# On Dynamical Behavior of the $p$-adic $\lambda$-Ising Model on Cayley Tree 

Mutlay Dogan

In the present paper, we continue to study some features of the mixed type $p$-adic $\lambda$-Ising model which was studied in [3]. In that study, the existence of the $p$-adic Gibbs measures and phase transitions were investigated for the model on the Cayley tree of order two. In the current paper, we study the dynamical behavior of the fixed points which have been found in [3]. As the main result, we proved that the fixed point $u_{0}$ is an attractor and the other fixed points $u_{1,2}$ are repellent fixed points for the mixed type $p$-adic $\lambda$-Ising model. In addition, the size of basin of attractor for the fixed point $u_{0}$ is described.

Key words: $p$-adic numbers, $p$-adic quasi Gibbs measure, dynamical systems, Cayley tree.

Mathematical Subject Classification 2010: 58F12, 46S10, 12J12, 11S99, $54 \mathrm{H} 20,30 \mathrm{D} 05$.

## 1. Introduction

The $p$-adic numbers were firstly described by German mathematician K. Hensel, and they have been attracting interest of scientists since then. By now, many theoretical and practical applications have being studied in the $p$ adic field. Also some features of dynamical behavior of a dynamical system were studied in $[2,7,10,13,17]$.

A number of scientists applied renormalization techniques in statistical mechanics, and they got lots of interesting results. As consequences of such results about phase transitions of spin models on hierarchical lattices represented that they make the exact calculation of multifarious physical quantities $[1,5]$. One of these useful hierarchical lattices is a Cayley tree or a Bethe lattice (see [19]). The lattice is not real but it provides to do some certain calculations of many physical quantities. And it helps to construct corresponding dynamical systems for many complicated models [21].

Renormalization methods were widely applied to study the Ising model [1]. At the same time, one of the generalizations of the Ising model is the so-called $\lambda$-model on the Cayley tree (see [11,20]). In [3], we combined these two most studied models and called " $p$-adic $\lambda$-Ising model". In the same work we proved the

[^0]existence and uniqueness of this kind of $p$-adic quasi Gibbs measures. Moreover, we proved the occurrence of the phase transition according to the $p$-adic $\lambda$-Ising model.

Currently, we continue studying different aspects of this model. In [3], we obtained a dynamical function and got three fixed points $u_{0}, u_{1,2}$. In the present paper, we will explore dynamical behavior of these fixed points.

As the main result of this paper, we have proved that $u_{0}$ is an attractor and $u_{1,2}$ are repellent fixed points for the mixed type $p$-adic $\lambda$-Ising model on the Cayley tree of order two.

The result is obtained for the $p$-adic case. However, it is not valid for the real case.

## 2. Preliminaries

2.1. $p$-adic numbers. Suppose that $p$ is a fixed prime number. The set $\mathbb{Q}_{p}$ is defined as a completion of the rational numbers $\mathbb{Q}$ with the norm $|\cdot|_{p}: \mathbb{Q} \rightarrow$ $\mathbb{R}$ given by

$$
|x|_{p}= \begin{cases}p^{-r}, & x \neq 0 \\ 0, & x=0\end{cases}
$$

where $x=p^{r} \frac{m}{n}$ with $r, m \in \mathbb{Z}, n \in \mathbb{N},(m, p)=(n, p)=1$. The absolute value $|\cdot|_{p}$ is called non-Archimedean norm, and it satisfies the strong triangle inequality

$$
|x+y|_{p} \leq \max \left\{|x|_{p},|y|_{p}\right\}
$$

This is the most crucial property of the norm, i.e., if $|x|_{p}>|y|_{p}$, then $\mid x+$ $\left.y\right|_{p}=|x|_{p}$. Notice that this very useful feature can be employed only in the non-Archimedean norm.

Any $p$-adic number $x \in \mathbb{Q}_{p}, x \neq 0$ is uniquely represented in the form

$$
\begin{equation*}
x=p^{\gamma(x)}\left(x_{0}+x_{1} p+x_{2} p^{2}+\cdots\right) \tag{2.1}
\end{equation*}
$$

where $\gamma=\gamma(x) \in \mathbb{Z}$ and $x_{j}$ are integers, $0 \leq x_{j} \leq p-1, x_{0}>0, j=0,1,2, \ldots$, in the case $|x|_{p}=p^{-\gamma(x)}$.

We recall that the $p$-adic integers set is

$$
\mathbb{Z}_{p}=\left\{x \in \mathbb{Q}_{p}:|x|_{p} \leq 1\right\}
$$

And a $p$-adic exponential function is defined by

$$
\exp _{p}(x)=\sum_{n=1}^{\infty} \frac{x^{n}}{n!}
$$

and is convergent for every $x \in B\left(0, p^{-1 /(p-1)}\right)$. It is known [9] that for any $x \in$ $B\left(0, p^{-1 /(p-1)}\right)$ one has

$$
\left|\exp _{p}(x)\right|_{p}=1, \quad\left|\exp _{p}(x)-1\right|_{p}=|x|_{p}<1
$$

Let

$$
\begin{equation*}
\mathcal{E}_{p}=\left\{x \in \mathbb{Q}_{p}:|x-1|_{p}<p^{-1 /(p-1)}\right\} \tag{2.2}
\end{equation*}
$$

Lemma 2.1 ([18]). The set $\mathcal{E}_{p}$ has the following properties:

1. $\mathcal{E}_{p}$ is a group under multiplication.
2. $|a-b|_{p}<1$ for all $a, b \in \mathcal{E}_{p}$.
3. If $a, b \in \mathcal{E}_{p}$, then it holds

$$
|a+b|_{p}= \begin{cases}\frac{1}{2} & \text { if } p=2 \\ 1 & \text { if } p \neq 2\end{cases}
$$

4. If $a \in \mathcal{E}_{p}$, then there is an element $h \in B\left(0, p^{-1 /(p-1)}\right)$ such that $a=\exp _{p}(h)$.
5. Let $x, y \in \mathbb{Q}_{p}$. If $|x y|_{p}=1$ and $|x+y|_{p}<1$, then $|x|_{p}=|y|_{p}=1$.
2.2. $p$-adic measure. Assume that $(X, \mathcal{B})$ is a measurable space, where $\mathcal{B}$ is an algebra of subsets of $X$. A function $\mu: \mathcal{B} \rightarrow \mathbb{Q}_{p}$ is called a $p$-adic measure if the equality

$$
\mu\left(\bigcup_{j=1}^{n} A_{j}\right)=\sum_{j=1}^{n} \mu\left(A_{j}\right)
$$

holds for any $A_{1}, \ldots, A_{n} \subset \mathcal{B}$ such that $A_{i} \cap A_{j}=\emptyset(i \neq j)$.
A $p$-adic measure is called a probability measure whenever $\mu(X)=1$.
2.3. Cayley tree. Let $\Gamma_{+}^{k}=(V, L)$ be a semi-infinite Cayley tree of order $k \geq 1$ with the root $x^{(0)}$. Here, the set of all vertices is $V$ and the set of all edges is $L$. A collection of the pairs $\left\langle x, x_{1}\right\rangle, \ldots,\left\langle x_{d-1}, y\right\rangle$ is called a path from the vertex $x$ to the vertex $y$. The distance $d(x, y), x, y \in V$, on the Cayley tree is the length of the shortest path from $x$ to $y$. The vertices $x$ and $y$ are the nearest neighbors denoted by $l=\langle x, y\rangle$ if $d(x, y)=1$. Two vertices $x, y \in V$ are considered as the next-nearest neighbors if $d(x, y)=2$. The next-nearest-neighbors vertices $x$ and $y$ are the prolonged next-nearest neighbors if $x \in W_{n-2}$ and $y \in W_{n}$ for some $n \geq 1$, which are shown by $\rangle x, y\langle$, and called one-level next-nearest-neighbors if $x, y \in W_{n}$ for some $n$ and shown by $\rangle \overline{x, y}\langle$.

