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# Pre-Expansivity in Cellular Automata 

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#### Abstract

We introduce the property of pre-expansivity for cellular automata (CA): it is the property of being expansive on asymptotic pairs of configurations (i.e. configurations that differ in only finitely many positions). Pre-expansivity therefore lies between expansivity and pre-injectivity, two important notions of CA theory.

We show that there exist one-dimensional pre-expansive CAs which are not (positively) expansive and they can be chosen reversible. We show however that no bi-dimensional CA which is linear over an Abelian group can be preexpansive. We also consider the finer notion of k-expansivity (expansivity over pairs of configurations with exactly k differences) and show examples of linear CA in dimension 2 and on the free group that are $k$-expansive depending on the value of $k$, whereas no (positively) expansive CA exists in this setting.


## 1. Introduction

The model of cellular automata is at the crossroads of several domains and is often the source of surprisingly complex objects in several senses (computationally, dynamically, etc).

From the dynamical systems and symbolic dynamics points of view, the theory of cellular automata is very rich [1, 2, [3, 4] and tells us, on the one hand, that CA are natural examples of chaotic systems that can perfectly fit the standard notions developed in a general context, and, on the other hand, that they have special properties allowing and justifying the development of a refined and dedicated theory. For instance, the structure of the space of configurations allows to define the notion of an asymptotic pair configurations: two configurations that differ only on finitely many positions of the lattice. The Garden of Eden theorem, which has a long history [1, 5, 6, [7, 4] and is emblematic of this CA specific theoretical
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development, then says that surjectivity is equivalent to pre-injectivity (injectivity on asymptotic pairs) if and only if the lattice is given by an amenable group.

Two important lines of questioning have been particularly developed and provide some of the major open problems of the field [8]:

- surjective CA and their dynamics;
- how does CA theory change when changing the lattice. giving both a rich theory in the one-dimensional case $\mathbf{9}, \mathbf{2}, \mathbf{1 0}$ and a general inexistence result in essentially any other setting [11, $\mathbf{1 2}$. Even in the onedimensional case where positive expansivity is equivalent to being conjugated to a one-sided subshift of finite type 13, it is interesting to note that outside the linear expansive on asymptotic pairs. Our motivation is to better understand surjective CA and expansive-like dynamics, in particular in higher-dimensional case or lattices where the classical notion of (positive)expansivity cannot be satisfied by any CA 11, 12 .

The paper is organized as follows. In section 2 we give the main definitions and results we need to work on cellular automata on groups. In section 3, we introduce pre-expansivity and $k$-expansivity, and we give some preliminary results which do not depend on the group defining the space. We also consider the particular case of linear cellular automata. In section [4, we restrict to the group $\mathbb{Z}$ and give
${ }_{40}$ examples of cellular automata which are pre-expansive but not positively expansive. In section 5, we consider the free group and show that $k$-expansivity is possible for infinitely many values of $k$ although positive expansivity is impossible. Finally, in section 6, we restrict to the group $\mathbb{Z}^{2}$ and we study some $k$-expansive examples for particular values of $k$ but also show that there is no pre-expansive cellular

## 2. Definitions, notations and classical results

We will work on cellular automata defined over a finitely generated group $\mathbb{G}$. We will consider abelian and non abelian groups, but since most of our examples are given for abelian groups, we will prefer the additive notation for $\mathbb{G}$.

Fixing a generator set $G$, that is closed under inversion, a norm can be defined in $\mathbb{G}$ : given $z \in \mathbb{G},\|z\|$ is the length of the shortest sequence $g_{1} g_{2} \ldots g_{n}$ of elements in $G$ such that $z=g_{1}+g_{2}+\ldots+g_{n}$. This norm induces a metric in $\mathbb{G}$ naturally, and given a non-negative integer $r$ we can define also the ball of radius $r$ and center
$z$ as the set $B_{r}(z)=\{x \in \mathbb{G} \mid\|-x+z\| \leq r\}$. Given a point $z \in \mathbb{G}$ and two sets $X, Y \subseteq \mathbb{G}$, we accept the following notation.

$$
z+S=\{z+x \mid x \in S\}, \quad \text { and } \quad X+Y=\{x+y \mid x \in X, y \in Y\}
$$

theorem [1, 4, every cellular automaton $F$ is characterized by a local function $f: Q^{V} \rightarrow Q$, where $V \subset \mathbb{G}$ is finite and called neighborhood of $F$, as follows:

$$
\forall c \in Q^{\mathbb{G}}, \forall z \in \mathbb{G}, F(c)(z)=f\left(\sigma_{-z}\left(\left.c\right|_{z+V}\right)\right)
$$

Every function defined in this way is a cellular automaton.
Basic properties of $F$ as surjectivity and injectivity have been considered. The weaker notion of pre-injectivity says that, for every pair of different asymptotics configurations $c$ and $d$, their image by $F$ are different:

$$
c \stackrel{\infty}{=} d \text { and } c \neq d \Rightarrow F(c) \neq F(d) .
$$

The so-called Garden-of-Eden theorem establishes that surjectivity is equivalent to pre-injectivity, which in particular implies that injective CAs are also bijective (equivalently reversible by Curtis-Hedlund theorem). It was first proved in particular cases [1, 5, 6] and later it was shown that it holds exactly when the group $\mathbb{G}$ is amenable, i.e. when it admits a finitely additive measure which is invariant under its action (4).

The pair $\left(Q^{\mathbb{G}}, F\right)$ is a dynamical system and can be studied from the point of view of topological dynamics. One of the important properties a CA can have is expansivity, which is a notion of strong sensitivity to initial conditions.

Given $\mathbb{T}$ be equal to either $\mathbb{N}$ or $\mathbb{Z}$, a CA $F$ is called $\mathbb{T}$-expansive if there exists a number $\delta>0$, called expansivity constant, such that for every $c \neq d$ there exists an instant $t \in \mathbb{T}$ such that $\Delta\left(F^{t}(c), F^{t}(d)\right) \geq \delta$. In this work we will frequently omit the prefix $\mathbb{T}$, assuming by default that $\mathbb{T}=\mathbb{N}$ (positive expansivity).

In $\mathbb{G}=\mathbb{Z}$, given a one-dimensional CA with local rule $f: Q^{[-l, r]} \rightarrow Q$ with $l, r>0$, we say that it is $L R$-permutif when for any $q_{-l}, \ldots, q_{r} \in Q$ the two following
maps are bijective:

$$
\begin{aligned}
a & \mapsto f\left(a, q_{-l+1}, \ldots, q_{r}\right) \\
a & \mapsto f\left(q_{-l}, \ldots, q_{r-1}, a\right)
\end{aligned}
$$

Definition 2. Let $(Q, \oplus)$ be a finite group and denote by $\bar{\oplus}$ the component-wise over $Q^{\mathbb{G}}$ is linear if

$$
\forall c, d \in Q^{\mathbb{G}}: F(c \bar{\oplus} d)=F(c) \bar{\oplus} F(d)
$$

The following lemma shows that linear CA over Abelian groups can be decomposed according to the structure of the group. It is a folklore knowledge that appears often in the particular case of cyclic groups $[\mathbf{1 7}, \mathbf{1 8}]$, and also in the more general Abelian case [19].

Recall that the product $F \times G$ of two CA $F$ and $G$ is the CA defined on the product alphabet and applying $F$ and $G$ on each component independently.

Lemma 1. Let $Q=Q_{p} \times Q^{\prime}$ be an Abelian group (law + and neutral element $(0,0)$ ) where $Q_{p}$ is a p-group (the order of every element is a power of $p$ ) for some prime

Then, any linear $C A F$ over $Q$ is isomorphic to $F_{p} \times F^{\prime}$ where $F_{p}$ is a linear $C A$ over $Q_{p}$ and $F^{\prime}$ is a linear $C A$ over $Q^{\prime}$.
Proof. By linearity of $F$, if a configuration $c$ is such that $n \cdot c=\underbrace{c+\cdots+c}_{n}=\overline{(0,0)}$, then $F(c)$ must be such that $n \cdot F(c)=\overline{(0,0)}$. We deduce that the subset of 95 states $Q_{1}=Q_{p} \times\left\{0_{Q^{\prime}}\right\}$ induces a subautomaton $F_{1}$ of $F$ because any configuration $c \in Q_{1}^{\mathbb{G}}$ is such that $p^{k} \cdot c=\overline{(0,0)}$ for some $k$ and no configuration in $\left(Q \backslash Q_{1}\right)^{\mathbb{G}}$ has this property. Moreover if $n \cdot c=\overline{(0,0)}$ for $n$ relatively prime with $p$, it implies that $c \in Q_{2}^{\mathbb{G}}=\left(\left\{0_{Q_{p}}\right\} \times Q^{\prime}\right)^{\mathbb{G}}$ (because the order of an element must divide the order of the group it belongs to). Therefore $Q_{2}$ induces a subautomaton $F_{2}$ of $F$.

Now, any $c \in Q^{\mathbb{G}}$ can be written $c=c_{1}+c_{2}$ where $c_{1} \in Q_{1}^{\mathbb{G}}$ and $c_{2} \in Q_{2}^{\mathbb{G}}$ through cellwise and componentwise decomposition, and $F(c)=F_{1}\left(c_{1}\right)+F_{2}\left(c_{2}\right)$. $F_{1}$ is isomorphic to a linear CA $F_{p}$ over $Q_{p}$ and $F_{2}$ to a linear CA $F^{\prime}$ over $Q^{\prime}$, and then $F$ is isomorphic to $F_{p} \times F^{\prime}$.

## 3. Pre-expansivity

Pre-expansivity is the property of expansivity restricted to asymptotic pairs of configurations.

Definition 3. Let $\mathbb{T}$ be either $\mathbb{N}$ or $\mathbb{Z}$ and let $F$ be a cellular automaton over $Q^{\mathbb{G}}$, supposed reversible in the case $\mathbb{T}=\mathbb{Z} . F$ is $\mathbb{T}$-pre-expansive if:

$$
\exists \delta>0: \forall c, d \in Q^{\mathbb{G}}, c \neq d \text { and } c \stackrel{\infty}{=} d \Rightarrow \exists t \in \mathbb{T}, \Delta\left(F^{t}(c), F^{t}(d)\right)>\delta
$$

The value $\delta$ is the expansivity constant.

In the sequel, when $\mathbb{T}$ is not explicitly mentioned, we will always refer to $\mathbb{T}=\mathbb{N}$. In particular this choice does not require the hypothesis of reversibility on the cellular automaton considered.

The notion of pre-expansivity can be further refined by considering only pairs of configurations with a fixed finite number of differences. Given $c, d \in Q^{\mathbb{G}}$, we denote $c \neq{ }_{k} d$ if $\#\{z \in \mathbb{G}: c(z) \neq d(z)\}=k$, i.e. if $c$ and $d$ differ in exactly $k$ positions.

Definition 4. Let $\mathbb{T}$ be either $\mathbb{N}$ or $\mathbb{Z}$, let $F$ be a cellular automaton over $Q^{\mathbb{G}}$ (supposed reversible in the case $\mathbb{T}=\mathbb{Z}$ ) and let $k>0$. $F$ is $\mathbb{T}$ - $k$-expansive $i f$ :

$$
\exists \delta>0: \forall c, d \in Q^{\mathbb{G}}, c \neq{ }_{k} d \Rightarrow \exists t \in \mathbb{T}, \Delta\left(F^{t}(c), F^{t}(d)\right)>\delta
$$

Denote by $T_{m}: Q^{\mathbb{G}} \rightarrow\left(Q^{B_{m}(0)}\right)^{\mathbb{N}}$ the trace function which to any configuration associates its orbit restricted to $B_{m}(0)$ :

$$
T_{m}(c)=\left(\left.t \mapsto\left(F^{t}(c)\right)\right|_{B_{m}}\right) .
$$

Proposition 1. Let $F$ be any $C A$ over $Q^{\mathbb{G}}$ where $\mathbb{G}$ is amenable. It holds:

- $F$ is pre-expansive $\Rightarrow \forall k>0 F$ is $k$-expansive,
- $F$ is $k$-expansive $\Rightarrow F$ is sensitive to initial configurations,
- $F$ expansive $\Rightarrow F$ pre-expansive $\Rightarrow F$ pre-injective $\Rightarrow F$ surjective,
- $F$ is pre-expansive $\Leftrightarrow T_{m}$ is pre-injective for some $m$.
- If $L$ is a $C A$ over $\tilde{Q}^{\mathbb{G}}$, then $F \times L$ is $k^{\prime}$-expansive for every $k^{\prime} \leq k$ if and only if $L$ and $F$ are $k^{\prime}$-expansive for every $k^{\prime} \leq k$.

Proof.

