The natural extension of the random beta-transformation

YOUNÈS TIERCE

Univ. Rouen Normandie, CNRS, LMRS, UMR 6085, F-76000 Rouen, France (e-mail: younes.tierce@gmail.com)

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Abstract. We construct a geometrico-symbolic version of the natural extension of the random β -transformation introduced by Dajani and Kraaikamp [Random β -expansions. *Ergod. Th. & Dynam. Sys.* **23**(2) (2003) 461–479]. This construction provides a new proof of the existence of a unique absolutely continuous invariant probability measure for the random β -transformation, and an expression for its density. We then prove that this natural extension is a Bernoulli automorphism, generalizing to the random case the result of Smorodinsky [β -automorphisms are Bernoulli shifts. *Acta Math. Acad. Sci. Hungar.* **24** (1973), 273–278] about the greedy transformation.

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1. Introduction

1.1. *Expansions in base* β . Throughout the paper, we fix a real number β , $1 < \beta < 2$. For $x \in [0, 1)$, an expansion of x in base β is a sequence $(x_n)_{n \ge 1}$ of $\{0, 1\}$ such that

$$x = \sum_{n=1}^{+\infty} \frac{x_n}{\beta^n}.$$

Rényi [13] introduced the *greedy map* T_{β} defined on [0, 1) by $T_{\beta}(x) = \beta x \mod 1$, which provides a particular expansion of any real number in [0, 1), given by the sequence $x_n := \mathbb{1}_{[1/\beta,1)}(T_{\beta}^{n-1}(x))$, for all $n \ge 1$. This expansion is called the *greedy expansion* of x. In fact, almost every $x \in [0, 1)$ has an infinite number of expansions in base β , and the greedy expansion is the greatest in the lexicographic order [7].

Rényi proved that T_{β} has a unique absolutely continuous invariant probability measure ν_{β} , and Parry [12] proved that its density is proportional to the function



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$$\sum_{n \ge 0} \frac{1}{\beta^n} \mathbb{1}_{[0, T^n_\beta(1)]}$$

and that the associated measure-preserving system is weakly mixing.

Smorodinsky then showed in [15] that the natural extension of this system is a Bernoulli automorphism. This result was also obtained by Dajani, Kraaikamp, and Solomyak in [6] through a geometric construction of this natural extension, in the form of a tower.

The greedy map can be extended to a map T_g defined on the interval $I_\beta := [0, 1/(\beta - 1)]$ by setting

$$T_g(x) := \begin{cases} T_\beta(x) & \text{for } x \in [0, 1[, \\ \beta x - 1 & \text{for } x \in [1, 1/(\beta - 1)]. \end{cases}$$

We will still refer to this extended map as the greedy map. The measure ν_{β} extended to I_{β} by setting $\nu_{\beta}([1, 1/(\beta - 1)]) = 0$ is still the unique absolutely continuous invariant probability measure of T_{g} on I_{β} .

We can obtain the smallest expansion of any real number of the interval I_{β} with the *lazy* map T_{ℓ} [7], defined by

$$T_{\ell}: I_{\beta} \to I_{\beta}$$

$$x \mapsto \begin{cases} \beta x & \text{if } x \leq \frac{1}{\beta(\beta - 1)}, \\ \beta x - 1 & \text{otherwise.} \end{cases}$$

For $x \in I_{\beta}$, the *lazy expansion* of x is given by the sequence (x'_n) defined for all $n \ge 1$ by $x'_n = \mathbb{1}_{(1/\beta(\beta-1),1/(\beta-1))}(T_{\ell}^{n-1}(x)).$

Let *s* be the symmetry on I_{β} , defined by $s(x) = 1/(\beta - 1) - x$. The pushforward measure $\tilde{v}_{\beta} := v_{\beta} \circ s^{-1}$ is the unique absolutely continuous invariant measure for T_{ℓ} , and *s* conjugates T_g and T_{ℓ} ; hence the systems $(I_{\beta}, \mathcal{B}, v_{\beta}, T_g)$ and $(I_{\beta}, \mathcal{B}, \tilde{v}_{\beta}, T_{\ell})$ are isomorphic [4].

In 2003, Dajani and Kraaikamp [5] introduced the random β transformation. We set $\Omega := \{g, \ell\}^{\mathbb{N} \cup \{0\}}$, the set of sequences of g and ℓ . We then define the transformation K_{β} on $\Omega \times I_{\beta}$ by

$$K_{\beta} : \frac{\Omega \times I_{\beta} \to \Omega \times I_{\beta}}{(\omega, x)} \mapsto (\sigma(\omega), T_{\omega_0}(x)),$$

where σ is the left shift on Ω . The sequence $\omega \in \Omega$ describes the successive transformations that will be applied to the real $x \in I_{\beta}$: once T_{ω_0} is applied to x, we shift the sequence ω to the left and then apply T_{ω_1} , and so on. If we fix $\omega \in \Omega$, we obtain a particular expansion of x in base β , by setting, for any $n \ge 1$:

$$x_{n}^{\omega} := \begin{cases} \mathbbm{1}_{\{\pi(K_{\beta}^{n-1}(\omega, x)) \ge 1/\beta\}} & \text{if } \omega_{n-1} = g, \\ \mathbbm{1}_{\{\pi(K_{\beta}^{n-1}(\omega, x)) > 1/\beta(\beta-1)\}} & \text{if } \omega_{n-1} = \ell, \end{cases}$$

where π is the projection on I_{β} . Any expansion of x can be obtain with a sequence $\omega \in \Omega$ [3].

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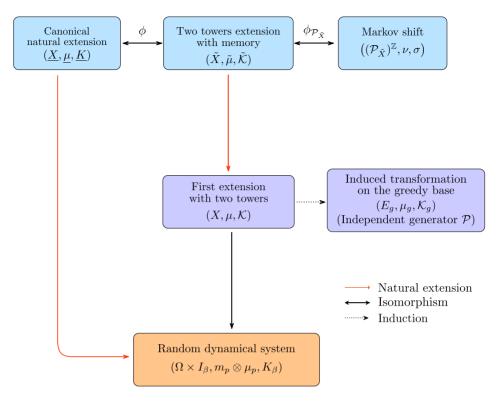


FIGURE 1. Extensions involved in this paper.

In the model we consider, the sequence ω will be drawn with the product Bernoulli measure m_p , where $p \in (0, 1)$ is a fixed parameter. In other words, we draw each transformation independently with probability p for T_g and 1 - p for T_ℓ . Dajani and De Vries proved in [3] that there is a unique absolutely continuous probability measure μ_p on I_β such that $m_p \otimes \mu_p$ is K_β -invariant. Kempton [11] obtained an expression of the density $\rho_{1/2}$ of $\mu_{1/2}$ with the construction of a natural extension of ($\Omega \times I_\beta, m_{1/2} \otimes \mu_{1/2}, \mathcal{K}_\beta$), using two symmetric towers. Suzuki [16] then generalized this expression for any p with the use of the Perron–Frobenius operators.

1.2. Roadmap of this paper. In this paper, we consider several extensions of the random β -transformation. See Figure 1. In §2, we construct a first geometrico-symbolic extension \mathcal{K} of the random transformation. This extension is defined on two towers (the greedy tower and the lazy tower), each tower having a 'base'. By providing a simple invariant measure on this extension and projecting it on I_{β} , we obtain the expression of the density of μ_p for any $p \in (0, 1)$ (formula (8)). We then study the transformation \mathcal{K}_g induced by \mathcal{K} on the base of the greedy tower, and define a particular partition of this base. We prove that this partition is an independent generator of the induced system (Proposition 10), and thus prove that this induced system is isomorphic to a unilateral Bernoulli shift. This provides a new proof of the ergodicity of the random system and of the uniqueness of μ_p . This first extension

is not yet invertible. In §3, we construct a geometrico-symbolic version of the natural extension of this first extension. With the use of the relatively independent joining above a factor, we prove that this new extension is in fact the natural extension of the random β -transformation (Theorem 21). We then prove that this extension is isomorphic to a bilateral Bernoulli shift (Theorem 26), by 'unfolding' the previous partition on the two towers.

2. The first extension

In this section, we extend the dynamics K_{β} on two towers. Kempton [11] implemented the same type of strategy in the case $p = \frac{1}{2}$, but his construction does not generalize to any p. We construct a different extension, valid for any $p \in (0, 1)$.

2.1. The domain of the extension: the greedy and lazy towers. As in [11], the towers are made of floors. We first define the 'base' of each tower. Let E_g be the base of the greedy tower \mathcal{G} , defined by

$$E_g := \Omega \times \{g\} \times [0, 1].$$

The label of the greedy base is (g), telling us that it is the first floor of the greedy tower \mathcal{G} . Likewise, E_{ℓ} is the base of the lazy tower \mathcal{L} , defined by

$$E_{\ell} := \Omega \times \{\ell\} \times \left[s(1), \frac{1}{\beta - 1}\right].$$

A floor in one of the towers is denoted by E_e , where *e* is the label of the floor. This floor is of the form

$$E_e := \Omega \times \{e\} \times I_e,$$

where I_e is a sub-interval of I_β that will be specified later. Thus, a point in one of the two towers is of the form $c = (\omega, e, x)$, where the following hold.

- The sequence ω ∈ {g, ℓ}^{ℕ∪{0}} describes the successive transformations T_g and T_ℓ to be applied to x.
- The label *e* is of the form

$$e = (h, \omega_{-n}, \omega_{-n+1}, \dots, \omega_{-1})$$

and characterizes the floor containing the point *c*. The letter *h* indicates the tower in which: if h = g, the floor is from the greedy tower; and if $h = \ell$, the floor is from the lazy tower. The integer *n* is the *level* of the floor E_e . Therefore, there are 2^n floors of level *n* in each tower, the bases being at level 0. Finally, the symbols $\omega_{-n}, \ldots, \omega_{-1}$ indicate the applied transformations since the last passage in one of the two bases; in other words, the maps $T_{\omega_{-n}}$ to $T_{\omega_{-1}}$ were applied to the real component of a point in E_e since the base E_h .