$$
\begin{aligned}
W_{n} & =\left\{x \in V: d\left(x, x^{(0)}\right)=n\right\}, \quad V_{n}=\bigcup_{m=0}^{n} W_{m} \\
L_{n} & =\left\{l=\langle x, y\rangle \in L: x, y \in V_{n}\right\} .
\end{aligned}
$$

The direct successors set of $x$ is defined by

$$
S(x)=\left\{y \in W_{n+1}: d(x, y)=1\right\}, \quad x \in W_{n}
$$

2.4. Dynamical systems in $\mathbb{Q}_{p}$. In this subsection, we remind some standard notions of dynamical systems. Let $r, s>0$ s.t. $r<s$, and $a \in \mathbb{Q}_{p}$. Then we define the $p$-adic balls and a $p$-adic sphere as follows:

$$
\begin{array}{ll}
B_{r}(a)=\left\{x \in \mathbb{Q}_{p}:|x-a|_{p}<r\right\}, & \bar{B}_{r}(a)=\left\{x \in \mathbb{Q}_{p}:|x-a|_{p} \leq r\right\} \\
B_{r}(a)=\left\{x \in \mathbb{Q}_{p}: r<|x-a|_{p}<s\right\}, & S_{r}(a)=\left\{x \in \mathbb{Q}_{p}:|x-a|_{p}=r\right\} \tag{2.4}
\end{array}
$$

It is obvious that $\bar{B}_{r}(a)=B_{r}(a) \cup S_{r}(a)$.
If any function $f: B_{r}(a) \rightarrow \mathbb{Q}_{p}$, which converges uniformly on the ball $B_{r}(a)$, can be represented by

$$
f(x)=\sum_{n=0}^{\infty} \frac{f^{(n)}(a)}{n!}(x-a)^{n}, \quad f \in \mathbb{Q}_{p}
$$

then it is said to be analytic.
Letting a dynamical system $(f, B)$ in $\mathbb{Q}_{p}$, where $f: x \in B \rightarrow f(x) \in B$, be an analytic function and $B=B_{r}(a)$ or $\mathbb{Q}_{p}$. Denote $x^{(n)}=f^{n}\left(x^{(0)}\right)$, where $x^{0} \in B$ and $f^{n}(x)=\underbrace{f \circ \cdots \circ f(x)}_{n}$. If $f\left(x^{(0)}\right)=x^{(0)}$, then $x^{(0)}$ is said to be a fixed point.

A fixed point $x^{(0)}$ is called an attractor if there exists a neighborhood $U\left(x^{(0)}\right)(\subset B)$ of $x^{(0)}$ such that for all points $y \in U\left(x^{(0)}\right)$ it holds $\lim _{n \rightarrow \infty} y^{n}=$ $x^{(0)}$ where $y^{n}=f^{n}(y)$. If $x^{(0)}$ is an attractor, then its basin of attraction is

$$
A\left(x^{(0)}\right)=\left\{y \in \mathbb{Q}_{p}: y^{n} \rightarrow x^{(0)}, \text { as } n \rightarrow \infty\right\}
$$

A fixed point $x^{(0)}$ is said to be a repellent if there exists a neighborhood $U\left(x^{(0)}\right)$ of $x^{(0)}$ such that $\left|f(x)-x^{(0)}\right|_{p}>\left|x-x^{(0)}\right|_{p}$ for $x \in U\left(x^{(0)}\right), x \neq x^{(0)}$.

For a fixed point $x^{(0)}$ of a function $f(x)$ a ball $B_{r}\left(x^{(0)}\right)$ (contained in $B$ ) is said to be a Siegel disc if each sphere $S_{\rho}\left(x^{(0)}\right), \rho<r$, is an invariant sphere of $f(x)$, i.e., if $x \in S_{\rho}\left(x^{(0)}\right)$, then all iterated points $x^{(n)} \in S_{\rho}\left(x^{(0)}\right)$ for all $n=1,2 \ldots$ The union of all Siegel discs with the center at $x^{(0)}$ is said to be a maximum Siegel disc and it is denoted by $S I\left(x^{(0)}\right)$.

In other words, let $x^{(0)}$ be a fixed point of an analytic function $f(x)$. Then

$$
\lambda=\frac{d}{d x} f\left(x^{(0)}\right)
$$

It is clear that $\lambda$ is a usual derivative of $f$ at $x^{(0)}$. Therefore, the point $x^{(0)}$ is attractive whenever $0 \leq|\lambda|_{p}<1$, neutral whenever $|\lambda|_{p}=1$, and repellent whenever $|\lambda|_{p}>1$.

## 3. Construction of $p$-adic Gibbs measure and dynamical function

In this section, we deal with the mixed type $p$-adic $\lambda$-Ising model where spins take values $\{ \pm 1\}$, and these values are assigned to the vertices of the Cayley tree $\Gamma_{+}^{k}=(V, L)$. A configuration $\sigma$ on $V$ is defined as a function of $\sigma: x \in V \rightarrow$
$\sigma(x) \in \Phi$; in a similar manner one defines configurations $\sigma_{n}$ and $\omega$ on $V_{n}$ and $W_{n}$, respectively. The set of all configurations on $V$ (respectively $V_{n}, W_{n}$ ) is the same with $\Omega=\Phi^{V}$ (respectively $\Omega_{V_{n}}=\Phi^{V_{n}}, \Omega_{W_{n}}=\Phi^{W_{n}}$ ) and it is easy to see that $\Omega_{V_{n}}=\Omega_{V_{n-1}} \times \Omega_{W_{n}}$. Due to the given configurations $\sigma_{n-1} \in \Omega_{V_{n-1}}$ and $\omega \in$ $\Omega_{W_{n}}$, their concatenation is defined by

$$
\left(\sigma_{n-1} \vee \omega\right)(x)= \begin{cases}\sigma_{n-1}(x), & \text { whenever } x \in V_{n-1} \\ \omega(x), & \text { whenever } x \in W_{n}\end{cases}
$$