- The first item follows directly from definitions.
- For the second item it is sufficient to note that for any configuration $c$, any $\delta>0$ and any $k \geq 1$ there always exist a configuration $c^{\prime}$ with $c \neq{ }_{k} c^{\prime}$ and $\Delta\left(c, c^{\prime}\right) \leq \delta$.
- For the third item, it is clear that expansivity implies pre-expansivity (restriction of the universal quantification). Then pre-expansivity implies pre-injectivity because if there is a pair of configurations $c, c^{\prime}$ with $c \stackrel{\infty}{=} c^{\prime}$ and $F(c)=F\left(c^{\prime}\right)$ then, eventually applying a translation, we can also suppose them such that $\Delta\left(c, c^{\prime}\right)$ is arbitrary small. Finally, since $\mathbb{G}$ is amenable we have that pre-injectivity implies surjectivity by Garden of Eden Theorem 4].
- For the fourth item, it is sufficient to note that the existence of some time $t$ such that $\Delta\left(F^{t}(c), F^{t}\left(c^{\prime}\right)\right)>\delta$ is equivalent to $T_{m}(c) \neq T_{m}\left(c^{\prime}\right)$ for a suitable choice of $m$.
- Finally, if $F \times L$ is $k^{\prime}$-expansive, it is enough to take two configurations with their differences in only one of their components, since both automata act independently, $k^{\prime}$ expansivity of $F \times L$, imply that the perturbations will arrive to the center at the same component, proving the $k^{\prime}$-expansivity of the corresponding automaton.
If now $L$ and $F$ are $k^{\prime}$-expansive for every $k^{\prime} \leq k$, we take two configurations with $k^{\prime}$ differences. They may lay in one or both of their components, in any way there will be $0<k^{\prime \prime} \leq k^{\prime}$ differences in one of the components of $F \times L$. By the $k^{\prime \prime}$-expansivity of the corresponding automaton, we show the expansivity of $F \times L$.

Note however that $k$-expansivity does not generally imply pre-injectivity or surjectivity as shown by the following example.

Proposition 2. For any $k \geq 1$ there exists a $C A$ which is not surjective but $k^{\prime}$ expansive for any $k^{\prime} \leq k$.

Proof. Consider any pre-expansive one-dimensional CA $F$ of radius 1 over state set $Q=\{0,1\}$ (for instance a bi-permutative CA), and define a CA $\Psi$ over state set $Q^{k+1}$ as follows. It has $k+1$ "layers" and to any configuration $c$ we associate its projection $\pi_{i}(c)$ on the $i$ th layer. Intuitively it behaves on the $k$ first layers as $k$ independent copies of $F$, except that the $(k+1)$ th layer induce a state flip in the image in the following way: if it has a 1 at position $z$ then, in the image, layer $i$ is flipped at position $z+3 i$. Moreover, $(k+1)$ th layer is uniformly reset to 0 after one step. Formally, $\Psi$ is defined by:

$$
\Psi(c)_{z}=\left(F\left(\pi_{1}(c)\right)_{z}+\pi_{k+1}(c)_{z-3} \bmod 2, \ldots, F\left(\pi_{k}(c)\right)_{z}+\pi_{k+1}(c)_{z-3 k} \bmod 2,0\right)
$$

First it is clear from the definition that it is not surjective since the image of any configuration is always 0 on layer $k+1$. Note also that, when reduced to state set $\{0,1\}^{k} \times\{0\}, \Psi$ is isomorphic to $F^{k}$ which is pre-expansive. Therefore, to show that $\Psi$ is $k^{\prime}$-expansive for any $1 \leq k^{\prime} \leq k$, it is sufficient to show that for any pair of configurations $c$ and $d$ with $c \not \mathcal{k}_{k^{\prime}} d$ we have $\Psi(c) \neq \Psi(d)$.

So consider such a pair $(c, d)$. $\Psi$ was defined such that, if $c$ and $d$ differ on the $(k+1)$ th layer at position $z$, then, on the $i$ th layer, $F(c)$ and $F(d)$ will differ at position $z+3 i$ as soon as $c$ and $d$ are the same on the $i$ th layer at positions $z+3 i-1, z+3 i$ and $z+3 i+1$. Therefore, supposing that $c$ and $d$ indeed differ on the $(k+1)$ th layer at position $z$, it implies that $F(c)$ and $F(d)$ differ because $c$ and $d$ having only $k^{\prime} \leq k$ differences, they can not differ at $z$ and at one of the positions $z+3 i-1, z+3 i$ or $z+3 i+1$ for each $1 \leq i \leq k$.

Finally, suppose that $c$ and $d$ are equal on the $(k+1)$ th layer. Then they must differ on some layer $i$ with $1 \leq i \leq k$. Therefore, we must have $F\left(\pi_{i}(c)\right) \neq F\left(\pi_{i}(d)\right)$. We deduce that $\Psi(c) \neq \Psi(d)$ because their respective $i$ th layers are $F\left(\pi_{i}(c)\right)$ and $F\left(\pi_{i}(d)\right)$ up to some modification by the $(k+1)$ th layer which are identical in $c$ and $d$.

The next lemma talks about linear CA. When $F$ is supposed to be linear (for law $\oplus)$, then $T_{m}$ is also linear, i.e. $T_{m}(c \bar{\oplus} d)=T_{m}(c) \Phi T_{m}(d)$ where $\bar{\oplus}$ denotes the component-wise application of $\oplus$ either on $Q^{\mathbb{G}}$ or on $\left(Q^{B_{m}}\right)^{\mathbb{N}}$.

Proposition 3. Let $F$ be a linear $C A$ for law $\oplus$ and neutral element 0 . Let $I$ be

Proof. First, by linearity of the trace functions $T_{m}$, we have that $T_{m}(c)=T_{m}\left(c^{\prime}\right)$ if and only if $T_{m}\left(c^{\prime} \bar{\oplus}(-c)\right)=T_{m}(\overline{0})$ where $-c$ is the configuration such that $c \bar{\oplus}(-c)=\overline{0}$. Moreover we also have $c \not \neq k_{k} c^{\prime}$ if and only if $c^{\prime} \bar{\oplus}(-c) \neq{ }_{k} \overline{0}$. Hence $F$ is pre-expansive (resp. $k$-expansive) if and only if there is $m$ such that no $c \neq \overline{0}$ with $c \stackrel{\infty}{=} \overline{0}$ (resp. $\left.c \neq{ }_{k} \overline{0}\right)$ can verify $T_{m}(c)=T_{m}(\overline{0})$.

From this we deduce the second item of the Proposition.
For the first item, consider $k_{1}$ and $k_{2}$ in $I$. From what we said above, for any $m_{1}$ there is $c_{1}$ such that $c_{1} F_{k_{1}} \overline{0}$ and $T_{m_{1}}\left(c_{1}\right)=T_{m_{1}}(\overline{0})$. Now choose $m_{2}$ large enough so that any non-zero state of $c_{1}$ appears at distance at most $m_{2}$ from the center. Let us remark that the differences between $c_{1}$ and $\overline{0}$ are outside $B_{m_{1}}$, otherwise $T_{m_{1}}\left(c_{1}\right) \neq T_{m_{1}}(\overline{0})$, thus $m_{2}>m_{1}$. Since $k_{2} \in I$ we deduce from what we said earlier that there is $c_{2}$ such that $c_{2} \not{ }_{k_{2}} \overline{0}$ and $T_{m_{2}}\left(c_{2}\right)=T_{m_{2}}(\overline{0})$. By choice of $m_{2}$, this implies that $T_{m_{1}}\left(c_{1} \bar{\oplus} c_{2}\right)=\mathbb{T}_{m_{1}}(\overline{0})$. Moreover $c_{1} \bar{\oplus} c_{2} \not{ }_{k_{1}+k_{2}} \overline{0}$. Since $m_{1}$ was arbitrary, we deduce that $F$ is not $\left(k_{1}+k_{2}\right)$-expansive.

## 4. 1-dimensional Cellular Automata

205 The main goal of this section is to show that the notion of $k$-expansivity, preexpansivity and expansivity all differ. More precisely we will show the following existential result.

Theorem 1. For each item in the following list, there exists a CA having the given properties:

- $\mathbb{N}$-pre-expansive and reversible and not $\mathbb{N}$-expansive,
- $\mathbb{N}$-pre-expansive and irreversible and not $\mathbb{N}$-expansive,
- 1-pre-expansive and reversible and not $\mathbb{N}$-pre-expansive, and
- 1-pre-expansive and irreversible and not $\mathbb{N}$-pre-expansive.

For this purpose it will be sufficient to focus on linear cellular automata.
notion.

Definition 5. Given two configurations $c \neq d$, and a $C A F$, we define the left and right propagation sequences as follows.

$$
\begin{aligned}
l_{t}^{d}(c) & =\min \left\{z \in \mathbb{Z}:\left(F^{t}(c)\right)(z) \neq\left(F^{t}(d)\right)(z)\right\} \\
r_{t}^{d}(c) & =\max \left\{z \in \mathbb{Z}:\left(F^{t}(c)\right)(z) \neq\left(F^{t}(d)\right)(z)\right\}
\end{aligned}
$$

Lemma 3. Given a CA $F$ of neighborhood $[-l, r]$ and $k \in \mathbb{N}$, the next assertion hold.

1. If $F$ is $k$-expansive, then $\forall c \not \neq k d,\left(l_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not lower bounded and $\left(r_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not upper bounded.
2. If $\forall k^{\prime} \leq k, \forall c \neq{ }_{k^{\prime}} d,\left(l_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not lower bounded and $\left(r_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not upper bounded, then $F$ is $k$-expansive with expansivity constant $2^{-\max \{l, r\}}$.

Proof.

1. Let us suppose that $F$ is $k$-expansive with expansivity constant $m$. Let $c \not \neq k d$ be two configurations, and let us assume that $l_{t}^{d}(c)>n$ for some $n \in \mathbb{Z}$; this means that $F^{t}(c)_{]-\infty, n]}=F^{t}(d)_{]-\infty, n]}$ for every $t \in \mathbb{N}$. Thus $T_{m}\left(\sigma_{n-m}(c)\right)=T_{m}\left(\sigma_{n-m}(d)\right)$, which is a contradiction. The analogous happens if $\left(r_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is upper bounded.
2. We need to prove that $F$ is $k$-expansive with expansivity constant $2^{-m}=$ $2^{-\max \{l, r\}}$. Let $c \neq{ }_{k} d$ be two configurations, and let us define the next two additional configurations.

$$
\begin{aligned}
& c^{l}(i)= \begin{cases}c(i) & \text { if } i<0 \\
d(i) & \text { if } i \geq 0\end{cases} \\
& c^{r}(i)= \begin{cases}d(i) & \text { if } i<0 \\
c(i) & \text { if } i \geq 0\end{cases}
\end{aligned}
$$

If $c_{[-m, m]} \neq d_{[-m, m]}, T_{m}(c) \neq T_{m}(d)$ and we are done, so let us suppose that $c_{[-m, m]}=d_{[-m, m]}$. Let $k^{\prime}$ and $k^{\prime \prime}$ be such that $c^{r} \neq k_{k^{\prime}} d, c^{l} \not{ }_{k^{\prime \prime}} d$ and $k^{\prime}+k^{\prime \prime}=k$.
If $k^{\prime} \neq 0,\left(l_{t}^{d}\left(c^{r}\right)\right)_{t \in \mathbb{N}}$ is not lower bounded, thus by lemma 2 and the fact that $c^{r}$ is equal to $d$ below position $m$, there is a minimal iteration $t_{r}$ such that $F^{t_{r}}\left(c^{r}\right)_{[m-r, m]} \neq F^{t_{r}}(d)_{[m-r, m]}$. Analogously, if $k^{\prime \prime} \neq 0$ there is a minimal iteration $t_{l}$ such that $F^{t_{l}}\left(c^{l}\right)_{[-m,-m+l]} \neq F^{t_{l}}(d)_{[-m,-m+l]}$.
Let us take $\bar{t}=\min \left\{t_{r}, t_{l}\right\}$, by the choice of $m$ we have that $F^{\bar{t}}(c)_{[0, m]}=$ $F^{\bar{t}}\left(c^{r}\right)_{[0, m]}$ and $F^{\bar{t}}(c)_{[-m, 0]}=F^{\bar{t}}\left(c^{l}\right)_{[-m, 0]}$, and at least one of them is different from $F^{\bar{t}}(d)$ between $[-m, m]$, thus $T_{m}(c) \neq T_{m}(d)$.

The last lemma establishes that the expansivity constant is uniform in dimension one (it does not depends on $k$ ), thus we can conclude the next corollary.

Corollary 1. If $F$ is $k$-expansive for every $k \in \mathbb{N}$ then it is pre-expansive.
Lemma 4. $F$ is pre-expansive if and only if for all $c \stackrel{\infty}{=} d,\left(l_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not lower bounded and $\left(r_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not upper bounded.

Proof.
$265 \quad(\Rightarrow)$ Let $c \not F_{k} d$ be two asymptotics configurations. Since $F$ is pre-expansive, it is also $k$-expansive, thus by lemma 3, $\left(l_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not lower bounded and $\left(r_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not upper bounded.
$(\Leftarrow)$ If for every $k \in \mathbb{N}$ and every pair $c \not \neq k_{k} d$ we have that $\left(l_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not lower bounded and $\left(r_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not upper bounded, then by lemma 3 $F$ is $k$ expansive, and by corollary we conclude that $F$ is pre-expansive.