• x represents the real component of the point c, on which will be applied the transformation T_{ω_0} . We have $x \in I_e$, where

$$\begin{cases} I_e := [0, T_{\omega_{-n}, \omega_{-n+1}, \dots, \omega_{-1}}(1^+)] & \text{if } h = g, \\ I_e := \left[T_{\omega_{-n}, \omega_{-n+1}, \dots, \omega_{-1}}(s(1)^-), \frac{1}{\beta - 1} \right] & \text{if } h = \ell, \end{cases}$$

where the + and – signs respectively refer to the right and left limit at the considered points, and $T_{\omega_{-n},\omega_{-n+1},...,\omega_{-1}}$ is the composition $T_{\omega_{-1}} \circ \cdots \circ T_{\omega_{-n}}$. For simplification, we will often denote T_v with $v \in \{g, \ell\}^n$ to refer to this type of composition.

We then define the greedy tower as the disjoint union of all floors whose label starts with a g, and the lazy tower as the disjoint union of all floors whose label starts with an ℓ :

$$\mathcal{G} := \bigsqcup_{\substack{e' \in \{g,\ell\}^* \\ \ell :=}} E_{g,\ell} E_{g,\ell'},$$

where the point \cdot represents the concatenation, and the set $\{g, \ell\}^*$ is the set of finite sequences (including the empty sequence) of g and ℓ . We then denote by $X := \mathcal{G} \sqcup \mathcal{L}$ the disjoint union of the two towers, on which will be defined \mathcal{K} . Finally, given a floor of X, we call the *length* of this floor the Lebesgue measure of the associated interval.

Applying the dynamics \mathcal{K} to a point in a tower consists in going up a level in the tower, or going back to the base of the tower, or going to the base of the other tower, depending on certain conditions that we will detail in the next section.

We represent the two towers in Figure 2, one above the other (the lazy tower is represented upside down). The ω -component is represented vertically. Two points on a same vertical line have the same real coordinate *x*.

2.2. The dynamics on the two towers. We define the dynamics \mathcal{K} on the greedy tower first. We consider a floor of \mathcal{G} , with label $e := (g, \omega_{-n}, \ldots, \omega_{-1})$. The floor E_e is split into two parts, depending on ω_0 :

$$E_e = (E_e \cap [g]_0) \sqcup (E_e \cap [\ell]_0),$$

where $E_e \cap [g]_0$ abusively refers to the set of points (ω, e, x) of E_e such that $\omega_0 = g$ (and in the same way for ℓ). The dynamics \mathcal{K} is different on each of these parts. In the following, we denote by *t* the upper bound of I_e .

(1) If $\omega_0 = g$ and if $t < 1/\beta$ (see Figure 3):

$$\mathcal{K}: \frac{E_e \cap [g]_0 \to E_{e \cdot g}}{(\omega, e, x) \mapsto (\sigma(\omega), e \cdot g, \beta x).}$$

Here, \mathcal{K} sends $E_e \cap [g]_0$ onto the floor $E_{e \cdot g}$.

(2) If $\omega_0 = g$ and if $t \ge 1/\beta$, the left part of the floor is sent back onto the base E_g (see Figure 4):

$$\mathcal{K}: E_e \cap \left\{ \omega_0 = g, x < \frac{1}{\beta} \right\} \to E_g$$
$$(\omega, e, x) \mapsto (\sigma(\omega), g, \beta x);$$

and the right part goes up onto $E_{e \cdot g}$:

$$\mathcal{K}: E_e \cap \left\{ \omega_0 = g, x \ge \frac{1}{\beta} \right\} \to E_{e \cdot g}$$
$$(\omega, e, x) \mapsto (\sigma(\omega), e \cdot g, \beta x - 1).$$

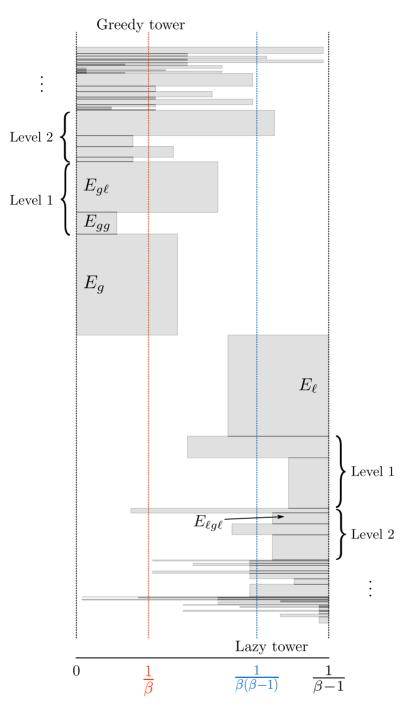
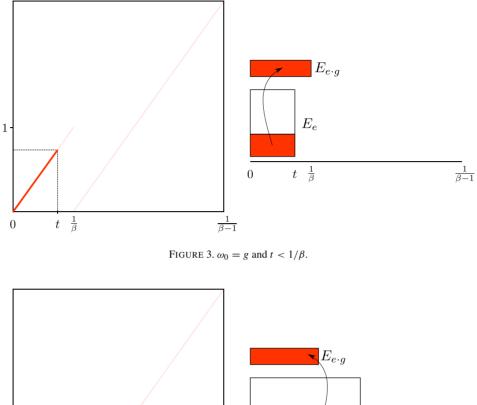


FIGURE 2. Greedy and lazy towers.



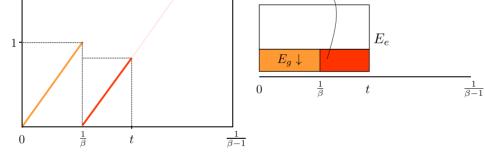


FIGURE 4. $\omega_0 = g$ and $t \ge 1/\beta$.

(3) If $\omega_0 = \ell$ and if $t < 1/\beta(\beta - 1)$ (see Figure 5):

$$\mathcal{K}: \frac{E_e \cap [\ell]_0 \to E_{e \cdot \ell}}{(\omega, e, x) \mapsto (\sigma(\omega), e \cdot \ell, \beta x).}$$

The part $E_e \cap [\ell]_0$ is sent onto the floor $E_{e \cdot \ell}$.

(4) If ω₀ = ℓ and if t ≥ 1/β(β − 1) (see Figure 6), the central part of the floor is sent onto the base E_ℓ:

$$\mathcal{K}: E_e \cap \left\{ \omega_0 = \ell, \frac{s(1)}{\beta} < x \leq \frac{1}{\beta(\beta - 1)} \right\} \to E_\ell$$
$$(\omega, e, x) \mapsto (\sigma(\omega), \ell, \beta x);$$

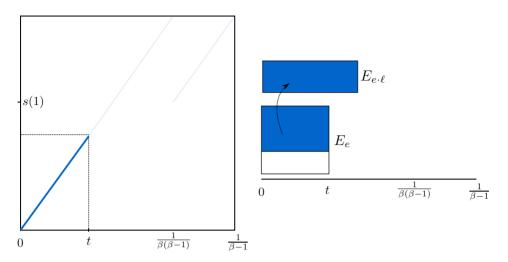


FIGURE 5. $\omega_0 = \ell$ and $t < 1/\beta(\beta - 1)$.

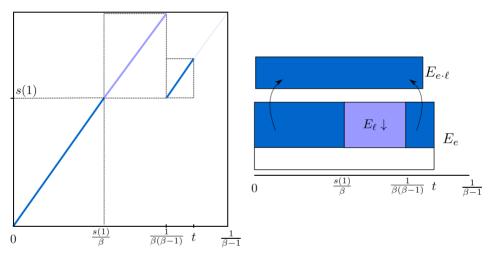


FIGURE 6. $\omega_0 = \ell$ and $t \ge 1/\beta(\beta - 1)$.

and the left and right parts go up onto $E_{e \cdot \ell}$:

$$\mathcal{K}: E_e \cap \left\{ \omega_0 = \ell, x \leqslant \frac{s(1)}{\beta} \right\} \to E_{e \cdot \ell}$$
$$(\omega, e, x) \mapsto (\sigma(\omega), e \cdot \ell, \beta x);$$
$$\mathcal{K}: E_e \cap \left\{ \omega_0 = \ell, x > \frac{1}{\beta(\beta - 1)} \right\} \to E_{e \cdot \ell}$$
$$(\omega, e, x) \mapsto (\sigma(\omega), e \cdot \omega_0, \beta x - 1).$$

On the lazy tower, the dynamics \mathcal{K} is defined in the same way by symmetry (see Figures 7 and 8).

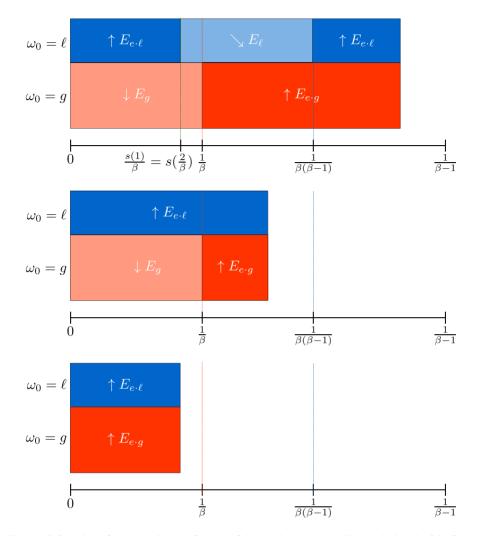


FIGURE 7. Overview of the dynamics on a floor E_e of the greedy tower, depending on the length of the floor.

We denote by $\pi : X \to \Omega \times I_{\beta}$ the projection onto $\Omega \times I_{\beta}$. By construction, we have

$$\pi \circ \mathcal{K} = K_{\beta} \circ \pi.$$

2.3. Construction of an invariant measure. The goal of this section is to define a \mathcal{K} -invariant measure μ on X, such that the projection of this measure on I_{β} is absolutely continuous. We denote the Lebesgue measure by λ , regardless of the interval on which we consider it. We recall that m_p is the product Bernoulli measure of parameter p on Ω (we draw g with probability p and ℓ with probability 1 - p, independently at each step).

We first define the measures μ_g and μ_ℓ with respective support \mathcal{G} and \mathcal{L} as follows.