It is obvious that $\sigma_{n-1} \vee \omega \in \Omega_{V_{n}}$.
$\lambda: \Phi \times \Phi \rightarrow \mathbb{Q}_{p}$ is a function defined on each edge $\langle x, y\rangle \in L$. However, the Hamiltonian of the $p$-adic $\lambda$-Ising model can be defined by

$$
\begin{equation*}
H_{n}(\sigma)=\sum_{\langle x, y\rangle \in L_{n}} \lambda(\sigma(x), \sigma(y))+J \sum_{\langle x, y\langle } \sigma(x) \sigma(y), \tag{3.1}
\end{equation*}
$$

where $H_{n}: \Omega_{V_{n}} \rightarrow \mathbb{Q}_{p}$.
In what follows, we suppose that $|\lambda(x, y)|_{p} \leq 1 / p$ and $|J|_{p} \leq 1 / p(p \geq 3)$. These conditions guarantee the existence of $\exp _{p}\left(H_{n}(\sigma)\right)$ for all $n \in \mathbb{N}$ and $\sigma . \mathbf{h}$ : $x \in V \backslash\left\{x^{(0)}\right\} \rightarrow h_{x} \in \mathbb{Q}_{p}$ is a mapping. For any $n \in \mathbb{N}$, a $p$-adic probability measure $\mu_{\mathbf{h}}^{(n)}$ on the entire configuration $\Omega_{V_{n}}$ is

$$
\begin{equation*}
\mu_{\mathbf{h}}^{(n)}(\sigma)=\frac{1}{Z_{n, h}^{(\mathbf{h})}} \exp _{p}\left(H_{n}(\sigma)\right) \prod_{x \in W_{n}}\left(h_{x}\right)^{\sigma(x)} \tag{3.2}
\end{equation*}
$$

Here, $\sigma \in \Omega_{V_{n}}$, and $Z_{n, h}^{(\mathbf{h})}$ is the related function called a partition function and given by

$$
\begin{equation*}
Z_{n, h}^{(\mathbf{h})}=\sum_{\sigma \in \Omega_{V_{n}}} \exp _{p}\left(H_{n}(\sigma)\right) \prod_{x \in W_{n}}\left(h_{x}\right)^{\sigma(x)} \tag{3.3}
\end{equation*}
$$

Remember [12] that one of the main results of the theory of probability concerns a construction of an infinite volume distribution with given finitedimensional distributions, which is a well-known Kolmogorov's Extension Theorem [23]. Recall that a p-adic probability measure $\mu$ on $\Omega$ is compatible if one holds

$$
\begin{equation*}
\mu\left(\sigma \in \Omega:\left.\sigma\right|_{V_{n}}=\sigma_{n}\right)=\mu_{\mathbf{h}}^{(n)}\left(\sigma_{n}\right) \quad \text { for all } \sigma_{n} \in \Omega_{V_{n}}, n \in \mathbb{N} \tag{3.4}
\end{equation*}
$$

The existence of the measure $\mu$ is guaranteed by Kolmogorov's theorem $[4,8]$. Namely, if the measures $\mu_{\mathrm{h}}^{(n)}, n \geq 1$ satisfy the compatibility condition, i.e.,

$$
\begin{equation*}
\sum_{\omega \in \Omega_{W_{n}}} \mu_{\mathbf{h}}^{(n)}\left(\sigma_{n-1} \vee \omega\right)=\mu_{\mathbf{h}}^{(n-1)}\left(\sigma_{n-1}\right) \tag{3.5}
\end{equation*}
$$

for any $\sigma_{n-1} \in \Omega_{V_{n-1}}$, then there is a unique $p$-adic Gibbs measure $\mu$ on $\Omega$ with (3.4).

Now, following [12], if the measures $\mu_{\mathbf{h}}^{(n)}$ satisfy the compatibility condition for a function $\mathbf{h}$, then there is a unique $p$-adic probability measure, which is denoted
by $\mu_{\mathbf{h}}$. The measure $\mu_{\mathbf{h}}$ is said to be a $p$-adic quasi Gibbs measure corresponding to the $p$-adic $\lambda$-Ising model. By $Q \mathcal{G}(H)$, we denote the set of all $p$-adic quasi Gibbs measures associated with functions $\mathbf{h}=\left\{\mathbf{h}_{x}, x \in V\right\}$.

Theorem 3.1. The measure $\mu_{h}^{(n)}\left(\sigma_{n}\right), n=1,2, \ldots$, satisfies Kolmogorov's consistency condition (3.5) if and only if the following equations hold for any $x \in V$ :

$$
\begin{equation*}
h_{x}=\sum_{y \in S(x)} F_{x, y}\left(h_{y}, \lambda, J\right) \tag{3.6}
\end{equation*}
$$

and

$$
\begin{equation*}
h_{x}^{2}=\frac{a^{2} \theta^{2} h_{y}^{2} h_{z}^{2}+a b\left(h_{y}^{2}+h_{z}^{2}\right)+b^{2} \theta^{2}}{c^{2} \theta^{2} h_{y}^{2} h_{z}^{2}+c d\left(h_{y}^{2}+h_{z}^{2}\right)+d^{2} \theta^{2}} \tag{3.7}
\end{equation*}
$$

where $a=\exp _{p}(\lambda(1,1)), \quad b=\exp _{p}(\lambda(1,-1)), \quad c=\exp _{p}(\lambda(-1,1)), \quad d=$ $\exp _{p}(\lambda(-1,-1)), \theta=\exp _{p}(J)$.

Proof. Necessity. Assume that (3.5) holds; we want to obtain (3.7). We substitute (3.2) in (3.5) and hence, for any configuration $\sigma_{n-1} \in \Omega_{V_{n}}$, we get

$$
\begin{array}{r}
Z_{n}^{-1} \sum_{\sigma^{(n)}} \exp _{p}\left[\sum_{\langle x, y\rangle \in L_{n}} \lambda_{x, y}(\sigma(x), \sigma(y))+J \sum_{\substack{\rangle x, y \backslash \\
x, y \in W_{n}}} \sigma(x) \sigma(y)+\sum_{x \in W_{n}} h_{x} \sigma(x)\right] \\
=Z_{n-1}^{-1} \sum_{\sigma^{(n)}} \exp _{p}\left[\sum_{\substack{\langle x, y\rangle \in L_{n-1}}} \lambda_{x, y}(\sigma(x), \sigma(y))\right. \\
\left.+J \sum_{\substack{\backslash x, y \backslash \\
x, y \in W_{n-1}}} \sigma(x) \sigma(y)+\sum_{x \in W_{n-1}} h_{x} \sigma(x)\right]
\end{array}
$$

It yields that

$$
\begin{aligned}
\exp _{p}\left(H_{n}\left(\sigma_{n-1}\right)\right) & \prod_{x \in W_{n-2}} \prod_{y \in S(x)}\left(h_{x y, \sigma(x) \sigma(y)}\right)^{\sigma(x) \sigma(y)} \\
= & L_{n} \sum_{\eta \in \Omega_{W_{n}}} \exp _{p}\left[H_{n}\left(\sigma_{n-1}\right)+\sum_{x \in W_{n-1}} \sum_{y \in S(x)} \lambda_{x, y}(\sigma(x), \sigma(y))\right. \\
& \left.+J \sum_{\substack{y, z \in S(x) \\
y \neq z}} \sigma(y) \sigma(z)\right] \prod_{x \in W_{n-1}} \prod_{y \in S(x)}\left(h_{x y, \sigma(x) \eta(y)}\right)^{\sigma(x) \eta(y)}
\end{aligned}
$$