Left and right propagation determines also expansivity. The next lemma can be proven by using the techniques from the last lemmas.

LEMMA 5. $F$ is expansive if and only if for any pair of different configurations $c, d$, if $c_{]-\infty, 0]}=d_{]-\infty, 0]}$ then $\left(l_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not lower bounded, and if $c_{[0, \infty[ }=d_{[0, \infty[ }$ then $\left(r_{t}^{d}(c)\right)_{t \in \mathbb{N}}$ is not upper bounded.

The following proposition shows that the simplest form of linearity is not sufficient to achieve the separation between expansivity and pre-expansivity. Let us first introduce some notation. Given $a \in Q$, we denote by $c^{a}$ the configuration that is equal to 0 everywhere except at cell 0 where its value is $a$.

Proposition 4. Let $\mathbb{Z}_{n}$ be the the group of integers modulo $n$ with addition, and let $F$ be a one-dimensional linear $C A$ over $\mathbb{Z}_{n}$. Then $F$ is 1 -expansive if and only if it is expansive.
Proof. First by Proposition 1 if $F$ is expansive it is in particular 1-expansive.
For the other direction, it is sufficient to consider the case $n=p^{k}$ with $p$ a prime number by lemma 1 because if some $F_{1} \times F_{2}$ is 1-expansive then both $F_{1}$ and $F_{2}$ must be 1-expansive.

By commutation with shifts we have $l_{t}^{\overline{0}}\left(\sigma_{n}\left(c^{a}\right)\right)=-n+l_{t}^{\overline{0}}\left(c^{a}\right)$ and the analogous for $r_{t}$. Moreover $c^{a}=a \cdot c^{1}$ because we are on a cyclic group.

Let us define $l_{t}^{U}\left(c^{a}\right)=\min \left\{i \in \mathbb{Z} \mid \operatorname{gcd}\left(F^{t}\left(c^{a}\right)_{i}, p\right)=1\right\}$. The next properties hold.

1. If $\operatorname{gcd}(a, p) \neq 1$, then $l_{t}^{U}\left(c^{a}\right)=\infty$ (respectively $\left.r_{t}^{U}\left(c^{a}\right)=-\infty\right)$. In fact, in this case the entire evolution of $F$ over $c^{a}$ is composed by multiples of $a$, which are multiples of $p$ too.
2. If $\operatorname{gcd}(a, p)=1$, then $l_{t}^{U}\left(c^{a}\right)=l_{t}^{U}\left(c^{1}\right)$ (respectively $r_{t}^{U}\left(c^{a}\right)=r_{t}^{U}\left(c^{1}\right)$ ). In fact, in this case $F^{t}\left(c^{a}\right)_{i}=a F^{t}\left(c^{1}\right)_{i}$ is coprime with $p$ if and only if $F^{t}\left(c^{1}\right)_{i}$ does.
3. $\quad l_{t}^{\overline{0}}\left(c^{a}\right) \leq l_{t}^{U}\left(c^{1}\right)$ (respectively $r_{t}^{\overline{0}}\left(c^{a}\right) \geq r_{t}^{U}\left(c^{1}\right)$ ). In fact, if $F^{t}\left(c^{1}\right)_{i}$ is coprime with $p$, then $F^{t}\left(c^{a}\right)_{i}=a F^{t}\left(c^{1}\right)_{i}$ is not null.
4. $\quad l_{t}^{\overline{0}}\left(c^{p^{k-1}}\right)=l_{t}^{U}\left(c^{1}\right)$ (respectively $r_{t}^{\overline{0}}\left(c^{p^{k-1}}\right)=r_{t}^{U}\left(c^{1}\right)$. In fact, $F^{t}\left(c^{p^{k-1}}\right)_{i}=$ $p^{k-1} F^{t}\left(c^{1}\right)=0 \Leftrightarrow \operatorname{gcd}\left(F^{t}\left(c^{1}\right)_{i}, p\right) \neq 1$.
From the last assertion and lemma 3 (first item) we have that $l_{t}^{U}\left(c^{1}\right)$ is not lower bounded and $r_{t}^{U}\left(c^{1}\right)$ is not upper bounded.

Now let us take a configuration $v \in\left(\mathbb{Z}_{p^{k}}\right)^{\mathbb{Z}}$, such that $v_{i}=0$ for every $i<0$ and $v_{0} \neq 0$. Let us define $j=\max \left\{i \in\{0, \ldots, k\}\left|\forall x \in \mathbb{Z}, p^{i}\right| v(x)\right\}$, and let us consider $u=v / p^{j}$. In this way, there is $y \geq 0$ with $u(y) \neq 0$ and $\operatorname{gcd}(p, u(y))=1$. Let $y$ be the smallest integer with this property.

$$
F^{t}(u)_{l_{t}^{U}\left(c^{1}\right)+y}=\sum_{x=0}^{l_{t}^{U}\left(c^{1}\right)+y+r t} F^{t}\left(c^{u_{x}}\right)_{l_{t}^{U}\left(c^{1}\right)+y-x}
$$

But $u(x)$ is a multiple of $p$ when $x<y$, thus $F^{t}\left(c^{u_{x}}\right)_{i}=0 \bmod p$ for any $i$. For $x>y, F^{t}\left(c^{u_{x}}\right)_{l_{t}^{U}\left(c^{1}\right)+y-x}$ is also a multiple of $p$ because the smallest index for which $F^{t}\left(c^{u_{x}}\right)$ is coprime with $p$ is $l_{t}^{U}\left(c^{u_{x}}\right)$ which is greater than or equal to $l_{t}^{U}\left(c^{1}\right)$. In this way we conclude that

$$
F^{t}(u)_{l_{t}^{U}\left(c^{1}\right)+y}=F^{t}\left(c^{u_{y}}\right)_{l_{t}^{U}\left(c^{1}\right)} \bmod p
$$

is coprime with $p$ and is non null. Therefore $F^{t}(v)_{l_{t}^{U}\left(c^{1}\right)+y}$ is non null as well. This imply that $l_{t}^{\overline{0}}(v)$ is not lower bounded. Symmetrically, $r_{t}^{\overline{0}}(w)$ is not upper bounded when $w$ is any configuration equal to zero on positive coordinates. Lemma 5 concludes that $F$ is (positively) expansive.

To establish a separation between expansivity and pre-expansivity, we will focus on linear CA obtained by what is often called "second order method" in the literature [20. The idea is to turn any CA into a reversible one by memorizing one
step of history and combining, in a reversible way, the memorized past step into the produced future step. The interest of this construction for our purpose is that reversible at the same time [21.

Let $Q=\{0, \ldots, n-1\}$ be equipped with some group law $\oplus$ and consider some CA $F$ over state set $Q$. The second-order CA associated to $F$ and $\oplus$, denoted $\mathcal{S O}(F, \oplus)$ is the CA over state set $Q \times Q$, which is conjugated through the natural

$$
(c, d) \mapsto(d, F(d) \bar{\oplus} c)
$$

The following proposition shows that second order construction is useful to separate expansivity from 1-expansivity. Some of the results in the next proposition can be deduced from more general results in [16, 22, 23, but we develop a new specific proof here.
${ }_{330}$ Proposition 5. Let $\oplus$ be a group law over $Q$ with neutral element 0 and $F$ be a $C A$ over $Q$ which is linear for $\oplus$. It holds:

- $\mathcal{S O}(F, \oplus)$ is bijective and linear for the law $\oplus \times \oplus$;
- if $F$ is LR-permutive then $\mathcal{S O}(F, \oplus)$ is $\mathbb{Z}$-expansive and 1-expansive;
- if $F$ is LR-permutive then for any $m>0$ the subshift of traces $T_{m}\left((Q \times Q)^{\mathbb{Z}}\right)$ is an SFT.

Proof. For the first item it is sufficient to check that the CA over the state set $Q \times Q$ is conjugated to the following map:

$$
(c, d) \mapsto(\bar{\iota}(F(c)) \bar{\oplus} d, c)
$$

is the inverse of $\mathcal{S O}(F, \oplus)$, where $\iota$ denotes the inverse function for the group law $\oplus$. Moreover $\mathcal{S O}(F, \oplus)$ is linear for $\oplus \times \oplus$ because it is component-wise linear

For the second item, suppose that $F$ is LR-permutive with neighborhood $\{-l, \ldots, r\}$ and denote $\Psi=\mathcal{S O}(F, \oplus)$. To prove $\mathbb{Z}$-expansivity of $\Psi$ it is sufficient to notice that $\Psi$ propagates to left and right when the second $Q$-component is nonnull and $\Psi^{-1}$ propagates to left and right when the first $Q$-component is non-null. ${ }_{35}$ In fact, let us consider a configuration $c \in(Q \times Q)^{\mathbb{Z}}$ equal to $(0,0)$ on negative coordinates but such that $c(0) \neq(0,0)$. If the second $Q$-component of $c(0)$ is non-null then, and since $F(0, \ldots, 0)=0$, the leftmost non-null cell of $\Psi(c)$ is at position $-r$ and it is its second $Q$-component which is non-null, i. e., if the leftmost difference from $(0,0)$ is in the second $Q$-component, this will be always like this 0 and the difference will propagate to the left. The same holds symmetrically for propagation to the right. In the same way, and given the form of $\Psi^{-1}$, differences in the first $Q$-component will propagate to the left and right through $\Psi^{-1}$, thus by lemma $5 \Psi$ is $\mathbb{Z}$-expansive.

Now let us take $c^{(a, b)}$ as a configuration equal to $(a, b)$ at 0 and $(0,0)$ everywhere else. By the previous arguments, if $b \neq 0,\left(l_{t}^{\overline{(0,0)}}\left(c^{(a, b)}\right)\right)_{t \in \mathbb{N}}$ is not lower bounded and $\left(r_{t}^{\overline{(0,0)}}\left(c^{(a, b)}\right)\right)_{t \in \mathbb{N}}$ is not upper bounded. But if $b=0$ and $a \neq 0$, then $(\Psi(c))(0)=(0, a)$ and null everywhere else, then again $\left(l_{t}^{\overline{(0,0)}}\left(c^{(a, b)}\right)\right)_{t \in \mathbb{N}}$ is not lower bounded and $\left(r_{t}^{(0,0)}\left(c^{(a, b)}\right)\right)_{t \in \mathbb{N}}$ is not upper bounded. Therefore, $\Psi\left(\right.$ and $\left.\Psi^{-1}\right)$ is 1 -expansive.

The proof of the third item will be performed in two steps. To simplify notations, for any pair of words $u, v \in Q^{*}$ of the same length, we will denote by $\binom{u}{v}$ the word over alphabet $Q \times Q$ whose projection on the first (resp. second) component is $u$ (resp. $v$ ).

Assertion 1L: For every word $u \in Q^{r}$ there exists a configuration $c$ such that $\left.\Psi^{t}(c)\right|_{[0, r-1]}=\binom{u}{0_{r}^{r}}$ and $\Psi^{k}(c)_{i}=(0,0)$ for every $0 \leq k \leq t$ and $0 \leq i<(t-k) r$.
Proof of Assertion 1L. By induction on $t$. If $t=0$ it is obvious, we just take $c$ equal to $\binom{u}{0^{r}}$ at $[0, r-1]$ and $(0,0)$ everywhere else. Now, since $F$ is LR-permutive, given a word $w \in Q^{l+r}$, let us define the permutation $\tau_{w}(a)=f(w a)$ for every $a \in Q$. Given a word $u \in Q^{r}$, we inductively define another word $v \in Q^{r}$ as follows: $v_{0}=\tau_{0^{l+r}}^{-1}\left(u_{0}\right), v_{i+1}=\tau_{0^{1+r-i} v_{[0, i]}}^{-1}\left(u_{i+1}\right)$. In this way, $f\left(0^{l+r} v\right)=u$. By induction hypothesis, there exists a configuration $c$ such that $\left.\Psi^{t}(c)\right|_{[0, r-1]}=\binom{v}{0^{r}}$ and $\Psi^{k}(c)_{i}=(0,0)$ for every $0 \leq k \leq t$ and $i<(t-k) r$. We take $d=\sigma_{-r}(c)$, then $\left.\Psi^{t}(d)\right|_{[r, 2 r-1]}=\binom{v}{0^{r}}$, and $(0,0)$ to the left of $r$; and $\left.\Psi^{t+1}(d)\right|_{[0, r-1]}=\binom{u}{0^{r}}$, and $(0,0)$ to the left of 0 . Moreover, $\Psi^{k}(d)_{i}=(0,0)$ for every $0 \leq k \leq t+1$ and $0 \leq i<(t+1-k) r$.