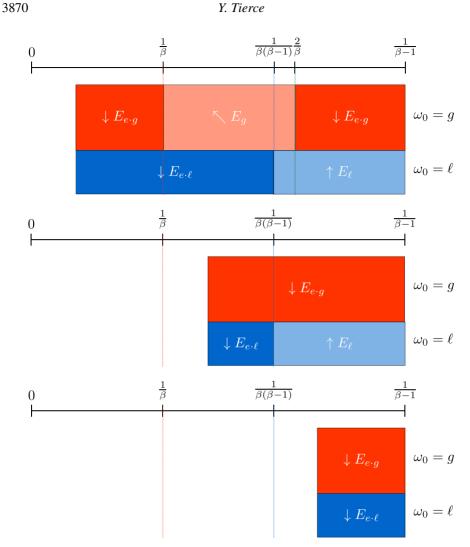


FIGURE 8. Overview of the dynamics on a floor E_e of the lazy tower (reversed), depending on the length of the floor.

On the bases, we set:

$$\mu_{g|E_g} := m_p \otimes \lambda;$$
$$\mu_{\ell|E_\ell} := m_p \otimes \lambda.$$

On the floor $E_{g,\omega_{-n},\ldots,\omega_{-1}}$, the measure μ_g is defined by

$$\mu_{g|E_{g,\omega_{-n},\dots,\omega_{-1}}} := \frac{1}{\beta^n} m_p([\omega_{-n},\dots,\omega_{-1}]_0^{n-1}) m_p \otimes \lambda,$$

where $[\omega_{-n}, \ldots, \omega_{-1}]_0^{n-1}$ is the cylinder of Ω containing the sequences ω whose *n* first terms are $(\omega_{-n}, \ldots, \omega_{-1})$.

Likewise on the floor $E_{\ell,\omega_{-n},\ldots,\omega_{-1}}$, the measure μ_{ℓ} is defined by

$$\mu_{\ell|E_{\ell,\omega_{-n},\dots,\omega_{-1}}} := \frac{1}{\beta^n} m_p([\omega_{-n},\dots,\omega_{-1}]_0^{n-1}) m_p \otimes \lambda.$$

On their respective tower, the measure μ_g and μ_ℓ are preserved by \mathcal{K} when it goes up in the tower, but we look for a measure μ on X, globally preserved by \mathcal{K} . The following theorem describes the situation in a more abstract framework, and shows the existence of the measure μ .

THEOREM 1. Let (X, \mathcal{B}) be a standard Borel space, \mathcal{K} a transformation on X, and \mathcal{G} and \mathcal{L} two disjoint subsets of X such that the following hold.

- (1) $X = \mathcal{G} \sqcup \mathcal{L}.$
- (2) There exist two sets $E_g \subset \mathcal{G}$ and $E_\ell \subset \mathcal{L}$ such that

$$\mathcal{K}(\mathcal{G}) \subset \mathcal{G} \cup E_{\ell} \tag{1}$$

and

$$\mathcal{K}(\mathcal{L}) \subset \mathcal{L} \cup E_g.$$

(3) There exist two finite measures μ_g and μ_ℓ with respective support in \mathcal{G} and \mathcal{L} such that

$$\mu_g(E_g) = \mu_\ell(E_\ell) = 1$$

and for any measurable A of $\mathcal{G} \setminus E_g$, we have

$$\mu_g(\mathcal{K}^{-1}(A)) = \mu_g(A), \tag{2}$$

and for any measurable A of $\mathcal{L} \setminus E_{\ell}$, we have

$$\mu_{\ell}(\mathcal{K}^{-1}(A)) = \mu_{\ell}(A).$$

(4) There exists positive constants g₀, l₀, g₁, l₁ such that for any measurable A of X included in E_g, we have

$$\mu_g(\mathcal{K}^{-1}(A) \cap \mathcal{G}) = g_0 \mu_g(A) \tag{3}$$

and

$$\mu_{\ell}(\mathcal{K}^{-1}(A) \cap \mathcal{L})) = \ell_0 \mu_g(A), \tag{4}$$

and for any measurable A of X included in E_{ℓ} , we have

$$\mu_g(\mathcal{K}^{-1}(A) \cap \mathcal{G}) = g_1 \mu_\ell(A)$$

and

$$\mu_{\ell}(\mathcal{K}^{-1}(A) \cap \mathcal{L})) = \ell_1 \mu_{\ell}(A).$$

Then there exist a unique probability measure μ on X, linear combination of μ_g and μ_ℓ , such that μ is \mathcal{K} -invariant.

Proof. Let μ be a measure of the form

$$\mu = \Gamma_g \mu_g + \Gamma_\ell \mu_\ell,$$

with Γ_g and Γ_ℓ two positive constants. We wish to determine the values of Γ_g and Γ_ℓ such that μ is a \mathcal{K} -invariant probability measure on X.

Let *A* be a measurable subset of *X* included in $\mathcal{G} \setminus E_g$. We have $\mu(A) = \Gamma_g \mu_g(A)$. Then, by the assumptions in equations (1) and (2), we have

$$\mu(\mathcal{K}^{-1}(A)) = \Gamma_g \mu_g(\mathcal{K}^{-1}(A)) = \Gamma_g \mu_g(A) = \mu(A).$$

Similarly, for any measurable *A* of *X* included in $\mathcal{L} \setminus E_{\ell}$:

$$\mu(\mathcal{K}^{-1}(A)) = \Gamma_{\ell}\mu_{\ell}(\mathcal{K}^{-1}(A)) = \Gamma_{\ell}\mu_{\ell}(A) = \mu(A).$$

Let A be a measurable subset of X included in E_g . We have, by the assumptions in equations (3) and (4):

$$\mu(\mathcal{K}^{-1}(A)) = \mu(\mathcal{K}^{-1}(A) \cap \mathcal{G}) + \mu(\mathcal{K}^{-1}(A) \cap \mathcal{L})$$

= $\Gamma_g \mu_g(\mathcal{K}^{-1}(A) \cap \mathcal{G}) + \Gamma_\ell \mu_\ell(\mathcal{K}^{-1}(A) \cap \mathcal{L})$
= $g_0 \Gamma_g \mu_g(A) + \ell_0 \Gamma_\ell \mu_g(A)$
= $(g_0 \Gamma_g + \ell_0 \Gamma_\ell) \mu_g(A).$

Therefore, $\mu(\mathcal{K}^{-1}(A)) = \mu(A)$ if and only if

$$g_0\Gamma_g + \ell_0\Gamma_\ell = \Gamma_g. \tag{5}$$

Likewise, if A is a measurable subset of X included in E_{ℓ} , $\mu(\mathcal{K}^{-1}(A)) = \mu(A)$ if and only if

$$g_1 \Gamma_g + \ell_1 \Gamma_\ell = \Gamma_\ell. \tag{6}$$

We need both the equations (5) and (6) to be satisfied. In other words, we solve the system:

$$\begin{cases} (g_0 - 1)\Gamma_g + \ell_0 \Gamma_\ell = 0, \\ g_1 \Gamma_g + (\ell_1 - 1)\Gamma_\ell = 0. \end{cases}$$
(7)

From the assumption in equation (1), we have

$$\mu_g(\mathcal{G}) = \mu_g(\mathcal{K}^{-1}(\mathcal{G} \setminus E_g)) + \mu_g(\mathcal{K}^{-1}(E_g) \cap \mathcal{G}) + \mu_g(\mathcal{K}^{-1}(E_\ell) \cap \mathcal{G})$$

and since $\mu_g(\mathcal{K}^{-1}(\mathcal{G} \setminus E_g)) = \mu_g(\mathcal{G} \setminus E_g)$, we get

$$\mu_g(E_g) = \mu_g(\mathcal{K}^{-1}(E_g) \cap \mathcal{G}) + \mu_g(\mathcal{K}^{-1}(E_\ell) \cap \mathcal{G}).$$

Since $\mu_g(E_g) = 1$, we deduce that

$$g_0 + g_1 = 1.$$

Similarly,

$$\ell_0 + \ell_1 = 1.$$

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Therefore, the system in equation (7) can be reduced to

$$\Gamma_g = \frac{l_0}{g_1} \Gamma_\ell$$

For μ to be a probability measure, we must also have

$$\Gamma_g \mu_g(\mathcal{G}) + \Gamma_\ell \mu_\ell(\mathcal{L}) = 1.$$

We finally obtain the (positive) values

$$\Gamma_g = \frac{\ell_0}{\ell_0 \mu_g(\mathcal{G}) + g_1 \mu_\ell(\mathcal{L})}$$

and

$$\Gamma_{\ell} = \frac{g_1}{\ell_0 \mu_g(\mathcal{G}) + g_1 \mu_{\ell}(\mathcal{L})}$$

With this choice of constants, the measure μ is a \mathcal{K} -invariant probability measure on X. \Box

Let us prove that the extension built in the previous section satisfies the assumptions of Theorem 1.

- (1) The greedy and lazy towers are disjoint.
- (2) Under the action of \mathcal{K} , any element of \mathcal{G} can go up in the tower \mathcal{G} , or go back to the base E_g , or be sent in the lazy base E_ℓ , and the situation is symmetric for the elements of the lazy tower.
- (3) The measures μ_g and μ_ℓ have their respective support in G and L, and μ_g(E_g) = μ_ℓ(E_ℓ) = 1. Let A be a measurable subset of E_{g,ω_{-n},...,ω₋₁} of the form A := [u]₀^{|u|-1} × {(g, ω_{-n}, ..., ω₋₁)} × [a, b], where u is a finite sequence of elements of {g, ℓ}, and |u| denotes the length of this sequence. Let us prove that A and K⁻¹(A) have the same measure. We have

$$\mu_g(\mathcal{K}^{-1}(A)) = \frac{1}{\beta^{n-1}} m_p([\omega_{-n}, \dots, \omega_{-2}]_0^{n-2}) m_p([\omega_{-1} \cdot u]_0^{|u|}) \frac{1}{\beta} (b-a)$$
$$= \frac{1}{\beta^n} m_p([\omega_{-n}, \dots, \omega_{-1}]_0^{n-1}) m_p([u]_0^{|u|-1}) (b-a).$$

We then deduce that for any measurable set A of $\mathcal{G} \setminus E_g$, we have $\mu_g(\mathcal{K}^{-1}(A)) = \mu_g(A)$. Similarly, for any measurable set A of $\mathcal{L} \setminus E_\ell$, we have $\mu_\ell(\mathcal{K}^{-1}(A)) = \mu_\ell(A)$.