From here we get

$$
\frac{Z_{n-1}}{Z_{n}} \sum_{\sigma(n)} \exp _{p}\left[\sum_{x \in W_{n-1}} \sum_{y \in S(x)} \lambda_{x, y}(\sigma(x), \sigma(y))+J \sum_{\substack{y, z \in S(x) \\ y \neq z}} \sigma(y) \sigma(z)\right.
$$

$$
\left.+\sum_{y \in S(x)} h_{y} \sigma(y)\right]=\prod_{x \in W_{n-1}} \exp _{p}\left(h_{x} \sigma(x)\right)
$$

Then, we find

$$
\begin{aligned}
& \frac{Z_{n-1}}{Z_{n}} \prod_{x \in W_{n-1}} \prod_{y \in S(x)} \sum_{\sigma(y) \in \pm 1} \exp _{p}\left[\sum_{x \in W_{n-1}} \lambda_{x, y}(\sigma(x), \sigma(y))\right. \\
& \left.+J \sum_{\substack{y, z \in S(x) \\
y \neq z}} \sigma(y) \sigma(z)+h_{y} \sigma(y)\right]=\prod_{x \in W_{n-1}} \exp _{p}\left(h_{x} \sigma(x)\right)
\end{aligned}
$$

In the equation above, if we substitute $\sigma(x)=1$ and $\sigma(x)=-1$ respectively, then, after dividing the corresponding equations by each other, we get equation (3.7), which means that (3.6) holds.

Sufficiency. Assume that (3.6) is valid. Then it implies the existence of $a(x) \in$ $\mathbb{Q}_{p}$ such that

$$
\begin{align*}
\prod_{y \in S(x)} \sum_{\tilde{\sigma}(y) \in\{ \pm 1\}} & \exp _{p}[\lambda(\sigma(x), \tilde{\sigma}(y)) \\
& \left.+J \sum_{\substack{z \in S(x) \\
z \neq y \\
\eta(z) \in \Phi}} \tilde{\sigma}(y) \eta(z)+\tilde{\sigma}(y) h_{y}\right]=a(x) \exp _{p}\left(\sigma(x) h_{x}\right) \tag{3.8}
\end{align*}
$$

where $\sigma(x) \in\{ \pm 1\}$. From the last equality, one gets

$$
\begin{align*}
& \prod_{x \in W_{n-1}} \prod_{y \in S(x)} \sum_{\tilde{\sigma}(y) \in\{ \pm 1\}} \exp _{p}[\lambda(\sigma(x), \tilde{\sigma}(y)) \\
& \left.\quad+J \sum_{\substack{z \in S(x) \\
z \neq y \\
\eta(z) \in \Phi}} \tilde{\sigma}(y) \eta(z)+\tilde{\sigma}(y) h_{y}\right]=\prod_{x \in W_{n-1}} a(x) \exp _{p}\left(\sigma(x) h_{x}\right) \tag{3.9}
\end{align*}
$$

Now, multiplying both sides of (3.9) by $\exp _{p}\left(H_{n-1}(\sigma)\right)$ and denoting

$$
\begin{equation*}
A_{n}(x)=\prod_{x \in W_{n}} a(x) \tag{3.10}
\end{equation*}
$$

we get

$$
\begin{aligned}
& U_{n} \exp _{p}\left(H_{n-1}(\sigma)\right) \prod_{x \in W_{n-2}} \prod_{y \in S(x)}\left(h_{x y, \sigma(x) \sigma(y)}\right)^{\sigma(x) \sigma(y)} \\
&=\exp _{p}\left(H_{n-1}(\sigma)\right) \prod_{x \in W_{n-2}} \prod_{y \in S(x)} \prod_{z \in S(y)} \sum_{\eta(z) \in\{ \pm 1\}} \exp _{p}[\lambda(\sigma(x), \tilde{\sigma}(y))
\end{aligned}
$$

$$
\left.+J \sum_{\substack{z \in S(x) \\ z \neq y \\ \eta(z) \in \Phi}} \tilde{\sigma}(y) \eta(z)+\tilde{\sigma}(y) h_{y}\right]
$$

From (3.2) it follows that

$$
\begin{equation*}
U_{n-1} Z_{n-1}^{(\mathbf{h})} \mu_{\mathbf{h}}^{(n-1)}(\sigma)=Z_{n}^{(\mathbf{h})} \sum_{\eta} \mu_{\mathbf{h}}^{(n)}(\sigma \vee \eta) \tag{3.11}
\end{equation*}
$$

Since the measure $\mu_{\mathbf{h}}^{(n)},(n \geq 1)$ is a probability measure for each $n \in \mathbb{N}$, i.e.,

$$
\sum_{\sigma \in \Omega_{V_{n-1}}} \mu_{\mathbf{h}}^{(n-1)}(\sigma)=\sum_{\sigma \in \Omega_{V_{n-1}}} \sum_{\eta \in \Omega_{W_{n}}} \mu_{\mathbf{h}}^{(n)}(\sigma \vee \eta)=1
$$

Hence, from (3.11) we obtain

$$
\begin{equation*}
Z_{n}^{(\mathbf{h})}=U_{n-1} Z_{n-1}^{(\mathbf{h})} \tag{3.12}
\end{equation*}
$$

(3.11) and (3.12) imply (3.5). The proof is completed.

This proof yields the uniqueness of the $p$-adic quasi Gibbs measures for our model.

In this study, our main goal is to analyze the $p$-adic dynamical behavior of the fixed points of the dynamic function (3.7). And the existence of three non-trivial fixed points belonging to the model was proved in [3].

Recall that a function $\mathbf{h}=\left\{\mathbf{h}_{x}\right\}_{x \in V \backslash\left\{x^{0}\right\}}$ is called translation-invariant if $\mathbf{h}_{x}=$ $\mathbf{h}_{y}$ for all $x, y \in V$. A $p$-adic Gibbs measure $\mu_{\mathbf{h}}$, corresponding to a translationinvariant function $\mathbf{h}$, is called a translation-invariant p-adic quasi Gibbs measure.

By Theorem 3.1, we reduced the existence of translation-invariant p-adic quasi Gibbs measures to the solutions of equation (3.7). In generally, solutions of such kind of equations are very complicated. Therefore, firstly we need to reduce the equation to the simplest form possible. Then we suppose that $h_{x}=h_{y}=h_{z}=h$ since $h$ is translation invariant for all $x, y, z \in V$. Therefore we rewrite equation (3.7) as following:

$$
\begin{equation*}
h^{2}=\frac{a^{2} \theta^{2} h^{4}+2 a b h^{2}+b^{2} \theta^{2}}{c^{2} \theta^{2} h^{4}+2 c d h^{2}+d^{2} \theta^{2}} \tag{3.13}
\end{equation*}
$$

For more simplicity of (3.13), let $h^{2}=u=f(u)$. Then the last equation is reduced to the following function, and we call it a dynamical function:

$$
\begin{equation*}
f(u)=\frac{a^{2} \theta^{2} u^{2}+2 a b u+b^{2} \theta^{2}}{c^{2} \theta^{2} u^{2}+2 c d u+d^{2} \theta^{2}} \tag{3.14}
\end{equation*}
$$

where $a, b, c, d, \theta \in \mathcal{E}_{p}$. Now we need to analyze fixed points of (3.14).
Lemma 3.2. Let $p \geq 3$, then for any $a, b \in \mathcal{E}_{p}$ one has $|a+1|_{p}=1, \mid a+$ $\left.b\right|_{p}=1$.