Assertion 1R: For every word $u \in Q^{l}$ there exists a configuration $c$ such that $\left.\Psi^{t}(c)\right|_{[-l+1,0]}=\binom{u}{0^{t}}$ and $\Psi^{k}(c)_{i}=(0,0)$ for every $0 \leq k \leq t$ and $i>(k-t) l$. The proof of Assertion $1 R$ is analogous to the proof of Assertion $1 L$.

Assertion 2: A sequence $\left(w_{t}\right)_{t=0}^{n}$, with $w_{t} \in(Q \times Q)^{2 m+1}$, is a finite subsequence of a trace in $T_{m}\left((Q \times Q)^{\bar{Z}}\right)$ if and only if for any $t$, there are extensions $w_{R} \in(Q \times Q)^{r}$ and $w_{L} \in(Q \times Q)^{l}$ verifying

$$
\psi\left(w_{L} \cdot w_{t} \cdot w_{R}\right)=w_{t+1}
$$

where $\psi$ denotes the action of $\Psi$ over finite words.
Proof of Assertion 2. In one direction, it is clear, so let $\left(w_{t}\right)_{t=0}^{n}$ be a sequence such that for any $t$, there are extensions $w_{R} \in(Q \times Q)^{r}$ and $w_{L} \in(Q \times Q)^{l}$ verifying $\psi\left(w_{L} \cdot w_{t} \cdot w_{R}\right)=w_{t+1}$, and let us prove that it is a subsequence of a trace of $\Psi$. We perform the proof by induction on $n$. If $n=0$ there is nothing to prove, of course any sequence of length 1 can be part of a trace. Now let $c$ be a configuration such that $\left.\Psi^{k}(c)\right|_{[-m, m]}=w_{k}$, for every $k \in\{0, \ldots, n-1\}$. By locality, only the values of $c$ between $-m-n l$ and $m+n r$ are relevant to this hypothesis, and we take $c(i)=(0,0)$ outside these
limits. Let $w_{R}$ and $w_{L}$ be such that $\psi\left(w_{L} \cdot w_{n-1} \cdot w_{R}\right)=w_{n}$. Let us remark that the first $Q$-component of $w_{L}$ and $w_{R}$ can be chosen arbitrarily, given the form of $\Psi$. We will suppose that $\pi_{1}\left(w_{L}\right)=\pi_{1}\left(\Psi^{n-1}(c)_{[-m-l,-m-1]}\right)$ and $\pi_{1}\left(w_{R}\right)=\pi_{1}\left(\Psi^{n-1}(c)_{[m+1, m+r]}\right)$, thus we can take $\binom{u_{L}}{0^{l}}=w_{L}-$ $\left.\Psi^{n-1}(c)\right|_{[-m-l,-m-1]}$ and $\binom{u_{R}}{0^{r}}=w_{R}-\left.\Psi^{n-1}(c)\right|_{[m+1, m+r]}$. We take from Assertion 1L a configuration $c^{R}$ that produces the word $\binom{u_{R}}{0^{r}}$ at time $n-1$ at position $[m+1, m+r]$ and $(0,0)$ to the left of the light cone that starts at $m+n r$ with slope $-1 / r$. From Assertion 1R we take a configuration $c^{L}$ that produces the word $\binom{u_{L}}{0^{r}}$ at time $n-1$ at position $[-m-l,-m-1$ ] and $(0,0)$ to the right of the light cone that starts at $-m-n l$ with slope $1 / l$. By linearity, $\left.\Psi^{n-1}\left(c^{L} \bar{\oplus} c^{R} \bar{\oplus} c\right)\right|_{[-m-l, m+r]}=w_{L} w_{n-1} w_{R}$, and then $\left.\Psi^{n}\left(c_{L} \bar{\oplus} c_{R} \bar{\oplus} c\right)\right|_{[-m, m]}=w_{n}$, moreover $\left.\Psi^{k}\left(c_{L} \bar{\oplus} c_{R} \bar{\oplus} c\right)\right|_{[-m, m]}=w_{k}$, for every $0 \leq k<n$, which completes the proof.

LEMMA 6. Let $c$ be a configuration in $(Q \times Q)^{\mathbb{Z}}$. Then, for any $k \geq 0$, any $t \geq 0$ and any $z \in \mathbb{Z}$ we have:

$$
\Psi_{z}^{2 \cdot 3^{k}+t}(c)=\Psi_{z}^{t} \oplus \Psi_{z-3^{k}}^{3^{k}+t} \oplus \Psi_{z+3^{k}}^{3^{k}+t}(c)
$$

Proof. First it is straightforward to check that

$$
\Psi_{0}^{2}(c)=\Psi_{-1}^{1} \oplus I d \oplus \Psi_{1}^{1}(c)
$$

Then, by Lucas' lemma, we have $(a+b)^{3^{k}}=a^{3^{k}}+b^{3^{k}}$ when doing the arithmetics modulo 3 . This identity naturally extends to $\oplus$ and therefore we have:

$$
\Psi_{0}^{2 \cdot 3^{k}}(c)=\Psi_{-3^{k}}^{3^{k}} \oplus I d \oplus \Psi_{3^{k}}^{3^{k}}(c)
$$

Finally, by linearity of both $\sigma$ and $\Psi$ with respect to $\oplus$ we can compose both sides of the above equality by $\Psi_{z}^{t}$ and the lemma follows.

Using the above lemma, we can show that $\Psi$ has a simple dependency structure at some space-time locations.


Figure 1. Some space-time diagrams of $\Psi$ (state $(0,0)$ is represented by empty space)

425 LEMMA 7. Let us consider a configuration $c=c^{(a, b)}$ for some pair $(a, b) \in Q \times Q$. For any $k \geq 0$, let $d_{k}=\Psi^{3^{k+1}}(c)$. Then we have:

- $d_{k}\left(-3^{k}\right)=d_{k}\left(3^{k}\right)=\phi(a, b)$
- $d_{k}(i)=(0,0)$ for $-3^{k}<i<3^{k}$
where $\phi$ is an automorphism of $Q \times Q$ which does not depend on $k$.
Proof. Denote by $c_{z}^{t}$ the state $\left(\Psi^{t}(c)\right)(z)$. First, by a simple recurrence we can show that $c_{-n}^{n}=c_{n}^{n}=\pi(a, b)$ where $\pi$ is the projection $\pi(a, b)=(0, b)$. Then, applying lemma 6 on $\Psi_{-3^{k}}^{2 \cdot 3^{k}}$, we obtain:

$$
d_{k}\left(-3^{k}\right)=c_{-2 \cdot 3^{k}}^{2 \cdot 3^{k}}+c_{-3^{k}}^{3^{k}}+c_{0}^{2 \cdot 3^{k}}
$$

Applying lemma 6 on $\Psi_{0}^{2 \cdot 3^{k}+0}$ we then have:

$$
\begin{aligned}
d_{k}\left(-3^{k}\right) & =c_{-2 \cdot 3^{k}}^{2 \cdot 3^{k}}+2 \cdot c_{-3^{k}}^{3^{k}}+c_{3^{k}}^{3^{k}}+c_{0}^{0} \\
& =(0, b)+2(0, b)+(0, b)+(a, b) \\
& =(a, 2 b)
\end{aligned}
$$

where $\phi(a, b)=(a, 2 b)$ is an automorphism. The same equality holds for $d_{k}\left(3^{k}\right)$ by symmetry and the first item of the lemma is shown.

For the second item, first note that $c_{z}^{t}=(0,0)$ whenever $z<-t$ or $z>t$ because $\Psi$ has radius 1. Consider any $i$ with $-3^{k}<i<3^{k}$. Applying lemma 6 on $\Psi_{i}^{2 \cdot 3^{k}+3^{k}}$, we obtain:

$$
d_{k}(i)=c_{i-3^{k}}^{2 \cdot 3^{k}}+c_{i}^{3^{k}}+c_{i+3^{k}}^{2 \cdot 3^{k}}
$$

Applying lemma 6 on $\Psi_{i-3^{k}}^{2 \cdot 3^{k}+0}$ and $\Psi_{i+3^{k}}^{2 \cdot 3^{k}+0}$ we further get:

$$
d_{k}(i)=c_{i-2 \cdot 3^{k}}^{3^{k}}+c_{i-3^{k}}^{0}+3 \cdot c_{i}^{3^{k}}+c_{i+3^{k}}^{0}+c_{i+2 \cdot 3^{k}}^{3^{k}}
$$

${ }_{440}$ From what we said before and doing the arithmetics modulo 3 we deduce $d_{k}(i)=(0,0)$ and the lemma follows.

Proposition 6. $\Psi$ is pre-expansive.
Proof. Let $c$ be any configuration with $c \stackrel{\infty}{=} \overline{(0,0)}$. Denote $l=l_{0}^{\overline{(0,0)}}(c)$ and $r=r_{0}^{\overline{(0,0)}}(c)$ and consider any $k$ such that $3^{k}>\max (|l|,|r|)$. By a finite number of applications of Lemma 7 and by linearity and translation invariance of $\Psi$, we have:

$$
\begin{aligned}
\left(\Psi^{3^{k+1}}(c)\right)\left(r-3^{k}\right) & =\phi(c(r)) \\
\left(\Psi^{3^{k+1}}(c)\right)\left(l+3^{k}\right) & =\phi(c(l))
\end{aligned}
$$

Since $\phi$ is a permutation of $Q$ sending $(0,0)$ to itself and since $c(r)$ and $c(l)$ are both different from $(0,0)$ we deduce that $l_{3^{k+1}}^{\overline{(0,0)}}(c)<r-3^{k}$ and $r_{3^{k+1}}^{\overline{(0,0)}}(c)>l+3^{k}$ for any $k$ large enough. This shows that $L(c)$ is not lower-bounded and $R(c)$ is not upper-bounded and concludes the proof by Lemma 4


Figure 2. Some space-time diagrams of $\Upsilon$ (state ( 0,0 ) is not represented)

We now give an example of CA that is 1-expansive but not pre-expansive.
Example $2(\Upsilon)$ Let $Q=\{0,1\}$, + be the addition modulo 2 , and $F_{2}$ be the $C A$ defined over $Q^{\mathbb{Z}}$ by $F_{2}=\sigma \mp \sigma_{-1}$. We define $\Upsilon$ as the second order construction applied to $F_{2}$ :

$$
\Upsilon=\mathcal{S O}\left(F_{2},+\right)
$$

Proposition 7. $\Upsilon$ is not $k$-expansive when $k \geq 2$ and in particular $\Upsilon$ is not preexpansive.

Proof. Let $k \geq 2$ be fixed and for each $z \in \mathbb{Z}$ define the configuration $c^{z}$ by:

$$
c^{z}\left(z^{\prime}\right)= \begin{cases}(0,0) & \text { if } z^{\prime}<z \\ (0,1) & \text { if } z^{\prime}=z \\ (1,1) & \text { if } z<z^{\prime} \leq z+k-2 \\ (1,0) & \text { if } z^{\prime}=z+k-1 \\ (0,0) & \text { if } z^{\prime} \geq z+k\end{cases}
$$

We have $c^{z} \neq k^{(0,0)}$ and it is straightforward to check that $\Upsilon\left(c^{z}\right)=c^{z-1}$. We 455 conclude that $\Upsilon$ is not $k$-expansive by lemma 4

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

Proof. [Proof of Theorem 1 Let $F$ be any irreversible and expansive CA. It holds:

- $\Psi$ is $\mathbb{N}$-pre-expansive (by proposition (6) and it is reversible and not $\mathbb{N}$ expansive (by proposition [5and [21);
- therefore $\Psi \times F$ is $\mathbb{N}$-pre-expansive and irreversible and not $\mathbb{N}$-expansive;
- $\Upsilon$ is 1-pre-expansive and reversible (by proposition (5) but it is not $\mathbb{N}$-preexpansive (by proposition (7);
- therefore $\Upsilon \times F$ is 1-pre-expansive and irreversible and not $\mathbb{N}$-pre-expansive.


## 5. Cellular automata over the free group

Some of the properties proved in the last section come from the fact that $(\mathbb{Z},\{(i, i+j) \mid j \in V\})$ can be always disconnected by extracting a finite part from $\mathbb{Z}$. Any free group with any finite neighborhood has this property, and we would be able to extend some of the previous properties to the case of a cellular automaton over the free group. In particular, the expansivity constant is strictly related with the neighborhood size, as in $\mathbb{Z}$, and it does not depend on $k$ for a $k$-expansive CA. We denote by $\mathbb{F}_{n}$ the free group with $n$ generators $\left(\mathbb{F}_{1}\right.$ is $\left.\mathbb{Z}\right)$.

Proposition 8. If $F$ is a cellular automaton with a neighborhood $V \subseteq B_{r}(0)$ of radius $r$ over the free group $\mathbb{F}_{n}$ and it is $k^{\prime}$-expansive for all $k^{\prime} \leq k$, then $F$ is $k$-expansive with expansivity constant equal to $2^{-r}$.