(4) We set

$$g_{0} := \mu_{g}(\mathcal{K}^{-1}(E_{g}) \cap \mathcal{G}) = \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{\substack{v \in \{g,\ell\}^{n} \\ T_{v}(1^{+}) \ge 1/\beta}} m_{p}([v \cdot g]_{0}^{n}),$$
$$\ell_{0} := \mu_{\ell}(\mathcal{K}^{-1}(E_{g}) \cap \mathcal{L}) = \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{\substack{w \in \{g,\ell\}^{n} \\ T_{w}(s(1)^{-}) \le 1/\beta}} m_{p}([w \cdot g]_{0}^{n}).$$

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$$g_{1} := \mu_{g}(\mathcal{K}^{-1}(E_{\ell}) \cap \mathcal{G}) = \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{\substack{v \in \{g,\ell\}^{n} \\ T_{v}(1^{+}) \ge 1/\beta(\beta-1)}} m_{p}([v \cdot \ell]_{0}^{n}),$$
$$l_{1} := \mu_{\ell}(\mathcal{K}^{-1}(E_{\ell}) \cap \mathcal{L}) = \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{\substack{w \in \{g,\ell\}^{n} \\ T_{w}(s(1)^{-}) \le 1/\beta(\beta-1)}} m_{p}([w \cdot \ell]_{0}^{n})$$

Let $A \subset E_g$ of the form $A = [u]_0^{|u|-1} \times \{g\} \times]a$; b[, with $0 \le a \le b \le 1$. Then,

$$\mathcal{K}^{-1}(A) \cap \mathcal{G} = \bigsqcup_{n \ge 0} \bigsqcup_{\substack{v \in \{g,\ell\}^n \\ T_v(1^+) \ge 1/\beta}} E_{g \cdot v} \cap \left\{ x \in \left[\frac{a}{\beta}; \frac{b}{\beta} \right]; \omega \in [g \cdot u]_0^{|u|} \right\}.$$

Therefore,

$$\mu_{g}(\mathcal{K}^{-1}(A) \cap \mathcal{G}) = \sum_{n \ge 0} \sum_{\substack{v \in \{g,\ell\}^{n} \\ T_{v}(1^{+}) \ge 1/\beta}} \frac{1}{\beta^{n}} \times m_{p}([v]_{0}^{n-1}) \times m_{p}([g \cdot u]_{0}^{|u|}) \times \frac{1}{\beta}(b-a)$$
$$= g_{0}\mu_{g}(A).$$

We obtain the three other formulas with a similar computation.

Therefore, by setting $\Gamma_g := \ell_0/(\ell_0\mu_g(\mathcal{G}) + g_1\mu_\ell(\mathcal{L}))$ and $\Gamma_\ell := g_1/(\ell_0\mu_g(\mathcal{G}) + g_1\mu_\ell(\mathcal{L}))$, the measure $\mu := \Gamma_g\mu_g + \Gamma_\ell\mu_\ell$ is a \mathcal{K} -invariant probability measure on X.

On each floor of the two towers, the measure μ can be represented as the product of m_p on Ω and a multiple of λ . Thus, we get the following proposition.

PROPOSITION 2. By projecting the measure μ on $\Omega \times I_{\beta}$, we obtain the measure $m_p \otimes \mu_p$, where μ_p has for density

$$\rho_{p}(x) := \sum_{n=0}^{+\infty} \frac{1}{\beta^{n}} \bigg(\Gamma_{g} \sum_{\omega_{1},...,\omega_{n} \in \{g,\ell\}^{n}} m_{p}([\omega_{1},\ldots,\omega_{n}]_{0}^{n-1}) \mathbb{1}_{[0,T_{\omega_{1},...,\omega_{n}}(1^{+})]}(x) + \Gamma_{\ell} \sum_{\omega_{1},...,\omega_{n} \in \{g,\ell\}^{n}} m_{p}([\omega_{1},\ldots,\omega_{n}]_{0}^{n-1}) \mathbb{1}_{[T_{\omega_{1},...,\omega_{n}}(s(1)^{-});1/(\beta-1)]}(x) \bigg).$$
(8)

As a consequence, the natural projection $\pi : X \to \Omega \times I_{\beta}$ is a factor map from the system (X, μ, \mathcal{K}) to $(\Omega \times I_{\beta}, m_p \otimes \mu_p, K_{\beta})$.

2.4. Properties of the extension. The aim of this section is to prove that the system (X, μ, \mathcal{K}) is ergodic. For that, we first study the properties of \mathcal{K} when following paths on the towers. We then study the induced transformation of \mathcal{K} on the greedy base E_g , and prove that the induced system is Bernoulli. From this, we derive the ergodicity of the system (X, μ, \mathcal{K}) , which implies the ergodicity of the initial random system and the uniqueness of μ_p (as an absolutely continuous probability measure on I_β such that $m_p \otimes \mu_p$ is K_β -invariant).

To study the behavior of the dynamics along paths in the towers we first define a class of functions that will be useful in the following definition.

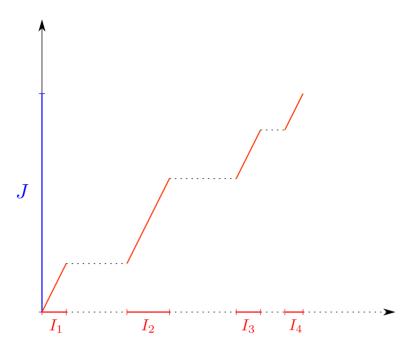


FIGURE 9. Function from the class \mathfrak{F}_n .

Definition 3. Let $n \in \mathbb{N}$. We define the class of functions \mathfrak{F}_n (see Figure 9) as the set of functions f satisfying, up to a set of zero measure:

- there exists a finite number of disjoint intervals I_1, \ldots, I_r and a sub-interval J of I_β such that $f: I_1 \sqcup \cdots \sqcup I_r \to J$;
- the function *f* is non-decreasing on its definition domain;
- on each interval I_i , f is a linear map with slope β^n ;
- $f(\bigsqcup_{i=1}^r I_i) = \bigsqcup_{i=1}^r f(I_i) = J.$

In particular, the intervals (I_i) are in the same order as the intervals $(f(I_i))$, and we have

$$\lambda(I_1 \sqcup \cdots \sqcup I_r) = \frac{1}{\beta^n} \lambda(J).$$
(9)

LEMMA 4. Let $n_1, n_2 \in \mathbb{N}$ and $f \in \mathfrak{F}_{n_1}$, $g \in \mathfrak{F}_{n_2}$ such that the definition domain D_f of f is included in the image of g. Then, $f \circ g_{|g^{-1}(D_f)} \in \mathfrak{F}_{n_1+n_2}$.

Proof. We denote $g: I_1 \sqcup \cdots \sqcup I_{r_1} \to J$ and $f: J_1 \sqcup \cdots \sqcup J_{r_2} \to L$ and $h = f \circ g_{|g^{-1}(D_f)}$. Then, $g^{-1}(D_f)$ is the disjoint union of the intervals $g^{-1}(J_i) \cap I_j$ (when this set is not empty). Therefore, h is defined on $\bigsqcup_{i,j} g^{-1}(J_i) \cap I_j$ with values in L. As the composition of non-decreasing functions, h is non-decreasing on its definition domain. Moreover, h is the composition of two linear maps of slope β^{n_1} and β^{n_2} on each non-empty $g^{-1}(J_i) \cap I_j$, so it is a linear map of slope $\beta^{n_1+n_2}$. Finally, for all $i \in \{1, \ldots, r_2\}$, we have

$$\bigsqcup_{j} g(g^{-1}(J_i) \cap I_j) = J_i,$$

and thus

$$\bigsqcup_{i,j} h(g^{-1}(J_i) \cap I_j) = L.$$

We denote by $\mathcal{E} := \{g, \ell\}^* \setminus \{\emptyset\}$ the set of non-empty finite sequences of $\{g, \ell\}$, in other words, the set of labels. For $e \in \mathcal{E}$, we denote by e(0) the first term of e and e(-1) its last term.

We define the oriented graph G on \mathcal{E} describing the set of admissible sequences of labels. A label $e \in \mathcal{E}$ can be followed by the label $e' \in \mathcal{E}$ (we then denote $e \to e'$) if $\mu(\mathcal{K}^{-1}(E_{e'}) \cap E_e) > 0$. In other words, we have $e \to e'$ if the map \mathcal{K} sends a subset of positive measure of the floor E_e onto the floor $E_{e'}$.

We use the new class of functions \mathfrak{F}_n to describe the dynamics along paths in the two towers.

PROPOSITION 5. Let $C = (e_0 = e, e_1, ..., e_n = e')$ $(n \ge 1)$ be a finite sequence of labels of \mathcal{E} corresponding to a path in the graph G. We have

$$E_C := \bigcap_{k=0}^n \mathcal{K}^{-k} E_{e_k} = [e_1(-1), \dots, e_n(-1)]_0^{n-1} \times \{e\} \times J_C$$

with:

(1) J_C is a finite union of disjoint intervals I_1, \ldots, I_r ;

(2) $T_{e_1(-1),...,e_n(-1)|J_C} \in \mathfrak{F}_n.$

In particular, the intervals (I_i) are in the same order as the intervals $(T_{e_1(-1),\ldots,e_n(-1)}(I_i)), \lambda(J_C) = 1/\beta^n \lambda(I_{e'})$ and the mapping $\mathcal{K}^n : E_C \to E_{e'}$ is a bijection.

Proof. We prove each point by induction on the length of the path *C*. Suppose that C = (e, e') is an edge in the graph of transitions of \mathcal{E} . Set $\omega_0 = e'(-1)$. Without loss of generality, we can suppose that e(0) = g. Finally, set $I_e = [0, t]$, with $t \in [0, 1/(\beta - 1)]$.