The proof is obvious from strong triangle inequality.
3.1. Uniqueness of $p$-adic quasi Gibbs measures. Due to [6], the following lemma guarantees the existence of at least one solution of the dynamical function (3.14) for the mixed type $\lambda$-Ising model.

Lemma 3.3 ([3]). Let $p \geq 3, a, b, c, d, \theta \in \mathcal{E}_{p}$, and $f$ be as given in (3.14).

1. The function $f$ has a unique fixed point $u_{0} \in \mathcal{E}_{p}$.
2. Then $f\left(\mathcal{E}_{p}\right) \subset \mathcal{E}_{p}$ and $|f(u)-f(v)|_{p} \leq \frac{1}{p}|u-v|_{p}$ for all $u, v \in \mathcal{E}_{p}$.

For the proof, we suggest you to refer to [3].
Hence the function $f$ satisfies the Banach contraction principle. This proof guarantees the existence of a unique fixed point $u_{0}$ of $f$ in $\mathcal{E}_{p}$. By using the results of $[14,15]$, we can state the following theorem.

Theorem 3.4. If $p \geq 3$ and $a, b, c, d, \theta \in \mathcal{E}_{p}$, then there exists a unique translational invariant p-adic quasi Gibbs measure $\mu_{0}$ associated to the unique fixed point $u_{0}$ for the model (3.1) on Cayley tree of order two.

Now we need to study other fixed points of (3.14) if they exist. For simplicity, let us assume that $b=\lambda(1,-1)=\lambda(-1,1)=c$ and let us take it as $b$. Then we rewrite (3.14) as follows:

$$
\begin{equation*}
b^{2} \theta^{2} u^{3}+\left(2 b d-a^{2} \theta^{2}\right) u^{2}+\left(d^{2} \theta^{2}-2 a b\right) u-b^{2} \theta^{2}=0 \tag{3.15}
\end{equation*}
$$

where $a, b, d, \theta \in \mathcal{E}_{p}$. Moreover, by Lemma 3.3, there exists at least one fixed point $u_{0} \in \mathcal{E}_{p}$. Therefore we can write (3.15) in the form

$$
\begin{equation*}
\left(u-u_{0}\right)\left(b^{2} \theta^{2} u^{2}+A u+B\right)=0 \tag{3.16}
\end{equation*}
$$

where $A=2 b d+b^{2} \theta^{2} u_{0}-a^{2} \theta^{2}$ and $B=\frac{b^{2} \theta^{2}}{u_{0}}$ since $u_{0}$ is a fixed point of (3.14).
To find other fixed points of (3.16), which are different from $u_{0}$, we need to analyze

$$
\begin{equation*}
u^{2}+\frac{2 b d+b^{2} \theta^{2} u_{0}-a^{2} \theta^{2}}{b^{2} \theta^{2}} u+\frac{1}{u_{0}}=0 \tag{3.17}
\end{equation*}
$$

where $a, b, d, \theta, u_{0} \in \mathcal{E}_{p}$. To find out all fixed points of (3.16), let us perform the results developed in [22]. Therefore, we rewrite (3.17) in the form

$$
\begin{equation*}
u^{2}+e u=f \tag{3.18}
\end{equation*}
$$

where $e=\frac{2 b d+b^{2} \theta^{2} u_{0}-a^{2} \theta^{2}}{b^{2} \theta^{2}}$ and $f=-\frac{1}{u_{0}}, a, b, d, \theta, u_{0} \in \mathcal{E}_{p}$. And the discriminant of (3.18) is $\Delta=e^{2}+4 f$.

To find other fixed points of $(3.18)$ in $\mathbb{Q}_{p}$, we will need the proposition below without the proof.

Proposition 3.5 ([3]). Let $p>3$ and assume that $\sqrt{\Delta}=\sqrt{e^{2}+4 f}$ exists, where $e=\frac{2 b d+b^{2} \theta^{2} u_{0}-a^{2} \theta^{2}}{b^{2} \theta^{2}}, f=-\frac{1}{u_{0}}$, and $a, b, d, \theta, u_{0} \in \mathcal{E}_{p}$. Then equation (3.18) has two fixed points $u_{1,2} \in \mathbb{Q}_{p}$ which are different from $u_{0} \in \mathcal{E}_{p}$.

Due to [3], we get easily two fixed points $u_{1,2} \in \mathbb{Q}_{p}$ which are different from $u_{0} \in \mathcal{E}_{p}$,

$$
\begin{equation*}
u_{1}=\frac{1}{2}(-e+\sqrt{\Delta}), \quad u_{2}=\frac{1}{2}(-e-\sqrt{\Delta}) \tag{3.19}
\end{equation*}
$$

And the conditions $\left|u_{1,2}\right|_{p}=1, \sqrt{|\Delta|_{p}}<1$ and $\left|u_{1,2}-1\right|_{p}=|e \pm \sqrt{\Delta}+2|_{p}=1$ hold.

Now, we are ready to study the dynamical behavior of the fixed points of (3.14).

## 4. Dynamical behavior of the fixed points

In this section, we are going to investigate the $p$-adic norm of the dynamic system at the fixed points $u_{0}$ and $u_{1,2}$. Due to the notion of dynamical system in $\mathbb{Q}_{p}$, the point $x^{(0)}$ is called attractive if $0 \leq|\lambda|_{p}<1$, neutral if $|\lambda|_{p}=1$, and repellent if $|\lambda|_{p}>1$, where $\lambda$ is a usual derivative of the dynamical function.

To study the dynamical behavior of the fixed points $u_{0}, u_{1,2}$ under the dynamic function (3.14), we need to find the derivative of (3.14). After that we are going to check the boundedness of the derivative of (3.14) at the fixed points $u_{0}, u_{1,2}$.

Firstly, before stating the main theorem, we need the following lemma:
Lemma 4.1. Let $p>3, a, b, c, d, \theta \in \mathcal{E}_{p}$ and $u_{1,2} \in \mathbb{Q}_{p}-\mathcal{E}_{p}$. Then $\mid u_{1,2}+$ $\left.1\right|_{p}<1$ and $\left|u_{1}\right|_{p}=\left|u_{2}\right|_{p}=1$.