Proof. The proof takes the ideas of lemma 3. Let $c \neq{ }_{k} d$ be two configurations in $\mathbb{F}_{n}$. Let us call $S$ the set of generators of $\mathbb{F}_{n}$, including their inverses (i. e. $|S|=2 n$ ), and for each $s \in S$, let us call $R_{s}$ the branch of $\mathbb{F}_{n}$ that hangs from $s$; we mean the set of elements whose shortest description in terms of $S$ start with $s$. In this way $\mathbb{F}_{n}=\{0\} \sqcup\left(\sqcup_{s \in S} R_{s}\right)$.

Now let us define $D=\left\{i \in \mathbb{F}_{n} \mid c(i) \neq d(i)\right\}$ and $D_{s}=D \cap R_{s}$. We want to prove that $T_{r}(c) \neq T_{r}(d)$ so let us suppose the opposite. This imply that $D=\sqcup_{s \in S} D_{s}$, and we can consider $k_{s}=\left|D_{s}\right| \leq|D|=k$. As in the case of lemma 3, we define configurations $c^{s}$ which are equal to $d$ everywhere except on branch $s$, as follows.

$$
c^{s}(i)= \begin{cases}c(i) & \text { if } i \in R_{s} \\ d(i) & \text { otherwise }\end{cases}
$$

At the beginning, $c^{s}$ differs from $d$ only on branch $R_{s}$. We will see that this will be always the case. Let us suppose that, for some $t \in \mathbb{N}, F^{t}\left(c^{s}\right)_{i}=F^{t}(c)_{i}$ for all $i \in R_{s}$ and $F^{t}\left(c^{s}\right)_{j}=F^{t}(d)_{j}$, for all $j \notin R_{s}$. We remark that, we assumed that if $j \in B_{r}(0), F^{t}\left(c^{s}\right)_{j}=F^{t}(d)_{j}=F^{t}(c)_{j}$. Since $\mathbb{F}_{n}$ is a tree and $F$ is a CA of radius $r$, if $i \in R_{s}, B_{r}(i) \subset R_{s} \cup B_{r}(0)$, thus $F^{t+1}\left(c^{s}\right)_{i}=F^{t+1}(c)_{i}$. If $j \in\{0\} \sqcup\left(\sqcup_{s^{\prime} \in S \backslash\{s\}} R_{s^{\prime}}\right)$, $B_{r}(j) \subset\left(\sqcup_{s^{\prime} \in S \backslash\{s\}} R_{s^{\prime}}\right) \cup B_{r}(0)$, then $F^{t+1}\left(c^{s}\right)_{j}=F^{t+1}(d)_{j}$.

We conclude that $T_{r}\left(c^{s}\right)=T_{r}(c)=T_{r}(d)$, for every $s \in S$.

But we know, by hypothesis, that $F$ is $k^{\prime}$-expansive, let us take its expansivity

95 $c^{s}$ by $m s \in \mathbb{F}_{n}$. By construction, $\sigma_{-m s}\left(c^{s}\right)$ and $\sigma_{-m s}(d)$ are equal over $B_{m}(0)$. By the $k^{\prime}$-expansivity of $F$, there exists a time $t$ and $j \in B_{m}(0)$ such that $F^{t}\left(\sigma_{-m s}\left(c^{s}\right)\right)_{j} \neq F^{t}\left(\sigma_{-m s}(d)\right)_{j}$. But, as shown before, $F^{t}\left(c^{s}\right)$ and $F^{t}(d)$ differ only on branch $R_{s}$. Therefore $F^{t}\left(\sigma_{-m s}\left(c^{s}\right)\right)$ and $F^{t}\left(\sigma_{-m s}(d)\right)$ differ only on branch $R_{m s}$ which is disjoint from $B_{m}(0)$ : this is a contradiction.

Not all the properties survives from $\mathbb{F}_{1}=\mathbb{Z}$ to $\mathbb{F}_{n}$, when $n>1 ; k$-expansivity is possible for infinitely many $k$ 's in $\mathbb{F}_{n}$ even without pre-expansivity, as the next example shows.

Example $3\left(\Lambda_{n}\right)$ Let $Q=\{0,1\}$, + be the addition modulo 2 , and $\Lambda_{n}$ be the $C A$ defined over $Q^{\mathbb{F}_{n}}$ by

$$
\Lambda_{n}(c)_{i}=c(i)+\sum_{j \in S} c(i+j)
$$

Lemma 8. If $\|x\|=\|y\|$, then for every $t \in \mathbb{N} \Lambda_{n}^{t}\left(c^{1}\right)_{x}=\Lambda_{n}^{t}\left(c^{1}\right)_{y}$, moreover $\Lambda_{n}^{\|x\|}\left(c^{1}\right)_{x}=1$.

Proof. We prove by induction on $l$ that for every $t \leq l$ and every $x, y \in B_{l}(0)$, $\left[\|x\|=\|y\| \Rightarrow \Lambda_{n}^{t}\left(c^{1}\right)_{x}=\Lambda_{n}^{t}\left(c^{1}\right)_{y}\right]$ and that $\Lambda_{n}^{l}\left(c^{1}\right)_{x}=1$ if $\|x\|=l$.

For $l=0$ is clear since in this case $x=0=y$ and $\Lambda_{n}^{0}\left(c^{1}\right)_{0}=1$. Now let us suppose it true for some $l$, and let us prove it for $l+1$.

Proposition 9. $\Lambda_{n}$ is $k$-expansive for every $k$ odd and it is not 2-expansive if $n \geq 2$.
${ }_{525}$ Proof. Let $c \not \neq k \overline{0}$. Let $D=\{i \mid c(i) \neq 0\}$ and let $D_{l}=D \cap\left(B_{l}(0) \backslash B_{l-1}(0)\right)$. It is clear that $c=\sum_{i \in D} \sigma_{-i}\left(c^{1}\right)$. Since $|D|=k$ is odd, there exists some $l$
such that $\left|D_{l}\right|$ is odd, let us take $\bar{l}$ as the smallest one. For every $x, y \in D_{l}$, $T_{0}\left(\sigma_{-x}\left(c^{1}\right)\right)=T_{0}\left(\sigma_{-y}\left(c^{1}\right)\right)$, thanks to lemma 8. Therfore, given a cell $y \in D_{\bar{l}}$,

$$
\begin{aligned}
\Lambda_{n}^{\bar{l}}(c)_{0} & =\Lambda_{n}^{\bar{l}}\left(\sum_{l \in \mathbb{N}} \sum_{x \in D_{l}} \sigma_{-x}\left(c^{1}\right)\right)_{0} \\
& =\Lambda_{n}^{\bar{l}}\left(\sum_{l=0}^{\bar{l}} \sum_{x \in D_{l}} \sigma_{-x}\left(c^{1}\right)\right)_{0} \\
& =\sum_{l=0}^{\bar{l}} \Lambda_{n}^{\bar{l}}\left(\sum_{x \in D_{l}} \sigma_{-x}\left(c^{1}\right)\right)_{0} \\
& =\Lambda_{n}^{\bar{l}}\left(\sum_{x \in D_{\bar{l}}} \sigma_{-x}\left(c^{1}\right)\right)_{0} \\
& =\Lambda_{n}^{\bar{l}}\left(\sigma_{-y}\left(c^{1}\right)\right)_{0}
\end{aligned}
$$

By lemma 8 this last term is equal to 1 which proves the $k$-expansivity when $k$ is odd.

The second part of the proposition is almost direct from lemma 8. In fact, let $m \in \mathbb{N}$ be any natural number and let us take $z=m s$ for some fixed generator $s$. Now let $s^{\prime}$ be another generator, different from $s$ and $-s$ and define $x=z+s^{\prime}$ and $y=z-s^{\prime}$. This imply that $\|x\|=\|y\|=\|z\|+1=m+1$. Lemma 8 says that $T_{0}\left(\sigma_{x}\left(c^{1}\right)\right)=T_{0}\left(\sigma_{y}\left(c^{1}\right)\right)$, but also that $T_{m}\left(\sigma_{x}\left(c^{1}\right)\right)=T_{m}\left(\sigma_{y}\left(c^{1}\right)\right)$, because $x$ and $y$ are equidistant from $z$, as well as from all the other members of $B_{m}(0)$.
6. Cellular Automata on $\mathbb{Z}^{n}$, with $n \geq 2$

Expansivity is not possible in dimension $n \geq 2$ or more, due to combinatorial reasons: the number of possible $n$-dimensional patterns grows too quickly to be uniformly conveyed into a 1-dimensional array without loss (see [12 for a general result of innexistence of expansive CA).

This argument does not apply to pre-expansivity because the definition allows loss of information in some cases: only finite differences should be propagated, but nothing is required for non-asymptotic pairs of configurations.

Nevertheless, in linear CA over abelian groups, the information propagates in a very regular way, and pre-expansivity is impossible as we will show.
6.1. No pre-expansivity for linear $C A$ in dimension 2 or higher Given a linear CA $F$, spot configurations (i.e. configurations $c$ with $c \neq 1 \overline{0}$ ) form a basis of the whole set of configurations and we get a complete knowledge on the orbit of any configuration from the orbits of spot configurations. The main argument of this section comes from the fact that the space time diagram of spot configurations in a linear CA is substitutive.

Lemma 9 (Main Theorem of $\mathbf{1 9}$ ) Let $F$ be any d-dimensional linear $C A$ on a p-group $Q$. Then there exists a substitution of factor $M$ describing space-time dependency, that is to say, there exists:

- $M \geq 2$,
- a finite set $E$,
- $e: \mathbb{Z}^{d} \times \mathbb{N} \rightarrow E$,
- $\Psi: E \rightarrow(Q \rightarrow Q)$,
- for any $0 \leq a_{1}, . ., a_{d}, s<M$, a function $\Phi_{\vec{a}}^{s}: E \rightarrow E$, such that,

$$
e(M \vec{z}+\vec{a}, M t+s)=\Phi_{\vec{a}}^{s}(e(\vec{z}, t))
$$

- $\Gamma_{\vec{z}}^{t}=\Psi(e(\vec{z}, t))$
where $\Gamma_{z}^{t}$ is the space-time dependency function given by:

$$
\Gamma_{z}^{t}: q \mapsto \sigma_{\vec{z}} \circ F^{t}\left(c^{q}\right) .
$$

Proof. [Proof sketch] We review the proof of [19] written in the one-dimensional setting to show that it works straightforwardly for CA over $\mathbb{Z}^{d}$.

First, as shown by proposition 1 of [19, it is sufficient to consider $Q$ a $p$-group of the form $Q=\mathbb{Z}_{p^{2}}^{D}$ since any linear CA on a $p$-group can be embedded in a linear CA on a group of this form. Then, a linear CA in that case can be viewed as a $D \times D$ matrix whose coefficients are Laurent polynomials with $d$ variables $u_{1}, \ldots, u_{d}$ and coefficients in $\mathbb{Z}_{p^{p}}$.

Formally, we denote by $\mathbb{Z}_{p^{p}}\left[u_{i}, u_{i}^{-1}\right]_{1 \leq i \leq d}$ the Laurent polynomials with variables $u_{1}, \ldots, u_{d}$, i.e. the ring of linear combinations of monomials made with positive or negative powers of the variables. A monomial corresponds to a vector of $\mathbb{Z}^{d}$, hence we use the notation $u^{\vec{i}}$ for any $\vec{i}=\left(i_{1}, \ldots, i_{d}\right) \in \mathbb{Z}^{d}$ to denote the monomial $u_{1}^{i_{1}} \cdots u_{d}^{i_{d}}$. A linear CA is identified with some $T \in \mathcal{M}_{D}\left(\mathbb{Z}_{p^{\imath}}\left[u_{i}, u_{i}^{-1}\right]_{1 \leq i \leq d}\right)$ where the coefficient $a_{\vec{z}} \in \mathbb{Z}_{p^{2}}$ of the monomial $u^{\vec{z}}$ of the coefficient $T_{i, j}$ of the matrix $T$ means intuitively that, when applying the CA, the layer $j$ of cell $\overrightarrow{z_{0}}$ receives $a_{\vec{z}}$ times the content of the layer $i$ of cell $\vec{z}+\overrightarrow{z_{0}}$, and all these individual contributions are summed-up. This matrix representation works well with the powers of $F$ in the sense that $F^{n}$ is represented by $T^{n}$.