- If $\omega_0 = g$, then $E_C = [g]_0 \times \{e\} \times J_C$ with - if $t < 1/\beta$, then $J_C = I_e$, - if $t \ge 1/\beta$, either e' = (g), then $J_C = [0, 1/\beta[$, or $e' = e \cdot g$ then $J_C = [1/\beta, t]$.
- If $\omega_0 = \ell$, then $E_C = [\ell]_0 \times \{e\} \times J_C$ with
 - if $t < 1/\beta(\beta 1)$, then $J_C = I_e$,
 - if $t \ge 1/\beta(\beta 1)$, either $e' = (\ell)$ then $J_C = [s(1)/\beta; 1/\beta(\beta 1)]$, or $e' = e \cdot \ell$ then $J_C = [0, s(1)/\beta] \sqcup [1/\beta(\beta - 1), t]$.

In each case, J_C is a finite union of disjoint intervals. Moreover, $T_{\omega_0}: J_C \to I_{e'}$ is a function of \mathfrak{F}_1 .

Suppose now that $C = (e = e_0, e_1, ..., e_{n+1} = e')$. Set $C' = (e_1, ..., e_{n+1})$. The path C' is a path of length n in the graph of transitions of \mathcal{E} , such that $E_{C'}$ satisfies the conclusions of the proposition.

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- $E_{C'} = [e_2(-1), \ldots, e_{n+1}(-1)]_0^{n-1} \times \{e_1\} \times J_{C'}$ where $J_{C'}$ is a finite union of disjoint intervals.
- $T_{e_2(-1),...,e_{n+1}(-1)|J_{C'}} \in \mathfrak{F}_n.$

We then have $E_C = E_e \cap \mathcal{K}^{-1}(E_{C'})$.

The component in Ω is indeed $[e_1(-1), e_2(-1), ..., e_{n+1}(-1)]_0^n$.

Set $g = T_{e_1(-1)|J_{e_0,e_1}}$ and $f = T_{e_2(-1),\dots,e_{n+1}(-1)|J_{C'}}$. We have $g \in \mathfrak{F}_1$ and $f \in \mathfrak{F}_n$. By setting $J_C := g^{-1}(J_{C'})$, we do have that J_C is a finite union of disjoint intervals and

$$E_C = [e_1(-1), e_2(-1), \dots, e_{n+1}(-1)]_0^n \times \{e\} \times J_C.$$

Furthermore, $T_{e_1(-1),e_2(-1),...,e_{n+1}(-1)|J_C} = f \circ g_{|J_C}$. Applying Lemma 4, the function $T_{e_1(-1),e_2(-1),...,e_{n+1}(-1)|J_C}$ is in \mathfrak{F}_{n+1} .

We now study the induced transformation of \mathcal{K} on the base E_g . We denote by \mathcal{K}_g this induced transformation. The measure $\mu(.|E_g)$ is the measure μ_g , and is preserved by the induced map \mathcal{K}_g [1, Proposition 3.6.1, pp. 58]. In the following, we identify the base E_g with $\Omega \times [0, 1]$.

We consider the set \mathscr{C} of possible paths for a first return to the base E_g . In other words, \mathscr{C} is the set of paths of the form $C = (e_0 = (g), e_1, \ldots, e_n = (g))$ (where |C| := n is the length of the path *C*), with $e_k \neq (g)$ for $1 \leq k \leq n - 1$ such that

$$\mu\bigg(\bigcap_{k=0}^n \mathcal{K}^{-k}(E_{e_k})\bigg) > 0.$$

We then define the set \mathcal{P} of subsets of E_g by

$$\mathcal{P} := \{ E_C, C \in \mathscr{C} \}.$$

LEMMA 6. The set \mathcal{P} is a countable partition of E_g , up to a set of zero measure.

Proof. If we consider two distinct paths $C, C' \in \mathcal{C}$, then it is clear that E_C and $E_{C'}$ are disjoint. Additionally, since μ_g -almost every a in E_g has a finite return time, there exists $C \in \mathcal{C}$ such that $a \in E_C$.

For almost every $a \in E_g$, we call \mathcal{P} -name of a the unique sequence $(P_n)_{n \ge 0}$ of atoms of \mathcal{P} such that for any $n \ge 0$, $K_g^n(a) \in P_n$.

PROPOSITION 7. The sequence of partitions $(\mathcal{K}_g^{-n}(\mathcal{P}))_{n \ge 0}$ is independent for the measure μ_g , that is, for any integer $n \in \mathbb{N}$, the partitions $\mathcal{K}_g^{-n}\mathcal{P}$ and $\bigvee_{0 \le k \le n-1} \mathcal{K}_g^{-k}\mathcal{P}$ are independent.

Proof. Consider a path

$$C = (e_0 = (g), e_1, \dots, e_N = (g))$$

starting and ending with the label of the greedy base E_g . Then, this path is a unique concatenation of paths $C_0, \ldots, C_k \in \mathcal{C}$, and

$$E_C = \bigcap_{i=0}^k \mathcal{K}_g^{-i} E_{C_i} = \bigcap_{i=0}^N \mathcal{K}^{-i}(E_{e_i}) = [e_1(-1), \dots, e_N(-1)]_0^{N-1} \times J_C$$

with $\lambda(J_C) = (1/\beta^N)\lambda(I_{(g)}) = 1/\beta^N$, from Proposition 5. Then, we have

$$\mu_g(E_C) = m_p([e_1(-1), \dots, e_N(-1)]_0^{N-1}) \times \frac{1}{\beta^N},$$

and we have a similar equality for all C_j , $0 \le j \le k$. Therefore, it is clear that

$$\mu_g(E_C) = \prod_{j=0}^k \mu_g(E_{C_j}).$$

From Proposition 5, each atom of \mathcal{P} is of the form $[\omega_0, \ldots, \omega_{r-1}] \times J_C$. Likewise, for each $N \ge 0$, each atom of $\bigvee_{i=0}^N \mathcal{K}_g^{-i} \mathcal{P}$ is again of this form.

Definition 8. Let $N \in \mathbb{N} \cup \{0\}$ and $(\omega, x) \in E_g$ such that there exists an atom of $\bigvee_{j=0}^{N} \mathcal{K}_g^{-j} \mathcal{P}$ containing (ω, x) . We denote *C* the associated path. Then, we define the set $A_N(\omega, x)$ included in the interval [0, 1] by

$$A_N(\omega, x) := J_C.$$

LEMMA 9. For almost every $(\omega, x) \in E_g$, for all $n \in \mathbb{N} \cup \{0\}$ and $m \in \mathbb{N}$, there exists $N \ge n$ and $t \in [0, 1]$ such that

$$A_{N+m}(\omega, x) = A_N(\omega, x) \cap [0, t].$$

Proof. Let $P_0 = [g]_0 \times [0, 1/\beta[$. The set P_0 is an atom of the partition \mathcal{P} , of positive measure, corresponding to the set of points returning to the base E_g in one step. The process that maps a point in E_g to its \mathcal{P} -name (defined for almost every $(\omega, x) \in E_g$) is a Bernoulli process. Therefore, for almost every $(\omega, x) \in E_g$, the orbit of (ω, x) under \mathcal{K}_g meets infinitely often the atom P_0 . Even better, for any $m \in \mathbb{N}$, the orbit of (ω, x) meets the atom $P_0 m$ times in a row, infinitely often.

Let $n \in \mathbb{N} \cup \{0\}$. For such a point (ω, x) , let $N \ge n$ be an integer such that

$$\mathcal{K}_g^N(\omega, x), \ldots, \mathcal{K}_g^{N+m-1}(\omega, x) \in P_0.$$

Observe that the atom of $\bigvee_{j=0}^{N} K_g^{-j} \mathcal{P}$ containing (ω, x) is of the form

$$[\omega_0,\ldots,\omega_{r-1}]_0^{r-1}\times A_N(\omega,x).$$

If we note $A_N(\omega, x) = I_1 \sqcup \cdots \sqcup I_m$, Proposition 5 implies

$$[0,1] = T_{\omega_0,\ldots,\omega_{r-1}}I_1 \sqcup \cdots \sqcup T_{\omega_0,\ldots,\omega_{r-1}}I_m,$$

where $T_{\omega_0,...,\omega_{r-1}}I_i$ is an interval for every $1 \le i \le m$, and the order of the intervals $(T_{\omega_0,...,\omega_{r-1}}I_i)$ is the same as the order of the intervals (I_i) .

Since $\mathcal{K}_{g}^{N}(\omega, x) \in P_{0}$, we have

$$A_{N+1}(\omega, x) = \left\{ y \in A_N(\omega, x), T_{\omega_0, \dots, \omega_{r-1}}(y) \in \left[0, \frac{1}{\beta}\right] \right\}.$$

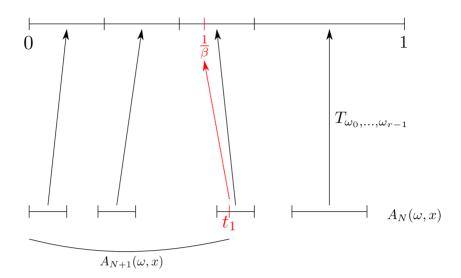


FIGURE 10. Passage from $A_N(\omega, x)$ to $A_{N+1}(\omega, x)$.

Then, there exists $t_1 \in [0, 1]$ such that $T_{\omega_0, \dots, \omega_{r-1}}(A_N(\omega, x) \cap [0, t_1]) = [0, 1/\beta[$. See Figure 10. Hence,

$$A_{N+1}(\omega, x) = A_N(\omega, x) \cap [0, t_1].$$

Finally, since $\mathcal{K}_g^N(\omega, x), \ldots, \mathcal{K}_g^{N+m-1}(\omega, x) \in P_0$, we can apply the same reasoning *m* times, which implies the existence of $t \in [0, 1]$ such that

$$A_{N+m}(\omega, x) = A_N(\omega, x) \cap [0, t].$$

This real number t is such that $\lambda(A_N(\omega, x) \cap [0, t]) = (1/\beta^m)\lambda(A_N(\omega, x)).$

Let \mathcal{F} be the factor σ -algebra $\mathcal{F} := \bigvee_{j=0}^{+\infty} \mathcal{K}_g^{-j} \mathcal{P}$. The partition \mathcal{P} provides a Bernoulli process, and we want to show that this process generates the induced system by proving the following proposition.