Proof. From Proposition 3.5 and the strong triangle inequality, we get

$$
\begin{aligned}
\left|u_{1,2}+1\right|_{p} & =\left|\frac{-2 b d-b^{2} \theta^{2} u_{0}+a^{2} \theta^{2} \pm b^{2} \theta^{2} \sqrt{\Delta}+2 b^{2} \theta^{2}}{2 b^{2} \theta^{2}}\right|_{p} \\
& =\frac{\left|-2(b d-1)-\left(b^{2} \theta^{2} u_{0}-1\right)+\left(a^{2} \theta^{2}-1\right)\right|_{p}}{\left|2 b^{2} \theta^{2}\right|_{p}} \\
& \pm \frac{\left|b^{2} \theta^{2} \sqrt{\Delta}+2\left(b^{2} \theta^{2}-1\right)\right|_{p}}{\left|2 b^{2} \theta^{2}\right|_{p}}<\frac{1}{p}
\end{aligned}
$$

and from (3.19), we get $\left|u_{1}\right|_{p}=\left|u_{2}\right|_{p}=1$.
Now we are ready to state the main theorem.
Theorem 4.2. Let $p>3$ and $a, b, c, d, \theta \in \mathcal{E}_{p}$. Then the p-adic dynamical function (3.14) has three fixed points which are $u_{0} \in \mathcal{E}_{p}$ and $u_{1,2} \in \mathbb{Q}_{p}-\mathcal{E}_{p}$ under the condition $b=\exp _{p}(\lambda(1,-1))=\exp _{p}(\lambda(-1,1))=c$, and
(i) $u_{0}$ is an attractor;
(ii) $u_{1,2}$ are repellent.

Proof. From Lemma 3.3 and Proposition 3.5, the dynamical function (3.14) has three fixed points which are $u_{0} \in \mathcal{E}_{p}$ and $u_{1,2} \in \mathbb{Q}_{p}-\mathcal{E}_{p}$. In here, to prove (i) and (ii), let

$$
\begin{equation*}
f(u)=\frac{a^{2} \theta^{2} u^{2}+2 a b u+b^{2} \theta^{2}}{b^{2} \theta^{2} u^{2}+2 b d u+d^{2} \theta^{2}} \tag{4.1}
\end{equation*}
$$

where $a, b, d, \theta \in \mathcal{E}_{p}$. If we find $\lambda=\frac{d f}{d u}$, then we get

$$
\begin{equation*}
f^{\prime}(u)=\frac{2\left[\left(a^{2} b d \theta^{2}-a b^{2} \theta^{2}\right) u^{2}+\left(a^{2} d^{2} \theta^{2}-b^{4} \theta^{4}\right) u+a b d^{2} \theta^{2}-b^{3} d \theta^{2}\right]}{\left[b^{2} \theta^{2} u^{2}+2 b d u+d^{2} \theta^{2}\right]^{2}} \tag{4.2}
\end{equation*}
$$

Hereafter, to prove (i) and (ii), we are going to check the $p$-adic norm of $\lambda=$ $f^{\prime}(u)$ at $u_{0}, u_{1,2}$.
(i) Let us plug $u_{0}$ in (4.2). Then we get

$$
\begin{align*}
\left.\lambda\right|_{u_{0}}= & f^{\prime}\left(u_{0}\right) \\
= & 2 \frac{\left(a^{2} b d \theta^{2}-a b^{2} \theta^{2}\right) u_{0}^{2}+\left(a^{2} d^{2} \theta^{2}-b^{4} \theta^{4}\right) u_{0}+a b d^{2} \theta^{2}-b^{3} d \theta^{2}}{\left[b^{2} \theta^{2} u_{0}^{2}+2 b d u_{0}+d^{2} \theta^{2}\right]^{2}} \\
= & 2 \frac{\left(a^{2} b d \theta^{2}-1-\left(a b^{2} \theta^{2}-1\right)\right) u_{0}^{2}+\left(a^{2} d^{2} \theta^{2}-1-\left(b^{4} \theta^{4}-1\right)\right) u_{0}}{\left[\left(b^{2} \theta^{2}-1\right) u_{0}^{2}+u_{0}^{2}+2(b d-1) u_{0}+u_{0}+d^{2} \theta^{2}-1+1\right]^{2}} \\
& +2 \frac{a b d^{2} \theta^{2}-1-\left(b^{3} d \theta^{2}-1\right)}{\left[\left(b^{2} \theta^{2}-1\right) u_{0}^{2}+u_{0}^{2}+2(b d-1) u_{0}+u_{0}+d^{2} \theta^{2}-1+1\right]^{2}} \tag{4.3}
\end{align*}
$$

From here, taking the $p$-adic norm of (4.3), we obtain

$$
\left|f^{\prime}\left(u_{0}\right)\right|_{p} \leq \frac{1}{p}<1
$$

since $|2|_{p}=1$ and $\left|u_{0}+1\right|_{p}=1$.
Then $u_{0}$ is an attractor.
(ii) Let us plug $u_{1,2}$ in (4.2). Then we get

$$
\begin{align*}
\left.\lambda\right|_{u_{1,2}} & =f^{\prime}\left(u_{1,2}\right) \\
& =2 \frac{\left(a^{2} b d \theta^{2}-a b^{2} \theta^{2}\right) u_{1,2}^{2}+\left(a^{2} d^{2} \theta^{2}-b^{4} \theta^{4}\right) u_{1,2}+a b d^{2} \theta^{2}-b^{3} d \theta^{2}}{\left[b^{2} \theta^{2} u_{1,2}^{2}+2 b d u_{1,2}+d^{2} \theta^{2}\right]^{2}} \\
& =2 \frac{\left(a^{2} b d \theta^{2}-1-\left(a b^{2} \theta^{2}-1\right)\right) u_{1,2}^{2}+\left(a^{2} d^{2} \theta^{2}-1-\left(b^{4} \theta^{4}-1\right)\right) u_{1,2}}{\left[\left(b^{2} \theta^{2}-1\right) u_{1,2}^{2}+u_{1,2}^{2}+2(b d-1) u_{1,2}+u_{1,2}+d^{2} \theta^{2}-1+1\right]^{2}} \\
& +2 \frac{a b d^{2} \theta^{2}-1-\left(b^{3} d \theta^{2}-1\right)}{\left[\left(b^{2} \theta^{2}-1\right) u_{1,2}^{2}+u_{1,2}^{2}+2(b d-1) u_{1,2}+u_{1,2}+d^{2} \theta^{2}-1+1\right]^{2}} \tag{4.4}
\end{align*}
$$

Since

$$
\left|\left(b^{2} \theta^{2}-1\right) u_{1,2}^{2}+u_{1,2}^{2}+2(b d-1) u_{1,2}+u_{1,2}+d^{2} \theta^{2}-1+1\right|_{p}=\left|u_{1,2}+1\right|_{p}
$$

taking the $p$-adic norm of (4.4), we obtain

$$
\left|f^{\prime}\left(u_{1,2}\right)\right|_{p} \geq p^{3}>1
$$

from Lemma 4.1 and the strong triangle inequality.
In this case, we conclude that $u_{1,2}$ are repellent.

Then, as a conclusion, we state that the fixed point $u_{0}$ is an attractor and the other fixed points $u_{1,2}$ are repellent.