We now mimic step by step the proof of 19:

- by the Cayley-Hamilton theorem (Laurent polynomials form an Abelian ring), the characteristic polynomial of $T$ gives a relation of the form:

$$
T^{m}=\sum_{j=0}^{m-1} \sum_{\vec{i} \in I} \lambda_{\vec{i}, j} \vec{u}^{\vec{i}} T^{j}
$$

for some $m$, some finite $I \subseteq \mathbb{Z}^{d}$, and where $\lambda_{\vec{i}, j} \in \mathbb{Z}_{p^{p}}$;

- first, by standard techniques (binomial theorem and binomial coefficients modulo $p$ ), we have the following identity on any commutative algebra of characteristic $p^{l}$ :

$$
\left(\sum_{i} X_{i}\right)^{p^{n+l-1}}=\left(\sum_{i} X_{i}^{p^{n}}\right)^{p^{l-1}}
$$

for any positive $n$. Then, applying this to the expression of $T^{m}$ obtained above we get:

$$
T^{m \cdot p^{n+2(l-1)}}=\left(\sum_{j=0}^{m-1}\left(\sum_{\vec{i} \in I} \lambda_{\vec{i}, j}^{p^{n}} u^{p^{n} \cdot \vec{i}}\right)^{p^{l-1}} T^{p^{n+l-1} \cdot j}\right)^{p^{l-1}}
$$

- noting that the sequence $\left(\lambda_{i, j}^{p^{n}}\right)_{n}$ is ultimately periodic, with common period $N$ for all $\vec{i}, j$, denoting $k=p^{N}$ and expanding the above equality, we have:

$$
T^{k^{n} m^{\prime}}=\sum_{j=0}^{m^{\prime}-1} \sum_{\vec{i} \in I^{\prime}} \mu_{\vec{i}, j} u^{k^{n} \vec{i}} T^{k^{n} j}
$$

for some $m^{\prime}$ large enough and $n \geq 1$, some finite $I^{\prime} \subseteq \mathbb{Z}^{d}$, and where $\mu_{\vec{i}, j} \in \mathbb{Z}_{p^{l}}$;

- we are interested in the space-time dependency function $T_{i}^{j}$ of $T$ which is the matrix of $\mathcal{M}_{D}(Q)$ corresponding to the terms in $u^{\vec{i}}$ of the matrix $T^{j}$, so that $T^{j}=\sum_{\vec{i}} T_{\vec{i}}^{j} u^{\vec{i}}$; we therefore have the following:

$$
T_{\vec{x}}^{k^{n} m^{\prime}+y}=\sum_{\vec{i} \in I^{\prime}} \sum_{j=0}^{m^{\prime}-1} \mu_{\vec{i}, j} T_{\vec{x}-k^{n} \vec{i}}^{y+k^{n}}
$$

- choosing $n=\left\lfloor\log _{k} \frac{y}{m^{\prime}}\right\rfloor$, we can rewrite it for all $y>m^{\prime}$ as:

$$
\begin{equation*}
T_{\vec{x}}^{y}=\sum_{\vec{i}, j} \mu_{\vec{i}, j} T_{\vec{x}+f_{\vec{i}, j}(y)}^{g_{\vec{i}, j}(y)} \tag{1}
\end{equation*}
$$

where $\vec{i}$ and $j$ range over fixed finite sets, and where for any $0 \leq t<k$ we have:

1. $g_{\vec{i}, j}(y)<y$,
2. $\quad f_{\vec{i}, j}(k y+t)=k f_{\vec{i}, j}(y)$,
3. $g_{\vec{i}, j}(k y+t)=k g_{\vec{i}, j}(y)+t$;

- starting from any point $(\vec{x}, y) \in \mathbb{Z}^{d} \times \mathbb{N}$ and applying recursively the equation above as much as we can, we get:

$$
\begin{equation*}
T_{\vec{x}}^{y}=\sum_{\vec{i}, j<m^{\prime}} \alpha_{\vec{i}, j}(\vec{x}, y) T_{\vec{i}}^{j} \tag{2}
\end{equation*}
$$

we take this as a definition of $\alpha_{\vec{i}, j}(x, y)$;

- now, by the properties of the functions $f_{\vec{i}, j}$ and $g_{\vec{i}, j}$ we have that

$$
\alpha_{\vec{i}, j}(\vec{x}, y)=\alpha_{\overrightarrow{0}, j}(\vec{x}-\vec{i}, y),
$$

and in the equation (11) the transformation $(\vec{x}, y) \mapsto(k \vec{x}+\vec{s}, k y+t)$ (for $t<k$ and $\left.\|s\|_{\infty}<k\right)$ applied to the left term $T_{\vec{x}}^{y}$, translates into the same
transformation on each term $T_{\vec{a}}^{b}$ of the right hand side; therefore, denoting $\alpha_{j}=\alpha_{\overrightarrow{0}, j}$, the recursive application of (11) in (2) can be reproduced identically starting from $T_{k \vec{x}+\vec{s}}^{k y+t}$ so that we have

$$
\begin{equation*}
T_{k \vec{x}+\vec{s}}^{k y+t}=\sum_{\overrightarrow{i_{1}}, j<m^{\prime}} \alpha_{j}\left(\vec{x}-\overrightarrow{i_{1}}, y\right) T_{k \overrightarrow{i_{1}}+\vec{s}}^{k j+t} \tag{3}
\end{equation*}
$$

that we want to compare to the final decomposition of $T_{k \vec{x}+\vec{s}}^{k y+t}$ which is by definition of the $\alpha_{j}$ :

$$
\begin{equation*}
T_{k \vec{x}+\vec{s}}^{k y+t}=\sum_{\overrightarrow{i_{2}}, j<m^{\prime}} \alpha_{j}\left(k \vec{x}+\vec{s}-\overrightarrow{i_{2}}, k y+t\right) T_{\overrightarrow{i_{2}}}^{j} \tag{4}
\end{equation*}
$$

- since $k j+t$ is bounded by $k(m+1)$, going from (3) to (4) involves only a bounded number of applications of (1) to each term of the right hand side; therefore, there is some finite set $\Delta \subseteq \mathbb{Z}^{d}$ such that each term $T_{\vec{a}}^{b}$ on the right will produce only terms in $T_{\vec{a}+\vec{i}}^{j}$ for $j<m^{\prime}$ and $\vec{i} \in \Delta$.
Therefore, $\alpha_{j}\left(k \vec{x}+\vec{s}-\overrightarrow{i_{2}}, k y+t\right)$ only depends on the $\alpha_{j}\left(\vec{x}+\frac{\vec{s}-\overrightarrow{i_{2}}}{k}-\overrightarrow{i^{\prime}}, y\right)$ for $\overrightarrow{i^{\prime}} \in \Delta$;
- we choose $X \subseteq \mathbb{Z}^{d}$ large enough so that:

1. $\quad \overrightarrow{i_{2}} \in X$ implies $\overrightarrow{i^{\prime}}-\frac{\vec{s}-\overrightarrow{i_{2}}}{k} \in X$ for any $\overrightarrow{i^{\prime}} \in \Delta$;
2. $\quad T_{\vec{i}}^{j}$ is null for $j<m^{\prime}$ and $\vec{i} \notin X$;
with such an $X$ we finally set $M=k$ and

$$
e(\vec{x}, y)=\left(\alpha_{j}(\vec{x}-\vec{i}, y): j<m^{\prime}, \vec{i} \in X\right)
$$

The lemma follows.
The existence of this substitution has strong consequences on the structure of traces: the trace of a finite configuration is determined by a prefix of linear size in the distance of the farthest non-zero cell. Let us first define some notation.

- For $z \in \mathbb{Z}^{d},\|z\|_{\infty}=\max \left\{\left|z_{i}\right|: i \in\{1, . ., d\}\right.$.
- The size of a configuration $c \stackrel{\infty}{=} \overline{0}$ is the smallest $n \in \mathbb{N}$ such that $c(z) \neq 0$ imply $\|z\|_{\infty} \leq n$.

Lemma 10. Let $F$ be any d-dimensional linear $C A$ on a p-group. Let $m>0$ and denote by $T_{m}$ the trace function associated to $F$ and $m$. There exist a function $\lambda: \mathbb{N} \rightarrow \mathbb{N}$ with $\lambda \in O(n)$ and such that for any $n$ and for any pair of configurations $c_{1}, c_{2}$ with:

- the size of $c_{i}$ is less than or equal to $n$,
- $T_{m}\left(c_{1}\right)(t)=T_{m}\left(c_{2}\right)(t)$ for any $t \leq \lambda(n)$,
then $T_{m}\left(c_{1}\right)=T_{m}\left(c_{2}\right)$.

Proof. First $F$ fullfils the hypothesis of Lemma 9 so we have the existence of the substitution and adopt the notations of the Lemma.

Let's focus on the substitution given by the function $e$ and consider $k \geq 0$, $t \geq M^{k}$, and $\vec{z} \in \mathbb{Z}^{d}$ with $\|z\|_{\infty} \leq M^{k}$. By $k-1$ applications of the substitution ${ }_{640}$ we get the following expression for $e(\vec{z}, t)$ :

$$
e(\vec{z}, t)=\Phi_{\vec{z} \bmod M}^{t \bmod M} \circ \cdots \circ \Phi_{\vec{z} / M^{k-1} \bmod M}^{t / M^{k-1} \bmod M}\left(e\left(\rho(\vec{z}), t / M^{k}\right)\right)
$$

where $\|\rho(\vec{z})\|_{\infty} \leq M$, and where the division/modulus correspond to the standard Euclidean division on $\mathbb{Z}^{d}$.

The sequence of superscripts in this expression only depends on $t \bmod M^{k}$. The sequence of subscripts depends only on $\vec{z}$. Therefore we can write this functional dependency of $e(\vec{z}, t)$ on $e\left(\rho(\vec{z}), t / M^{k}\right)$ in the following way:

$$
\begin{equation*}
e(\vec{z}, t)=\chi_{\vec{z}}^{t \bmod M^{k}}\left(e\left(\rho(\vec{z}), t / M^{k}\right)\right) \tag{5}
\end{equation*}
$$

Now consider a time $t_{0}$ sufficiently large to see any possible vector of the form $\left(e\left(\overrightarrow{z_{0}}, t\right)\right)_{\left\|\overrightarrow{z_{0}}\right\|_{\infty} \leq M}$ before $t_{0}$, precisely:

$$
\forall t, \exists t^{\prime} \leq t_{0}, \forall \overrightarrow{z_{0}},\left\|\overrightarrow{z_{0}}\right\|_{\infty} \leq M: e\left(\overrightarrow{z_{0}}, t\right)=e\left(\overrightarrow{z_{0}}, t^{\prime}\right)
$$

Given an index set $I$, consider any uple $\left(\vec{z}_{i}\right)_{i \in I}$, with $\left\|\vec{z}_{i}\right\|_{\infty} \leq M^{k}$, and any $\mathcal{P} \subseteq E^{I}$. For any time $t$, we can define the property $\mathcal{P}_{I}(t)$ by:

$$
\mathcal{P}_{I}(t) \Leftrightarrow\left(e\left(\vec{z}_{i}, t\right)\right)_{i \in I} \in \mathcal{P} .
$$

Claim: if $\mathcal{P}_{I}(t)$ holds for every $t \leq\left(t_{0}+1\right) \cdot M^{k}$ then $\mathcal{P}_{I}(t)$ holds for every $t \in \mathbb{N}$.
Indeed, take some time $t>\left(t_{0}+1\right) \cdot M^{k}$. Then by choice of $t_{0}$ there exists $t^{\prime} \leq t_{0}$ such that:

$$
\forall \overrightarrow{z_{0}},\left\|\overrightarrow{z_{0}}\right\|_{\infty} \leq M: e\left(\overrightarrow{z_{0}}, t / M^{k}\right)=e\left(\overrightarrow{z_{0}}, t^{\prime}\right)
$$

Now we can choose $t^{\prime \prime} \leq\left(t_{0}+1\right) \cdot M^{k}$ with

$$
\begin{aligned}
t^{\prime \prime} / M^{k} & =t^{\prime} \text { and } \\
t^{\prime \prime} \bmod M^{k} & =t \bmod M^{k}
\end{aligned}
$$

and equation 5 yields the equalities:

$$
e\left(\overrightarrow{z_{i}}, t\right)=\chi_{\vec{z}}^{t} \bmod M^{k}\left(e\left(\rho\left(\overrightarrow{z_{i}}\right), t / M^{k}\right)\right)=\chi_{\vec{z}}^{t^{\prime \prime}} \bmod M^{k}\left(e\left(\rho\left(\overrightarrow{z_{i}}\right), t^{\prime}\right)\right)=e\left(\overrightarrow{z_{i}}, t^{\prime \prime}\right)
$$

It shows that $\mathcal{P}_{I}(t) \Leftrightarrow \mathcal{P}_{I}\left(t^{\prime \prime}\right)$ and the claim follows.
Since the space-time dependency function is completely determined by the substitution $e$ (Lemma (9), the fact that the trace of a finite configuration at time $t$ is null can be expressed by a property of the form $\mathcal{P}_{I}(t)$. More precisely, for any configuration $c$ of size $M^{k}$, we can define

$$
\begin{aligned}
D & =\left\{z \in \mathbb{Z}^{d}: c(z) \neq 0\right\} \\
I & =\bigcup_{z \in D} B_{m}(z) \\
\mathcal{P} & =\left\{f \in E^{I}: \forall x \in B_{m}(0), \sum_{z \in D} \Psi\left(f_{z+x}\right)(c(z))=0\right\}
\end{aligned}
$$

We then have that

$$
\begin{aligned}
\mathcal{P}_{I}(t) & \Leftrightarrow \quad \forall x \in B_{m}(0), \sum_{z \in D} \Psi(e(z+x, t))(c(z))=0 \\
& \Leftrightarrow \quad \forall x \in B_{m}(0), \sum_{z \in D} \Gamma_{z+x}^{t}(c(z))=0 \\
& \Leftrightarrow \sum_{z \in D} T_{m}\left(c^{c(z)}\right)=0 \\
& \Leftrightarrow T_{m}(c)_{t}=0 .
\end{aligned}
$$

Note that this does not avoid a priori the existence of a linear CA which is $k$-expansive for any $k \in \mathbb{N}$ or for infinitely many $k$.
6.2. Linear $C A$ with $Q=\mathbb{Z}_{p}$ In general, in a CA with neighborhood $V \subset \mathbb{G}$, we can remark that the influence of the cell 0 is restricted to the set generated by
We deduce that if $F^{t}(c)$ is null on $B_{m}(0)$ until time $\left(t_{0}+1\right) \cdot M^{k}$ then it is null forever. By linearity of $F$, equality of two traces is equivalent to nullity of their difference. We have thus shown the Lemma for $m \geq 0$ by choosing $\lambda(n)=\left(t_{0}+1\right) \cdot M^{k}$ for $k=\left\lceil\log _{M}(n)\right\rceil$.