PROPOSITION 10. The σ -algebra \mathcal{F} is (up to zero measure sets) the Borel σ -algebra on E_g . In other words, the partition \mathcal{P} is a generator of the induced system, which is isomorphic to a one-sided Bernoulli shift.

LEMMA 11. \mathcal{F} contains the σ -algebra generated by the ω -component.

Proof. Let $N \in \mathbb{N}$. As said earlier, every atom of $\bigvee_{j=0}^{N} \mathcal{K}_{g}^{-j} \mathcal{P}$ is of the form $[\omega_{0}, \ldots, \omega_{r-1}]_{0}^{r-1} \times J_{C}$, with $r \ge N$. Therefore, for any $N \in \mathbb{N}$, $(\omega_{0}, \ldots, \omega_{N-1})$ is \mathcal{F} -measurable.

To prove Proposition 10, it suffices to show that for any continuous function φ : $E_g \to \mathbb{R}_+$, the conditional expectation $E[\varphi|\mathcal{F}]$ satisfies

$$E[\varphi|\mathcal{F}] = \varphi, \quad \mu_g \text{ almost surely.}$$

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The existence of the regular conditional probability implies that for almost every $(\omega, x) \in E_g$, there exists a measure $\mu_{(\omega,x)}$ on E_g such that for any positive continuous function $\varphi : E_g \to \mathbb{R}_+$,

$$E[\varphi|\mathcal{F}](\omega, x) = \int \varphi \, d\mu_{(\omega, x)}$$

We want to prove that for almost every $(\omega, x) \in E_g$, the measure $\mu_{(\omega,x)}$ equals the Dirac measure on (ω, x) , which leads to the announced result. Lemma 11 implies that, for almost every $(\omega, x) \in E_g$, the measure $\mu_{(\omega,x)}$ is of the form $\delta_{\omega} \otimes \tilde{\mu}_{(\omega,x)}$, where $\tilde{\mu}_{(\omega,x)}$ is a measure on [0, 1]. In the following, we identify the measures $\mu_{(\omega,x)}$ and $\tilde{\mu}_{(\omega,x)}$. Therefore, it remains to prove that $\mu_{(\omega,x)}$ is the Dirac measure on *x*.

Let $\varphi : E_g \to \mathbb{R}_+$ be a continuous function depending only on the real variable *x*, that is, $\varphi(\omega, x) = f(x)$, where *f* is a continuous function from [0, 1] into \mathbb{R}^+ .

By the martingale convergence theorem, we have the almost sure convergence:

$$E\left[\varphi \middle| \bigvee_{j=0}^{N} \mathcal{K}_{g}^{-j} \mathcal{P}\right] \to E[\varphi|\mathcal{F}].$$

Thus, for almost every $(\omega, x) \in E_g$, we have the convergence

$$E\left[\varphi \middle| \bigvee_{j=0}^{N} \mathcal{K}_{g}^{-j} \mathcal{P}\right](\omega, x) \to \int f(y) \, d\mu_{(\omega, x)}(y).$$

With *N* fixed, and for almost every $(\omega, x) \in E_g$, we have

$$E\left[\varphi\right|\bigvee_{j=0}^{N}\mathcal{K}_{g}^{-j}\mathcal{P}](\omega,x) = \frac{1}{\lambda(A_{N}(\omega,x))}\int_{A_{N}(\omega,x)}f(y)\,dy.$$

Set $\mu_{A_N(\omega,x)} := (1/\lambda(A_N(\omega, x)))\lambda_{|A_N(\omega,x)}$ as the normalized Lebesgue measure on $A_N(\omega, x)$. For almost every $(\omega, x) \in E_g$, the sequence of measures $(\mu_{A_N(\omega,x)})_N$ converges weakly to the measure $\mu_{\omega,x}$.

We define the set \overline{E}_g of full measure as the set of points $(\omega, x) \in E_g$ satisfying the following conditions.

• There exists a probability measure $\mu_{(\omega,x)}$ on E_g such that for any continuous function $\varphi: E_g \to \mathbb{R}_+$, we have

$$E[\varphi|\mathcal{F}](\omega, x) = \int \varphi \, d\mu_{(\omega, x)}.$$

- For any continuous function $\varphi: E_g \to \mathbb{R}_+$, $E[\varphi| \bigvee_{j=0}^n \mathcal{K}_g^{-j} \mathcal{P}](\omega, x)$ tends to $E[\varphi|\mathcal{F}](\omega, x)$.
- For any integer N, the point (ω, x) belongs to an atom of the partition $\bigvee_{j=0}^{N} \mathcal{K}_g^{-j} \mathcal{P}$ and satisfies Lemma 9.

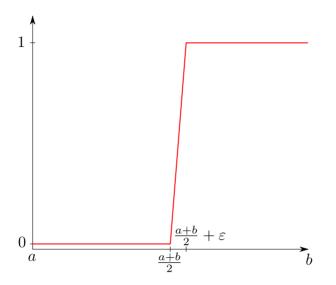


FIGURE 11. Graph of the function f_{ε} .

• For any integer *N* and for any continuous function $\varphi : E_g \to \mathbb{R}_+$ such that there exists a continuous function $f : \mathbb{R} \to \mathbb{R}_+$ such that $\varphi(\omega, x) = f(x)$ for any $(\omega, x) \in E_g$, we have

$$E\left[\varphi\middle|\bigvee_{j=0}^{N}\mathcal{K}_{g}^{-j}\mathcal{P}\right](\omega,x)=\int f(y)\,d\mu_{A_{N}(\omega,x)}(y).$$

LEMMA 12. Let $(\omega, x) \in \overline{E}_g$. Let $[a, b] \subset [0, 1]$ such that $\mu_{(\omega, x)}[a, b] = 1$. Then one of these two propositions is true:

- (i) $\mu_{(\omega,x)}[a, (a+b)/2] = 1;$
- (ii) $\mu_{(\omega,x)}](a+b)/2, b] = 1.$

Proof. Suppose that proposition (ii) does not hold. Set $\eta := \mu_{(\omega,x)}([a, (a+b)/2]) > 0$. Let *m* be a large enough integer so that $1/\beta^m < \eta/2$. Let $\varepsilon > 0$. Let f_{ε} be the continuous function defined on [0, 1] by $f_{\varepsilon} = 0$ on [a, (a+b)/2], $f_{\varepsilon} = 1$ on $[(a+b)/2 + \varepsilon, b]$, and f_{ε} is a linear map on $[(a+b)/2, (a+b)/2 + \varepsilon]$. See Figure 11.

Then,

$$\int f_{\varepsilon} d\mu_{(\omega,x)} \leqslant 1 - \eta$$

Therefore, since the sequence $(\mu_{A_n(\omega,x)})_n$ converges weakly to $\mu_{(\omega,x)}$, there exists $N_0 \in \mathbb{N}$, such that for any $n \ge N_0$,

$$\int f_{\varepsilon} d\mu_{A_n(\omega,x)} \leqslant 1 - \frac{\eta}{2}$$

Thus, for any $n \ge N_0$,

$$\mu_{A_n(\omega,x)}\left(\left[\frac{a+b}{2}+\varepsilon,b\right]\right)\leqslant 1-\frac{\eta}{2}\leqslant 1-\frac{1}{\beta^m}.$$

It then implies that

$$\mu_{A_n(\omega,x)}\left(\left[a,\frac{a+b}{2}+\varepsilon\right]\right) > \frac{1}{\beta^m}$$

Let $n \ge N_0$ such that $K_g^n(\omega, x), \ldots, K_g^{n+m-1}(\omega, x) \in P_0$. Then from Lemma 9, there exists a real number $t \in [0, 1]$ such that $A_{n+m}(\omega, x) = A_n(\omega, x) \cap [0, t]$. Moreover, $\lambda(A_{n+m}(\omega, x)) = (1/\beta^m)\lambda(A_n(\omega, x))$. It implies that $A_{n+m}(\omega, x) \subset [a, (a+b)/2 + \varepsilon[$, and hence, for any $n' \ge n + m$,

$$\mu_{A_{n'}(\omega,x)}\left(\left[\frac{a+b}{2}+\varepsilon,b\right]\right)=0.$$

Furthermore, we have

$$\int f_{2\varepsilon} d\mu_{A_n(\omega,x)} \to \int f_{2\varepsilon} d\mu_{(\omega,x)}$$

We then get that for any $\varepsilon > 0$, $\mu_{(\omega,x)}([(a+b)/2+2\varepsilon, b]) = 0$, and hence $\mu_{(\omega,x)}(](a+b)/2, b]) = 0$.

By dichotomy, we deduce that for any $(\omega, x) \in E_g$, the measure $\mu_{(\omega,x)}$ is the Dirac measure on *x*, which proves Proposition 10. In particular, the induced transformation \mathcal{K}_g is isomorphic to a one-sided Bernoulli shift.

LEMMA 13. There exists a floor of the greedy tower with length greater than or equal to $1/\beta(\beta - 1)$. In particular, it is always possible to go from the tower \mathcal{G} to the tower \mathcal{L} under \mathcal{K} .

Proof. We set

$$n_0 = \min\left\{n \in \mathbb{N} \cup \{0\}, T_{\ell}^n(1) \ge \frac{1}{\beta(\beta - 1)}\right\}.$$

Such an integer n_0 exists for any $\beta > 1$: if $1 < 1/\beta(\beta - 1)$, then $T_{\ell}(1) = \beta \times 1$. Iterating T_{ℓ} from 1 is just a multiplication by β , as long as the images are less than $1/\beta(\beta - 1)$. Thus, the length of the floor $E_{g\ell^{n_0}}$ is greater than or equal to $1/\beta(\beta - 1)$. Then we have

$$\mathcal{K}\left(E_{g\ell^{n_0}} \cap \left\{\omega_0 = \ell, x \in \left[\frac{s(1)}{\beta}, \frac{1}{\beta(\beta-1)}\right]\right\}\right) = E_{\ell}.$$

THEOREM 14. The system (X, μ, \mathcal{K}) is ergodic.

Proof. The induced transformation of \mathcal{K} on E_g is Bernoulli, and hence ergodic. Therefore, it suffices to prove that, up to a zero measure set,

$$\bigcup_{n \ge 1} \mathcal{K}^{-n}(E_g) = X.$$

In other words, we want to prove that almost every point in X reaches the base E_g in a finite number of iterations of \mathcal{K} .