Now we are going to describe the basin of attraction of the fixed point $u_{0}$,

$$
A\left(u_{0}\right)=\left\{u \in \mathbb{Q}_{p}: f^{n}(u) \rightarrow u_{0}\right\}
$$

Let

$$
\mathcal{B}=\left\{u \in S_{1}\left(u_{0}\right): \exists n_{0} \in \mathbb{N} f^{n_{0}}(x) \in K\right\}
$$

where $K=\left\{u \in S_{1}\left(u_{0}\right):\left|u_{0}-1\right|_{p}<1\right\}$.
One can see that

$$
\begin{equation*}
\mathcal{B}=\left(\bigcup_{n \geq 1} f^{-n}(K)\right) \bigcap S_{1}\left(u_{0}\right) \tag{4.5}
\end{equation*}
$$

Now we are going to describe the size of an attractor of the dynamic function as in [16].

Theorem 4.3. Let $a, b, d, \theta \in \mathcal{E}_{p}$. Then the following statements hold:
(i) If $p>3, \sqrt{\Delta}=\sqrt{\left(\frac{2 b d+b^{2} \theta^{2} u_{0}-a^{2} \theta^{2}}{b^{2} \theta^{2}}\right)^{2}-4 \frac{1}{u_{0}}}$, then one has

$$
A\left(u_{0}\right)=\mathcal{B} \cup\left(\mathbb{Q}_{p} \backslash S_{1}\left(u_{0}\right)\right) \cup B_{1}(0)
$$

(ii) Otherwise, $A\left(u_{0}\right)=\mathbb{Q}_{p}$.

Proof. (i) Due to Lemma 3.3, for any $u \in \mathcal{E}_{p}$, we infer that $u \in A\left(u_{0}\right)$, which means that $\mathcal{E}_{p} \subset A\left(u_{0}\right)$. We notice that $\mathcal{E}_{p}=B_{1}\left(u_{0}\right)$.

Consider several cases.
(I) Assume that $|u|_{p}<1$, i.e., $u \in B_{1}(0)$. Now, let us examine $|f(u)|_{p}$ and $|f(u)-1|_{p}$. One can see that

$$
\begin{aligned}
|f(u)|_{p} & =\frac{\left|\left(a^{2} \theta^{2}-1\right) u^{2}+u^{2}+2(a b-1) u+2 u+\left(b^{2} \theta^{2}-1\right)+1\right|_{p}}{\left|\left(b^{2} \theta^{2}-1\right) u^{2}+u^{2}+2(b d-1) u+2 u+\left(d^{2} \theta^{2}-1\right)+1\right|_{p}} \\
& =\frac{|u+1|_{p}^{2}}{|u+1|_{p}^{2}}=1
\end{aligned}
$$

and

$$
|f(u)-1|_{p}=\frac{\left|\left(a^{2} \theta^{2}-b^{2} \theta^{2}\right) u^{2}+2(a b-b d) u+b^{2} \theta^{2}-d^{2} \theta^{2}\right|_{p}}{|u+1|_{p}^{2}}<1
$$

which imply $f(u) \in \mathcal{E}$. Hence, $u \in A\left(u_{0}\right)$, which means that $B_{1}(0) \subset A\left(u_{0}\right)$. So,

$$
B_{1}(0) \cup \mathcal{E}_{p} \subset A\left(u_{0}\right)
$$

(II) Now assume that $|u|_{p}>1$ taking into account that $u \notin S_{1}\left(u_{0}\right)$. Then we have

$$
|f(u)|_{p}=\frac{|u+1|_{p}^{2}}{|u+1|_{p}^{2}}=1
$$

and since

$$
\max \left\{\left|\left(a^{2} \theta^{2}-1\right) u^{2}\right|_{p},\left|\left(b^{2} \theta^{2}-1\right) u^{2}\right|_{p}\right\}<|u|_{p}^{2} \quad \text { and } \quad|u+1|_{p}^{2}=|u|_{p}^{2}
$$

we obtain

$$
\begin{aligned}
|f(u)-1|_{p} & =\frac{\left|\left(a^{2} \theta^{2}-b^{2} \theta^{2}\right) u^{2}+2(a b-1) u-(b d-1) u+b^{2} \theta^{2}-d^{2} \theta^{2}\right|_{p}}{|u+1|_{p}^{2}} \\
& =\frac{\left|\left(a^{2} \theta^{2}-1\right) u^{2}-\left(b^{2} \theta^{2}\right) u^{2}+u^{2}+2(a b-b d) u+b^{2} \theta^{2}-d^{2} \theta^{2}\right|_{p}}{|u+1|_{p}^{2}} \\
& =\frac{\max \left\{\left|\left(a^{2} \theta^{2}-1\right) u^{2}\right|_{p},\left|\left(b^{2} \theta^{2}-1\right) u^{2}\right|_{p}\right\}}{|u|_{p}^{2}}<1,
\end{aligned}
$$

which means $f(u) \in \mathcal{E}_{p}$. So, $u \in A\left(u_{0}\right)$.
Consequently, if $|u|_{p}>1$, then $u \in A\left(u_{0}\right)$.
(III) Let us consider $|u|_{p}=1$. Then

$$
|f(u)|_{p}=\frac{|u+1|_{p}^{2}}{|u+1|_{p}^{2}}=1
$$

and since

$$
\begin{aligned}
& \max \left\{\left|\left(a^{2} \theta^{2}-1\right) u^{2}\right|_{p},\left|\left(b^{2} \theta^{2}-1\right) u^{2}\right|_{p}\right\}<|u|_{p}^{2}=1 \quad \text { and } \quad|u+1|_{p}^{2}=|u|_{p}^{2}=1, \\
& |f(u)-1|_{p}=\frac{\left|\left(a^{2} \theta^{2}-b^{2} \theta^{2}\right) u^{2}+2(a b-1) u-(b d-1) u+b^{2} \theta^{2}-d^{2} \theta^{2}\right|_{p}}{|u+1|_{p}^{2}} \\
& =\frac{\left|\left(a^{2} \theta^{2}-1\right) u^{2}-\left(b^{2} \theta^{2}\right) u^{2}+u^{2}+2(a b-b d) u+b^{2} \theta^{2}-d^{2} \theta^{2}\right|_{p}}{|u+1|_{p}^{2}} \\
& =\frac{\max \left\{\left|\left(a^{2} \theta^{2}-1\right) u^{2}\right|_{p},\left|\left(b^{2} \theta^{2}-1\right) u^{2}\right|_{p}\right\}}{|u|_{p}^{2}}<1 .
\end{aligned}
$$

Hence one gets $f(u) \in \mathcal{E}_{p}$, so $u \in A\left(u_{0}\right)$.
Therefore, we have

$$
\mathcal{B} \cup\left(\mathbb{Q}_{p} \backslash S_{1}\left(x_{0}\right)\right) \cup B_{1}(0) \subset A\left(x_{0}\right) .
$$

Now assume that $u \notin \mathcal{B} \cup\left(\mathbb{Q}_{p} \backslash S_{1}\left(u_{0}\right)\right) \cup B_{1}(0)$. Due to (4.5), it yields that $f^{n}(u) \notin K$ for all $n \in \mathbb{N}$. This means that

$$
\begin{equation*}
\left|f^{n}(u)-1\right|_{p} \nless 1, \quad n \in \mathbb{N} . \tag{4.6}
\end{equation*}
$$

If $f^{n}(u) \rightarrow u_{0}$ as $n \rightarrow \infty$, then from (4.6) one finds $\left|f^{n}(u)-1\right|_{p}<1$, which is a contradiction. Therefore, we infer that $u \notin A\left(u_{0}\right)$. So,

$$
A\left(u_{0}\right) \subset \mathcal{B} \cup\left(\mathbb{Q}_{p} \backslash S_{1}\left(u_{0}\right)\right) \cup B_{1}(0)
$$

This completes the part (i).
(ii) The proof of this part immediately follows from (i).