Theorem 2. If a CA of dimension $d \geq 2$ is linear over an Abelian group then it is not pre-expansive.

Proof. First, if $G$ is the Abelian group of the Theorem it can be decomposed in a direct product $G=G_{p} \times G^{\prime}$ where $G_{p}$ is a finite $p$-group for some prime $p$ and $G^{\prime}$ is a group whose order $m$ is such that $p$ doesn't divide $m$ (structure theorem for finite abelian groups, see [24]). Then $F$ is isomorphic to $F_{p} \times F^{\prime}$ according to Lemma 1, where $F_{p}$ is linear over $G_{p}$. Moreover if $F$ is pre-expansive, then $F_{p}$ must also be pre-expansive (by Proposition (1). It is therefore sufficient to show the Theorem for p-groups.

Now consider $F$ of dimension $d \geq 2$ linear over a $p$-group and some $m \geq 0$. By lemma 10 we know that the trace $T_{m}$ of finite configurations of size $n$ is determined by its prefix of size $\lambda(n)$ where $\lambda \in O(n)$. The number of such finite configurations grows like $\alpha^{n^{d}}$ for some $\alpha>0$ and the number of prefixes of $T_{m}$ of length $\lambda(n)$ grows like $\beta^{\lambda(n)}$ for some $\beta>0$ which depends only on $m, G_{p}$ and $d$. Since $d \geq 2$ and $\lambda$ is linear we deduce for $n$ large enough that two finite configurations of size $n$ have the same trace $T_{m}$. Therefore $T_{m}$ is not pre-injective and by Proposition 1 , $F$ is not pre-expansive. linear combinations of $-V$. More precisely, at time $t$, its influence is restricted to the following set:

$$
-V_{t}(0)=\left\{\sum_{i=1}^{t} v_{i} \mid(\forall i \in\{1, . ., t\}) v_{i} \in-V\right\}
$$

A perturbation in a cell $u \in \mathbb{G}$ can produce a change in the state of cells in $-V_{t}(u)=u-V_{t}(0)$ up to time $t$.

If $\mathbb{G}$ is commutative, for example $\mathbb{G}=\mathbb{Z}^{n}$ and $V=\left\{v_{1}, . ., v_{m}\right\}$, this set can be computed as follows.

$$
\begin{aligned}
-V_{t}(0) & =\left\{\sum_{k=1}^{m} n_{k}\left(-v_{k}\right) \mid \sum_{k=1}^{m} n_{k}=t \text { and for each } k, n_{k} \in \mathbb{N}\right\} \\
& =\left\{\left.\sum_{k=1}^{m} \frac{n_{k}}{t}\left(-t v_{k}\right) \right\rvert\, \sum_{k=1}^{m} n_{k}=t \text { and for each } k, n_{k} \in \mathbb{N}\right\} \\
& \subseteq\left\{\sum_{k=1}^{m} \lambda_{k}\left(-t v_{k}\right) \mid \sum_{k=1}^{m} \lambda_{k}=1 \text { and for each } k, \lambda_{k} \in[0,1]\right\} \\
& \subseteq c o(-t V)
\end{aligned}
$$

Where $t V=\{t v \mid v \in V\}$, and $c o(\cdot)$ stands for the convex hull (in $\mathbb{R}^{n}$ ).
In the simpler case where $G_{p}=\mathbb{Z}_{p}$, any linear CA $F$ can be expressed as

$$
F=\sum_{z \in V} a_{z} \sigma_{z}
$$

where $\left(a_{z}\right)_{z \in V}$ is a sequence of elements of $\mathbb{Z}_{p}$. When $p$ is prime, we have Lucas' lemma, which gives even stronger properties to linear CAs, more precisely.

$$
\begin{equation*}
F^{p^{k}}=\sum_{z \in V} a_{z} \sigma_{p^{k} z} \tag{6}
\end{equation*}
$$

The next lemma establishes that the constant of $k$-expansivity in a linear CA on $\mathbb{Z}_{p}$ depends only on the radius of the neighborhood. The radius of a neighborhood $V$ is the smallest number $r \in \mathbb{Z}$ such that $V \subseteq B_{r}(0)$.

Lemma 11 (Amplification) Let $F$ be a linear $C A$ with neighborhood $V \subset \mathbb{Z}^{n}$ of radius $r$, with state set $\mathbb{Z}_{p}$. If there exists a configuration $c \neq \overline{0}$ such that $T_{r}(c)=0$, then for any $m \geq r$ there exists a configuration $c^{\prime} \neq \overline{0}$ such that $T_{m}\left(c^{\prime}\right)=0$.

Proof. Let $c$ be such that $T_{r}(c)=0$ and let $k$ be such that $m \leq p^{k}-1$. We define $c^{\prime}$ by $c_{p^{k} x}^{\prime}=c_{x}$ for every $x \in \mathbb{Z}^{n}$ and 0 elsewhere.

From equation 6, it is easy to see that $F^{t p^{k}}\left(c^{\prime}\right)_{p^{k} x}=F^{t}(c)_{x}$, and 0 elsewhere. Therefore, for every $t \in \mathbb{N}$ and every $v \in B_{r}(0), F^{t p^{k}}\left(c^{\prime}\right)_{p^{k} v}=0$.

Now, between iterations $t p^{k}$ and $(t+1) p^{k}$, we know, from the former remarks, that only cells in $\Omega=\bigcup_{x \notin B_{r}(0)}\left(p^{k} x-V_{p^{k}}(0)\right)$ can have a state different from 0 . Since $-V \subseteq B_{r}(0)$, we have that $-V_{p^{k}}(0) \subseteq B_{r p^{k}}(0)$, and the complement of $\Omega$ contains
$B_{p^{k}-1}(0)$, which is what we were looking for, in fact,

$$
\begin{aligned}
y & \in \Omega=\left(\bigcup_{x \notin B_{r}} p^{k} x-V_{p^{k}}(0)\right) \\
& \Rightarrow \quad\left(\exists x \notin B_{r}(0)\right)\left(\exists v \in-V_{p^{k}}(0)\right) y=p^{k} x+v \\
& \Rightarrow\|y\| \geq\left\|p^{k} x\right\|-\|v\| \geq p^{k}(r+1)-p^{k} r=p^{k} \\
& \Rightarrow y \notin B_{p^{k}-1}(0) .
\end{aligned}
$$

Corollary 2. Let $F$ be a linear $C A$ in $\mathbb{Z}_{p}$. It holds:

- $F$ is $k$-expansive, if and only if $F$ is $k$-expansive with expansivity constant $2^{-r}$;
- $F$ is $k$-expansive for all $k \in \mathbb{N}$ if and only if $F$ is pre-expansive.

With this lemma we can establish $k$-expansivity just by looking at $T_{r}$. We will show a CA in that setting which is 1-expansive, 3 -expansive and non 2 -expansive, and another which is non 1-expansive (and so non $k$-expansive for every $k$ ).
6.2.1. The rule $\oplus_{2}$ with von Neumann neighborhood in $\mathbb{Z}^{2}$ The rule that simply sums the state of its 5 neighbors in the von Neumann neighborhood: $(0,0),(0,1),(1,0),(0,-1),(-1,0)$ is not 2-expansive. This can be seen through a simple picture: let us suppose that we start with the configuration $c$ that has a ' 1 ' in cell $\left(-2^{k}, 2^{k-1}\right)$ and in cell $\left(2^{k}, 2^{k-1}\right)$. By symmetry, the vertical line $\{0\} \times \mathbb{Z}$ will be always null. Thus, at iterations $t 2^{k}$ only cells at $2^{k}(\mathbb{Z} \backslash\{0\}) \times \mathbb{Z}$ will be activated. Between iterations $t 2^{k}$ and $(t+1) 2^{k}$ these cells cannot influence the ball $B_{2^{k-1}-1}(0,0)$ (see figure 3) and this ball will have a null trace: $T_{2^{k-1}}(c)=0$.

In order to establish the 3 -expansivity of this CA, we will start by proving some lemmas that describe the form of the traces $T_{1}\left(\sigma_{z}\left(c^{1}\right)\right)$ of the evolution of the configuration $c^{1}$ at the different points of $\mathbb{Z}^{2}$. In order to achieve this, we start by computing the partial traces $\left.T_{0}\left(\sigma_{z}\left(c^{1}\right)\right)\right|_{\left[0,2^{k}-1\right]}$ and $\left.T_{0}\left(\sigma_{z}\left(c^{1}\right)\right)\right|_{\left[2^{k}, 2^{k+1}-1\right]}$. We first give a way for compute them, and afterwards we prove they are effectively the partial traces.

Definition 6. Given $k \geq 0$ and $z \in B_{2^{k}-1}(0,0)$, we recursively define $u_{k}(z)$ and $v_{k}(z)$ as follows. Let us define $S_{k}=\left\{(0,0),\left(0,2^{k}\right),\left(2^{k}, 0\right),\left(0,-2^{k}\right),\left(-2^{k}, 0\right)\right\}$, the


Figure 3. Potentially active cells at iteration $t 2^{k}$ cannot influence $B_{2^{k-1}-1}(0)$ before iteration $(t+1) 2^{k}$. Big dots represent the initially active cells.
active cells of iteration $2^{k}$.
$u_{0}(z)=v_{0}(z)=1 ;$
$u_{k}(z)= \begin{cases}u_{k-1}(z) v_{k-1}(z) & \text { if } z \in B_{2^{k-1}-1}(0,0) \\ 0^{2^{k-1}} u_{k-1}(z-x) & \text { if } z \in B_{2^{k-1}-1}(x) \backslash B_{2^{k-1}-1}(0,0) \text { and } x \in S_{k-1} \\ 0^{2^{k}} & \text { otherwise }\end{cases}$
$v_{k}(z)= \begin{cases}u_{k-1}(z) u_{k-1}(z) & \text { if } z \in B_{2^{k-1}-1}(0,0) \\ u_{k-1}(z-x) u_{k-1}(z-x) & \text { if } z \in B_{2^{k-1}-1}(x) \backslash B_{2^{k-1}-1}(0,0) \text { and } x \in S_{k} \\ 0^{2^{k}} & \text { otherwise }\end{cases}$
Lemma 12. If $z \in B_{2^{k}-1}(0,0)$, then $u_{k}(z)$ and $v_{k}(z)$ represent the trace of $z$ from 0 to $2^{k}-1$ and from $2^{k}$ to $2^{k+1}-1$ respectively.

Proof. When $k=0, B_{0}(0,0)=\{(0,0)\}$, and the trace of $(0,0)$ is constant and equal to 1.

Let us suppose the lemma true for $k-1 \geq 0$. Let $z \in B_{2^{k}-1}(0,0)$.
Figure 4(a) depicts the first two cases in the definition of $u_{k}(z)$, the last one corresponds to cells over the central diagonals.

- Case 1) $z \in B_{2^{k-1}-1}(0,0)$. In this case, the induction hypothesis says that the trace until $2^{k}-1$ is given by $u_{k-1}(z) v_{k-1}(z)$
- Case 2) $z \in B_{2^{k}-1}(0,0) \backslash B_{2^{k-1}-1}(0,0)$. From 0 to $2^{k-1}-1, z$ has not been touched jet, thus its trace until $2^{k-1}-1$ is $0^{2^{k-1}}$. At iteration $2^{k-1}$, only the cells in $S_{k-1}$ are in state 1 , thus $z$ is influenced by only one of the cells in $S_{k-1}$, say $x$, its trace from $2^{k-1}$ to $2^{k}-1$ is equal to the trace of the cell $z-x$ from 0 to $2^{k-1}-1$, thus by induction again, it is equal to $u_{k-1}(z-x)$.
- Case 3) If $z$ belong to none of these balls, it is over one of the two diagonals that pass by $(0,0)$, and its trace is null.