We first prove that almost every point in X can be reached from the base E_g . By construction of \mathcal{K} , every floor of the greedy tower is the image of a part of E_g by a power of \mathcal{K} . More precisely, for any $e \in \mathcal{E}$ such that e(0) = g, we have

$$E_{e} = \mathcal{K}^{|e|-1} \left(\bigcap_{k=0}^{|e|-1} \mathcal{K}^{-k}(E_{e(0),\dots,e(k)}) \right)$$

Similarly, every floor of the lazy tower is the image by a power of \mathcal{K} of a part of the base E_{ℓ} . Finally, from Lemma 13, we can pass from the greedy tower to the lazy tower.

Additionally, almost every point in the base E_g returns to E_g in a finite number of iterations of \mathcal{K} (from Poincaré recurrence theorem). Denote by \mathcal{N}_g the set of points of E_g which do not return to E_g . Then \mathcal{N}_g is a zero measure set. Therefore, the set $\bigcup_{n \ge 0} \mathcal{K}^n(\mathcal{N}_g)$ is measurable, and of measure zero. Indeed, let $n \in \mathbb{N} \cup \{0\}$. Then

$$\mathcal{K}^n(\mathcal{N}_g) = \bigcup_{|C|=n} \mathcal{K}^n(\mathcal{N}_g \cap E_C).$$

For any path *C* of length *n*, the transformation \mathcal{K}^n is invertible on $\mathcal{N}_g \cap E_C$ from Proposition 5. Therefore, $\mu(K^n(\mathcal{N}_g \cap E_C)) = 0$ and $\mu(\mathcal{K}^n(\mathcal{N}_g)) = 0$.

Let $a \in X \setminus \bigcup_{n \ge 0} \mathcal{K}^n(\mathcal{N}_g)$. There exist $a_0 \in E_g$ and $n \in \mathbb{N} \cup \{0\}$ such that $\mathcal{K}^n(a_0) = a$. Moreover, $a_0 \notin \bigcup_{n \ge 0} \mathcal{K}^n(\mathcal{N}_g)$ so there exists a path $C_0^a \in \mathcal{C}$ such that $a_0 \in P_{C_0^a}$. It implies that *a* will return to the greedy base in a finite number of iterations of \mathcal{K} (following the path C_0^a).

Since the system $(\Omega \times I_{\beta}, m_p \otimes \mu_p, K_{\beta})$ is a factor of the system (X, μ, \mathcal{K}) , we have the following corollary.

COROLLARY 15. The system $(\Omega \times I_{\beta}, m_p \otimes \mu_p, K_{\beta})$ is ergodic.

COROLLARY 16. The measure μ_p is the unique absolutely continuous probability measure on I_β such that $m_p \otimes \mu_p$ is K_β -invariant.

Proof. From Lemma 13, for any $1 < \beta < 2$, there exists a floor of the greedy tower whose length is greater than or equal to $1/\beta(\beta - 1)$. Symmetrically, there exists a floor of the lazy tower whose left limit is less than or equal to $1/\beta$. It implies that the support of μ_p is the full interval I_{β} , and that $m_p \otimes \mu_p$ is equivalent to $m_p \otimes \lambda$.

Let v be an absolutely continuous probability measure (with respect to λ , and so with respect to μ_p) such that $m_p \otimes v$ is K_β -invariant. Then the measure $m_p \otimes v$ is absolutely continuous with respect to the ergodic measure $m_p \otimes \mu_p$, and hence $m_p \otimes v = m_p \otimes \mu_p$ and $v = \mu_p$.

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In [16], Suzuki proves that the density of μ_p is proportional to the function

$$B^{0}(\beta, p) \bigg(\mathbb{1}_{[0,1]} + \sum_{n=1}^{\infty} \frac{1}{\beta^{n}} \sum_{v \in \{g,\ell\}^{v}} m_{p}([v]_{1}^{n}) \mathbb{1}_{[0,T_{v}(1)]} \bigg) + B^{0}(\beta, 1-p) \bigg(\mathbb{1}_{[s(1),1/(\beta-1)]} + \sum_{n=1}^{\infty} \frac{1}{\beta^{n}} \sum_{w \in \{g,\ell\}^{n}} m_{p}([w])_{1}^{n} \mathbb{1}_{[T_{w}(s(1)),1/(\beta-1)]} \bigg), \quad (10)$$

where

$$B^{0}(\beta, p) = p\left(1 - \sum_{n=1}^{\infty} \frac{1}{\beta^{n}} \sum_{w \in \{g,\ell\}^{n-1}} \mathbb{1}_{S}(T_{w}(s(1)))m_{p}([w,\ell]_{1}^{n})\right)$$

By uniqueness of μ_p , we know that the density in equation (8) and the expression obtained by Suzuki (10) are equal. This result can be easily proved in the case where β is not a root of a polynomial of the form $X^{n_0} - X^{n_1} - \cdots - X^{n_k} - 1$ with $n_0 > n_1 > \cdots > n_k$ (which is equivalent to say that 1 has no finite expansion in base β).

Indeed, we prove that, in this case, $\ell_0 = B^0(\beta, p)$. On one hand, we recall that

$$\ell_0 = \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{w \in \{g,\ell\}^n, T_w(s(1)^-) \leqslant 1/\beta} m_p([w \cdot g]_0^n).$$

Since β is not a root of a polynomial of the form $X^{n_0} - X^{n_1} - \cdots - X^{n_k} - 1$, then $T_w(s(1)^-) = T_w(s(1))$ for any $w \in \{g, \ell\}^n$. Moreover, $T_w(s(1)^-) \leq 1/\beta$ is equivalent to $T_w(s(1)) < 1/\beta$. Therefore, after adjusting the index of summation, we have

$$\ell_0 = \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{w \in \{g,\ell\}^n} \mathbb{1}_R(T_w(1)) m_p([w \cdot g]_0^n).$$

On the other hand, after a new reindexing in ℓ_1 , we have

$$\ell_0 = 1 - \ell_1 = 1 - \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{w \in \{g,\ell\}^n} \mathbbm{1}_{S \cup R}(T_w(1)) m_p([w \cdot \ell]_0^n)$$
$$= 1 - \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{w \in \{g,\ell\}^n} \mathbbm{1}_S(T_w(1)) m_p([w \cdot \ell]_0^n) - \frac{1-p}{p} \ell_0.$$

Thus,

$$\ell_0 = p \left(1 - \sum_{n=0}^{+\infty} \frac{1}{\beta^{n+1}} \sum_{w \in \{g,\ell\}^n} \mathbb{1}_S(T_w(1)) m_p([w \cdot \ell]_0^n) \right) = B^0(\beta, p).$$

By symmetry, we also obtain that $g_1 = B^0(\beta, 1 - p)$, and the equality of the two densities follows.

In the case where 1 has a finite expansion in base β , there exists $n \in \mathbb{N}$ and a sequence $w \in \{g, \ell\}^n$ such that $T_w(1) = 0$, and the right limit $T_v(1^+)$ and the value $T_v(1)$ can differ, thus complicating the identification of the two formulas. For the simple case where $\beta = (1 + \sqrt{5})/2$, we refer to the computations in the appendix of [17].

3. The natural extension

3.1. Construction of the new extension. Any point in X which is not in a base has a unique preimage by \mathcal{K} . However, a point in one of the two bases can come from several floors. Indeed, at each return to a base, the map \mathcal{K} 'forgets' from where it comes. This lack of information prevents \mathcal{K} from being invertible.

We wish to construct a natural extension of \mathcal{K} , denoted by $\tilde{\mathcal{K}}$. In some sense, it is the smallest invertible extension of \mathcal{K} . Therefore, to construct $\tilde{\mathcal{K}}$, we need to extend the set X such that the information of the past floors is available. We recall the definition of a natural extension of a system (see for example [2]).

Definition 17. The system (Y, C, ν, F) is a natural extension of the system (X, \mathcal{B}, μ, T) if there exist two sets $X^* \in \mathcal{B}$ and $Y^* \in \mathcal{C}$ such that $\mu(X^*) = \nu(Y^*) = 1$ and a measurable mapping $\pi : Y^* \to X^*$ such that:

- (1) F is a bijection of Y^* ;
- (2) $\mu = \nu \circ \pi^{-1};$
- (3) $\pi \circ F = T \circ \pi$; and
- (4) $\mathcal{C} = \bigvee_{n \ge 0} F^n(\pi^{-1}(\mathcal{B})).$

Moreover, all natural extensions of a system are isomorphic [14]. We now construct the extension $\tilde{\mathcal{K}}$ of \mathcal{K} , and then prove that it is one of its natural extensions. In the initial definition of the towers \mathcal{G} and \mathcal{L} , each floor has a label e. More precisely, a label e is a finite sequence of g and ℓ , where the first term of e describes the present tower and the following terms of e (when there are) describe which transformations have been applied since the base: the label describes the recent past of a point in the floor. Therefore, instead of 'erasing' the label at each return to a base, we will keep track of every past label, so that we can uniquely determine the past orbit of almost any point in the two towers.

We denote by Z the set of left-infinite sequences generated by the graph G, that is,

$$Z := \{ (e_j)_{j \in \mathbb{Z}_{\leq 0}} \in \mathcal{E}^{\mathbb{Z}_{\leq 0}} : \text{for all } j \leq -1, e_j \to e_{j+1} \}$$

A sequence of labels in Z contains in particular the information of the sequence of past transformations. Given a sequence \mathbf{e} in Z, we will denote e_0 its term of index 0.