## 5. Conclusions

In [3], it is proven that the $p$-adic $\lambda$-Ising model has three fixed points, which means that in the model there exist three transition invariant $p$-adic quasi Gibbs measures on the Cayley tree of order two. In this work, we continued the investigation of the dynamical properties of these fixed points.

As a main result of this paper, we proved that the $p$-adic norm of the derivative of the dynamical function (3.14) are $\left|f^{\prime}\left(u_{0}\right)\right|_{p}<1$ and $\left|f^{\prime}\left(u_{1,2}\right)\right|_{p}>1$ at the fixed points. Hence, $u_{0}$ is an attractor and $u_{1,2}$ are repellent fixed points. Moreover, we described the size of the basin attractor of the fixed point, $u_{0}$.

Acknowledgments. The author thanks the reviewer for precious comments and help.

## References

[1] R.J. Baxter, Exactly Solved Models in Statistical Mechanics, Academic Press, London, 1982.
[2] S. De Smedt and A. Khrennikov, A p-adic behavior of dynamical systems, Rev. Mat. Complut. 12 (1999), 301-323.
[3] M. Dogan, Phase transition of the mixed type p-adic $\lambda$-Ising model on the Cayley tree, p-Adic Numbers Ultrametric Anal. Appl. 10 (2018), No. 4, 276-286.
[4] N.N. Ganikhodjaev, F.M. Mukhamedov, and U.A. Rozikov, Phase transitions of the Ising model on $\mathbb{Z}$ in the p-adic number field, Uzbek. Mat. Zh. 4 (1998), 23-29 (Russian).
[5] H.O. Georgii, Gibbs Measures and Phase Transitions, De Gruyter Studies in Mathematics, 9, Walter de Gruyter \& Co., Berlin, 1988.
[6] M. Khamraev, F.M. Mukhamedov, and U.A. Rozikov, On uniqueness of Gibbs measure for $p$-adic $\lambda$-model on the Cayley tree, Lett. Math. Phys. 70 (2004), 17-28.
[7] A. Khrennikov, Non-Archimedean Analysis: Quantum Paradoxes, Dynamical Systems and Biological Models, Mathematics and its Applications, 427, Kluwer Academic Publishers, Dordrecht, 1997.
[8] A.Yu. Khrennikov and S. Ludkovsky, Stochastic processes on non-Archimedean spaces with values in non-Archimedean fields, Markov Process. Related Fields 9 (2003), 131-162.
[9] N. Koblitz, p-Adic Numbers, p-Adic Analysis and Zeta-Function, Graduate Texts in Mathematics, 58, Springer-Verlag, New York-Heidelberg, 1977.
[10] J. Lubin, Non-Archimedean dynamical systems, Compositio Math Math. 94 (1994), 321-346.
[11] F. Mukhamedov, On factor associated with the unordered phase of $\lambda$-model on a Cayley tree, Rep. Math. Phys. 53 (2004), 1-18.
[12] F. Mukhamedov, A dynamical system appoach to phase transitions p-adic Potts model on the Cayley tree of order two, Rep. Math. Phys. 70 (2012), 385-406.
[13] F. Mukhamedov, On dynamical systems and phase transitions for $q+1$-state $p$-adic Potts model on the Cayley tree, Math. Phys. Anal. Geom. 16 (2013), 49-87.
[14] F. Mukhamedov, Recurrence equations over trees in a non-Archimedean context, $p$-Adic Numbers Ultrametric Anal. Appl. 6 (2014), 310-317.
[15] F. Mukhamedov and H. Akin, On non-Archimedean recurrence equations and their applications, J. Math. Anal. Appl. 423 (2015), 1203-1218.
[16] F. Mukhamedov, M. Dogan, and H. Akin, On chaotic behaviour of the p-adic generalized Ising mapping and its application, J. Difference Equ. Appl. 23 (2017), 1542-1561.
[17] F. Mukhamedov and U.A. Rozikov, On rational p-adic dynamical systems, Methods of Funct. Anal. and Topology 10 (2004), 21-31.
[18] F. Mukhamedov, M. Saburov, and O. Khakimov, On p-adic Ising-Vannimenus model on an arbitraray order Cayley tree, J. Stat. Mech. Theory Exp. (2015), No. 5, P05032.
[19] M. Ostilli, Cayley trees and Bethe lattices: A concise analysis for mathematicians and physicists, Phys. A 391 (2012), 3417-3423.
[20] U.A. Rozikov, Description of limit Gibbs measures for $\lambda$-models on the Bethe lattice, Siberian Math. J. 39 (1998), 373-380.
[21] U.A. Rozikov, Gibbs Measures on Cayley Trees, World Scientific Publishing Co. Pte. Ltd., Hackensack, NJ, 2013.
[22] M. Saburov and M.A. Khameini, Quadratic equations over p-adic fields and their applications in statistical mechanics, ScienceAsia 41 (2015), 209-215.
[23] A.N. Shiryaev, Probability, Nauka, Moscow, 1980 (Russian).
Received April 4, 2018, revised June 11, 2018.
Mutlay Dogan,
University of Bahamas, Faculty of Pure and Applied Sciences, Oakes Field Campus, N 4912, Nassau, Bahamas,
E-mail: mutlay74@gmail.com

## Про динамічну поведінку $p$-адичної $\lambda$-ізінгової моделі на дереві Кейлі

Mutlay Dogan

У даній статті ми продовжуємо вивчення деяких властивостей $p$ адичної $\lambda$-ізінгової моделі змішаного типу, що вивчалася в [3]. У тій роботі існування $p$-адичних мір Гібса і фазових переходів досліджувалося в моделі на дереві Кейлі другого порядка. У даній статті ми вивчаємо динамічну поведінку нерухомих точок, які було знайдено в [3]. Основним результатом є те, що ми довели, що нерухома точка $u_{0}$ є атрактором, а інші нерухомі точки $u_{1,2}$ репелентними нерухомими точками для $p$ адичної $\lambda$-ізінгової моделі змішаного типу. На додаток описано розмір басейну атрактора для нерухомої точки $u_{0}$.

Ключові слова: $p$-адичні числа, $p$-адична квазіміра Гіббса, динамічні системи, дерево Кейлі.


[^0]:    (C) Mutlay Dogan, 2019