Figure 4. (a) Represents the evolution from iteration 0 to $2^{k}$. Non null cells at iterations 0 and $2^{k-1}$ are marked with brown and red dots, respectively (cell ( 0,0 ) is active at both instants). The balls of radius $2^{k-1}$ around these cells are colored with similar colors. (b) Represents the evolution from iteration $2^{k}$ to $2^{k+1}$. Brown dots represent non null cells at iterations $2^{k}$ and $2^{k}+2^{k-1}$, while red dots represent the cells which are non null at iteration $2^{k}+2^{k-1}$. The balls of radius $2^{k-1}$ around these cells are colored with similar colors. Faded colors represent the cells outside the ball of radius $2^{k}$.

Figure 4 (b) depicts the three cases in the definition of $v_{k}(z)$.

- Case 1) $z \in B_{2^{k-1}}(0,0)$. At iteration $2^{k}$, only the cells in $S_{k}$ are in state 1 . Therefore, from $2^{k}$ to $2^{k}+2^{k-1}-1, z$ will be influenced only by the cell $(0,0)$, and its trace will be equal to $u_{k-1}(z)$. At iteration $2^{k}+2^{k-1}$, the active cells corresponds to red and brown cells in Figure $4(\mathrm{~b})$, and again only cell ( 0,0 ) reaches $z$, the process is repeated.
- Case 2) $z \in B_{2^{k}-1}(0,0) \backslash B_{2^{k-1}-1}(0,0)$. At iteration $2^{k}$, only the cells in $S_{k}$ are in state 1 , thus, before iteration $2^{k}+2^{k-1}, z$ is touched only if it is at distance less than $2^{k-1}$ from one of the cells in $S_{k}$, say $x$, (orange zone in Figure $4(\mathrm{~b})$ ), then its trace is equal to $u_{k-1}(z-x)$. At iteration $2^{k}+2^{k-1}$, the active cells are far again, and $z$ is influenced only by $x$ again.
- Case 3) If $z$ belong to none of these balls, its trace is null.

This lemma proves that the traces can be obtained through the substitution $u \rightarrow u v$ and $v \rightarrow u u$. The basic $u$ and $v$ for a given cell $z$ are obtained at iteration $2^{k+1}$ if $B_{2^{k}-1}(0,0)$ is the smallest ball containing $z$. From the next lemma we can conclude the 1-expansiveness of this automaton with expansivity constant equal to $2^{-1}$.

Lemma 13. If $i+j$ is odd and smaller than $2^{k}$, then the trace of the cell $z=(i, j)$, $T_{0}\left(\sigma_{z}\left(c^{1}\right)\right)$ is not null and its first non null index is odd and smaller than $2^{k}$, in $770 \quad$ particular $u_{k}(i, j)$ is not null.

Proof. For $k=1$, the trace of the odd cells inside $B_{1}(0,0)$ from 0 to 1 is $u=01$, the result holds. Let us assume the result true for $k \geq 1$, and let $z=(i, j)$ be an odd cell in $B_{2^{k+1}-1}(0,0) \backslash B_{2^{k}-1}(0,0)$. Since the cell is odd, it is not over the diagonals, and it belongs to the ball of radius $2^{k}$ of one of the four cells of $S_{k} \backslash\{(0,0)\}$, say $x$. ${ }_{775}$ Then, by lemma 12, its trace from $2^{k}$ to $2^{k+1}-1$ is given by $u_{k}(z-x)$. Since the cells of $S_{k}$ are even, $(z-x)$ is also odd, and the conclusion follows by induction.

Now we will prove several properties that will be useful to prove 3-expansivity.
Lemma 14. Given $k>0$, if $z \in B_{2^{k}-1}(0,0)$ then $v_{k}(z)$ is a square.
Proof. It is clear from Definition 6 and the fact that $k>0$.
Lemma 15. Given $k>0$, if $z \in B_{2^{k}-1}(0,0) \backslash\{(0,0)\}$ and $u_{k}(z) \neq 0^{2^{k}}$, then $u_{k}(z)$ is not a square.

Proof. By induction on $k$. For $k=1$, if $z$ is inside the von Neumann neighborhood of $(0,0)$, thus $u_{1}(z)=01$ which is not a square. Let us suppose the assertion true for $k-1 \geq 1$. Since $k \geq 1$, from definition 6, we recognize two cases for $u_{k}(z)$.

Case 1: $u_{k}(z)=0^{2^{k}} u_{k-1}(z-x)$, for some $x \in S_{k-1}$. The only way for $u_{k}(z)$ to be a square is to be equal to $0^{2^{k}}$.

Case 2: $u_{k}(z)=u_{k-1}(z) v_{k-1}(z)$. By the induction hypothesis $u$ is either null or not a square. From lemma $14 v_{k-1}(z)$ is a square, then $u_{k}(z)$ is a square if and only if $u_{k}(z)=0^{2^{k}}$.

Lemma 16. If $|i|=|j|$ and $(i, j) \neq 0$, then the trace of the cell $(i, j), T_{0}\left(\sigma_{(i, j)}\left(c^{1}\right)\right)$ is equal to $0^{\omega}$.

Proof. Cells in the diagonals systematically falls in the boundaries of the zones given by the substitution, then their traces are systematically assigned equal to 0 .

Lemma 17. If $i+j$ is even, smaller than $2^{k}$ and the trace of the cell $z=(i, j)$, $T_{0}\left(\sigma_{z}\left(c^{1}\right)\right)$ is not null, then the first non null index of the trace is even and it is smaller than $2^{k}$.

Proof. For $k=0$, the trace of the even cell inside $B_{0}(0,0)$ from 0 to 0 is $u=1$, the result holds. Let us assume the result true for $k \geq 0$, and let $z=(i, j)$ be an even cell in $B_{2^{k+1}-1}(0,0) \backslash B_{2^{k}-1}(0,0)$. Since the cell is attained, from Lemma 16, it is not over the diagonals, then it belongs to the ball of radius $2^{k}$ of one of the four cells in $S_{k}$, say $x$. Then, its trace from $2^{k}$ to $2^{k+1}-1$ is given by $u_{k}(z-x)$. Since
the cells in $S_{k}$ are even, $z-x$ is also even, and the conclusion follows by induction.

Now we are ready to prove that this automaton is 3-expansive.
Lemma 18. If $z_{1}, z_{2}$ and $z_{3}$ are three different cells and $T_{0}\left(\sigma_{z_{1}}\left(c^{1}\right)\right)+T_{0}\left(\sigma_{z_{2}}\left(c^{1}\right)\right)+$ $T_{0}\left(\sigma_{z_{3}}\left(c^{1}\right)\right)=0^{\omega}$, then there exists $z \in\left\{z_{1}, z_{2}, z_{3}\right\}$ such that $T_{0}\left(\sigma_{z}\left(c^{1}\right)\right)=0^{\omega}$.
Proof. We will prove a stronger assertion:
If $z_{1}, z_{2}$ and $z_{3}$ are three different cells in $B_{2^{k}-1}(0,0)$ and $u_{k}\left(z_{1}\right)+$ $u_{k}\left(z_{2}\right)+u_{k}\left(z_{3}\right)=0^{2^{k}}$, then there exists $z \in\left\{z_{1}, z_{2}, z_{3}\right\}$ such that $u_{k}(z)=0^{2^{k}}$.

It is stronger because from Lemmas 13 and 17, if $z \in B_{2^{k}-1}(0,0)$ and $u_{k}(z)=0^{2^{k}}$, then $T_{0}\left(\sigma_{z}\left(c^{1}\right)\right)=0^{\omega}$.

By contradiction, let $z_{1}, z_{2}$ and $z_{3}$ be three different cells in a ball $B_{2^{k}-1}(0,0)$ such that $u_{k}\left(z_{1}\right)+u_{k}\left(z_{2}\right)+u_{k}\left(z_{3}\right)=0^{2^{k}}$ and for all $i \in\{1,2,3\}, u_{k}\left(z_{i}\right) \neq 0^{2^{k}}$. Let us choose these cells such that $k$ is as small as possible.

Let be $t_{z_{1}}, t_{z_{2}}$ and $t_{z_{3}}$ the indices where the respective traces equals 1 by the first time. It is clear that two of these numbers are equal and smaller than the third. Let us suppose that $t_{z_{1}}=t_{z_{2}}<t_{z_{3}}$.

Case 1: $z_{1}, z_{2} \in B_{2^{k-1}-1}(0,0)$. In this case, the trace of $z_{3}$ from $2^{k-1}$ to $2^{k}-1$ is equal to $u_{k}\left(z_{3}-x\right)$, for some $x \in S_{k}$, thus $u_{k}\left(z_{3}\right)=0^{2^{k-1}} u_{k-1}\left(z_{3}-x\right)$. On the other hand, $u_{k}\left(z_{i}\right)=u_{k-1}\left(z_{i}\right) v_{k-1}\left(z_{i}\right)$ for $i \in\{1,2\}$. From Lemma 14, $v_{k-1}\left(z_{1}\right)$ and $v_{k-1}\left(z_{2}\right)$ are squares if $k>0$, which is the case, since $B_{0}(0,0)$ contains only one cell. Thus $v_{k-1}\left(z_{1}\right)+v_{k-1}\left(z_{2}\right)$ is also a square, then it cannot be equal to $u_{k-1}\left(z_{3}-x\right)$ which is not a square thanks to Lemma 15

Case 2: $z_{1}, z_{2}, z_{3} \notin B_{2^{k-1}-1}(0,0)$. In this case, for each $i \in\{1,2,3\}$ there exists some $x_{i} \in S_{k-1}$ such that $u_{k}\left(z_{i}\right)=0^{2^{k-1}} u_{k-1}\left(z_{i}-x_{i}\right)$. Of course $z_{i}-x_{i} \in$ $B_{2^{k-1}-1}(0,0)$ for each $i \in\{1,2,3\}$ and $u_{k-1}\left(z_{1}-x_{1}\right)+u_{k-1}\left(z_{2}-x_{2}\right)+$ $u_{k-1}\left(z_{3}-x_{3}\right)=0^{2^{k-1}}$, which contradicts the minimality of $k$.

With this last lemma we know that if three cells produces a null trace of radius 1 , thus one of them has a null trace of radius 0 , this means that this cell is even, and its four neighbors are odd. When looking at the neighbors of these cells, their sum is also null, for each if its neighbors. Thus at least one cell must have even neighbors with a null trace, but in this case, two of these cells are odd, and its four even neighbors cannot equal the odd neighbors of the first cell, and then the trace of radius 1 of the sum of the three cells cannot be null.
6.2.2. The rule $\oplus_{2}$ with triangular neighborhood The last rule is 1 and 3 expansive, now we present a linear rule which is even not 1-expansive. Thanks to Proposition 3, it implies in particular that is not $k$-expansive, for every $k \in \mathbb{N}$.


Figure 5. Simulation of $\oplus_{2}$ with a triangular neighborhood at iteration 25 starting from an initial configuration with a single 1 inside 0s: cells in state 1 are in red, cells that have been in state 1 between iteration 0 and 24 but not at 25 are in blue, and the others cells are not drawn.

It correspond to addition modulo 2 as the last one but with a triangle shaped neighborhood: $N=\{(-1,1),(1,1),(0,0),(0,-1)\}$.

Proposition 10. $T_{2}\left(\sigma_{(0,36)}\left(c^{1}\right)\right)$ is null.

Proof. We will prove, by induction, that $T_{2}\left(\sigma_{(0,36)}\left(c^{1}\right)\right)$ is null from 0 to $2^{k}$. For $k=0$ to 5 is clear since cell $(0,36)$ is too far to touch $B_{2}(0,0)$. At iteration $2^{k}$ the only active cells are $\left(-2^{k},-2^{k}+36\right),\left(2^{k},-2^{k}+36\right),\left(0,2^{k}+36\right)$ and $(0,36)$. The first three are too far to touch $B_{2}(0,0)$ from iterations $2^{k}$ to $2^{k+1}$. By induction hypothesis $(0,36)$ does not attain $B_{2}(0,0)$ in $2^{k}$ iterations, thus $B_{2}(0,0)$ remains null until iteration $2^{k+1}$.

## 7. Open Problems

We showed in this paper that dynamics like pre-expansivity or $k$-expansivity can exist without necessarily implying positive expansivity. We also showed that some combinations of the space structure and the local rule structure forbid preexpansivity (Theorem (2).

However, we left many open questions concerning pre-expansivity and $k$ expansivity:

- is there a pre-expansive cellular automata on $\mathbb{Z}^{d}$ when $d \geq 2$ ?
- is there a 2-expansive cellular automata on $\mathbb{Z}^{2}$ ? on the free group? more generally which is the set of integers $k$ such that a given group admits $k$ expansive cellular automata?
- are pre-expansive CA always transitive? mixing? open?
- is the property of pre-expansivity decidable?


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