{align*}
& N Q_{1}=Q_{0} \bmod 1  \tag{D6a}\\
& N Q_{1}=(N+1) Q_{0} \bmod 1 \quad N Q_{2}=(N+1) Q_{1} \bmod 1 \tag{D6b}
\end{align*}
$$

showing that $Q_{0}$ can be expressed in terms of $Q_{1}$, itself expressible in terms of $Q_{2}$. It is not difficult to show recursively that this pattern continues to hold, namely there exists, for all $i$ between 1 and $L-1$, an integer $X_{i}$ coprime with $N$ such that

$$
\begin{equation*}
N Q_{j}=X_{i} Q_{j-1} \bmod 1 \quad \text { for all } 1 \leqslant j \leqslant i \tag{D7}
\end{equation*}
$$

(From (D6) we have $X_{1}=1$ and $X_{2}=N+1$.) Equation (D7) implies that all invariants $Q_{0}, Q_{1}, \ldots, Q_{L-2}$ can be expressed in terms of $Q_{L-1}$ which is thus the only independent one and which provides a complete labelling of the set of recurrent configurations.

To compute the invariant $Q_{L-1}$, we note that the coefficients $c_{k}$ satisfy in fact a recurrence relation of order $N-1$ :
$N c_{k}+(N-1) c_{k-1}+\cdots+c_{k+1-N}=N c_{N-1}+(N-1) c_{N-2}+\cdots+c_{0}=1$.
For $N=2$ and 3, the solutions to (D8) read ( $k \geqslant 0$ )

$$
\begin{align*}
& N=2: c_{k}=\frac{1}{3}\left[1-\left(\frac{-1}{2}\right)^{k+1}\right]  \tag{D9a}\\
& N=3: c_{k}=\frac{1}{6}+\frac{1}{12}\left[\left(\frac{-1-\sqrt{-2}}{3}\right)^{k}+\left(\frac{-1+\sqrt{-2}}{3}\right)^{k}\right] . \tag{D9b}
\end{align*}
$$

Using these values in $Q_{L-1}$ of (D4) yields the desired result. For higher values of $N$, the coefficients $c_{k}$ are given by the generating function

$$
\begin{equation*}
f_{N}(t)=\sum_{k=0}^{\infty} c_{k}^{(N)} t^{k}=(1-t) \cdot\left[t^{N+1}-(N+1) t+N\right]^{-1} \tag{D10}
\end{equation*}
$$

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