For $e \in \mathcal{E}$, we set

 $Z_e := \{ \mathbf{e} \in Z : e_0 = e \}.$

We then define the base of the two towers by

$$\tilde{E}_g := \Omega \times Z_{(g)} \times [0, 1]$$

and

$$\tilde{E}_{\ell} := \Omega \times Z_{(\ell)} \times \left[s(1), \frac{1}{\beta - 1} \right].$$

If $e = (g, \omega_{-n}, \ldots, \omega_{-1}), n \in \mathbb{N}$, we set

$$\tilde{E}_e := \Omega \times Z_e \times [0, T_{\omega_{-n}, \dots, \omega_{-1}}(1^+)].$$

If $e = (\ell, \omega_{-n}, \ldots, \omega_{-1}), n \in \mathbb{N}$, we set

$$\tilde{E}_e := \Omega \times Z_e \times \left[T_{\omega_{-n},\dots,\omega_{-1}}(s(1)^-), \frac{1}{\beta - 1} \right].$$

We note

$$\tilde{\mathcal{G}} := \bigsqcup_{e \in \mathcal{E}, e(0) = g} \tilde{E}_e$$

and

$$\tilde{\mathcal{L}} := \bigsqcup_{e \in \mathcal{E}, e(0) = \ell} \tilde{E}_e$$

We finally note $\tilde{X} := \tilde{\mathcal{G}} \sqcup \tilde{\mathcal{L}}$. We set

$$\pi_{\tilde{X},X} : \frac{\tilde{X}}{(\omega, \mathbf{e}, x) \mapsto (\omega, e_0, x)}$$

as the projection from \tilde{X} onto X.

Let $a = (\omega, \mathbf{e}, x) \in \tilde{X}$. Then $\pi_{\tilde{X}, X}(a) = (\omega, e_0, x)$. We have $\mathcal{K}(\omega, e_0, x) = (\sigma(\omega), e', T_{\omega_0}(x))$, where the label e' depends on e_0, ω_0 , and x, according to the previous construction. By construction of \mathcal{K} , we obviously have $e_0 \to e'$ in the graph G.

Applying the new dynamics $\hat{\mathcal{K}}$ on *a* consists in concatenating the new label e' to the sequence **e** (we denote by $\mathbf{e} \cdot e'$ this new sequence), in addition to shifting the sequence ω and applying T_{ω_0} to *x*: we retain every past labels in memory.

In other words, the dynamics $\tilde{\mathcal{K}}$ is defined on \tilde{X} by

$$\tilde{\mathcal{K}}: \begin{array}{c} \tilde{X} \to \tilde{X} \\ (\omega, \mathbf{e}, x) \mapsto (\sigma(\omega), \mathbf{e} \cdot e', T_{\omega_0}(x)) \end{array}$$

where e' is the label associated with the point $\mathcal{K}(\omega, e_0, x) \in X$.

The projection $\pi_{\tilde{X}|X}$ is measurable and we have by construction

$$\mathcal{K} \circ \pi_{\tilde{X},X} = \pi_{\tilde{X},X} \circ \tilde{\mathcal{K}}$$

Denote by \mathcal{B} the Borel σ -algebra on X and by $\tilde{\mathcal{B}}$ the Borel σ -algebra on \tilde{X} . Let us prove that

$$\tilde{\mathcal{B}} = \bigvee_{n \ge 0} \tilde{\mathcal{K}}^n(\pi_{\tilde{X},X}^{-1}(\mathcal{B})).$$
(11)

The inclusion $\bigvee_{n \ge 0} \tilde{\mathcal{K}}^n(\pi_{\tilde{X},X}^{-1}(\mathcal{B})) \subset \tilde{\mathcal{B}}$ is clear.

Let k and n be two non-negative integers. We consider a set A of $\tilde{\mathcal{B}}$ of the form $A = [\omega_0, \omega_1, \dots, \omega_n]_0^n \times [e_{-k}, \dots, e_0]_{-k}^0 \times [a, b]$. We then have

$$A = (\Omega \times [e_{-k}, \dots, e_{-1}]_{-k}^{-1} \times [a, b]) \cap ([\omega_0, \omega_1, \dots, \omega_n]_0^n \times [e_0]_0 \times [a, b]).$$

Set $A_{-} = \Omega \times [e_{-k}, \dots, e_{-1}]_{-k}^{-1} \times [a, b]$ and $A_{+} = [\omega_0, \omega_1, \dots, \omega_n]_0^n \times [e_0]_0 \times [a, b]$.

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Then, $A_{-} \in \tilde{\mathcal{K}}^{k}(\pi_{\tilde{X},X}^{-1}(\mathcal{B}))$ and $A_{+} \in \pi_{\tilde{X},X}^{-1}(\mathcal{B})$, which implies that $A \in \bigvee_{n \ge 0} \tilde{\mathcal{K}}^{n}(\pi_{\tilde{X},X}^{-1}(\mathcal{B}))$, and proves equation (11).

We can now define the measure $\tilde{\mu}$ on \tilde{X} . For $A \in \bigvee_{n=0}^{N} \tilde{\mathcal{K}}^{n}(\pi_{\tilde{Y}}^{-1}(\mathcal{B}))$, we set

$$\tilde{\mu}(A) := \mu(\pi_{\tilde{X},X}(\tilde{\mathcal{K}}^{-N}(A)))$$

In some sense, the measure $\tilde{\mu}$ of this set *A* is obtained by 'shifting' the set *A* to the future, so that it can be viewed as a set in the first extension *X*. No information is lost by projecting it on *X*, allowing us to use the measure μ . The measure $\tilde{\mu}$ is well defined: if $A \in \bigvee_{n=0}^{N} \tilde{\mathcal{K}}^{n}(\pi_{\tilde{X}|X}^{-1}(\mathcal{B}))$, then for any integer $k \in \mathbb{N}$, we have

$$\mu(\pi_{\tilde{X},X}(\tilde{\mathcal{K}}^{-N}(A))) = \mu(\pi_{\tilde{X},X}(\tilde{\mathcal{K}}^{-N-k}(A)))$$

since μ is \mathcal{K} -invariant. For any $B \in \mathcal{B}$, we have $\tilde{\mu} \circ \pi_{\tilde{X},X}^{-1}(B) = \mu(B)$, and the measure $\tilde{\mu}$ is $\tilde{\mathcal{K}}$ -invariant.

Finally, by construction, $\tilde{\mathcal{K}}$ is one to one on \tilde{X} . It implies that $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ is a natural extension of (X, μ, \mathcal{K}) . Moreover, the natural extension of an ergodic system is ergodic, which provides the following result, as a direct consequence of Theorem 14.

COROLLARY 18. The system $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ is ergodic.

3.2. Natural extension of the initial system. The goal of this section is to prove that the previous extension is in fact a natural extension of the initial system ($\Omega \times I_{\beta}, m_p \otimes \mu_p, K_{\beta}$). To do so, we first introduce a canonic way to construct a natural extension of the initial system, then prove that the two extensions are isomorphic.

We adapt the construction described in [1, pp. 62], which provides a natural extension of K_{β} . We set

$$\underline{X} := \{(\underline{\omega}, \underline{x}) \in \{g, \ell\}^{\mathbb{Z}} \times I_{\beta}^{\mathbb{Z}} : \text{ for all } k \in \mathbb{Z}, T_{\omega_k} x_k = x_{k+1}\}.$$

Let $\underline{\mathcal{B}}$ be the σ -algebra generated by the cylinders of \underline{X} . We define the measure $\underline{\mu}$ on every set of the form $A = [\omega_k, \ldots, \omega_{n+k}]_k^{n+k} \times [I_k, \ldots, I_{n+k}]_k^{n+k}$, with $k \in \mathbb{Z}$, $n \in \mathbb{N}$, and for any $k \leq i \leq n+k$, $\omega_i \in \{g, \ell\}$, and I_i sub-interval of I_β by

$$\underline{\mu}(A) = m_p \otimes \mu_p \bigg([\omega_k, \ldots, \omega_{n+k}]_0^n \times \bigcap_{i=0}^n T_{\omega_{k+i}}^{-i}(I_{k+i}) \bigg).$$

We denote by \underline{K} the shift on \underline{X} .

PROPOSITION 19. [1] The dynamical system $(\underline{X}, \underline{\mu}, \underline{K})$ is a natural extension of the system $(\Omega \times I_{\beta}, m_p \otimes \mu_p, K_{\beta})$.

Let $(\omega, \mathbf{e}, x) \in \tilde{X}$. We denote by $\underline{\omega} = \underline{\omega}(\omega, \mathbf{e})$ the sequence of $\{g, \ell\}^{\mathbb{Z}}$ defined by $\underline{\omega}(k) = \omega_k$ if $k \ge 0$ and $\underline{\omega}(k) = \mathbf{e}_{k+1}(-1)$ if $k \le -1$. The sequence $\underline{\omega}$ is the bi-infinite sequence of applied transformations in the past, and transformations to be applied in the future. For any $k \in \mathbb{Z}$, we denote by x_k the real component of $\tilde{\mathcal{K}}^k(\omega, \mathbf{e}, x)$, and

 $\underline{x} := (x_k)_{k \in \mathbb{Z}}$. We define the application

$$\phi: \frac{\tilde{X} \to \underline{X}}{(\omega, \mathbf{e}, x) \to (\underline{\omega}, \underline{x})}$$

One easily checks the following lemma.

LEMMA 20. The application ϕ is a factor map from $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ to $(\underline{X}, \mu, \underline{K})$.

THEOREM 21. The factor map ϕ is an isomorphism. Therefore, the system $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ is a natural extension of $(\Omega \times I_{\beta}, m_p \otimes \mu_p, K_{\beta})$.

To prove this theorem, we use the notion of relatively independent joinings above a factor (see for example [8, pp. 126–127]).

Let \mathbb{P} be the relatively independent self-joining of the system $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ above its factor \underline{X} , that is, the measure on $\tilde{X} \times \tilde{X}$ such that for two measurable sets $A, B \in \tilde{X}$,

$$\mathbb{P}(A \times B) = \int_{\tilde{X}} \tilde{\mu}(A|\mathcal{F})\tilde{\mu}(B|\mathcal{F}) d\tilde{\mu},$$

where $\mathcal{F} = \phi^{-1}(\mathcal{B}(\underline{X}))$ is the factor σ -algebra associated to the factor $(\underline{X}, \underline{\mu}, \underline{K})$. We have the following properties:

- the measure \mathbb{P} is supported on the set $\{(a, b) \in \tilde{X} \times \tilde{X} : \phi(a) = \phi(b)\};$
- \mathbb{P} is supported on the diagonal of $\tilde{X} \times \tilde{X}$ if and only if the factor map ϕ is an isomorphism.

Let $a = (\omega, \mathbf{e}, x)$ and $b = (\omega', \mathbf{e}', x')$ be two elements of \tilde{X} such that $\phi(a) = \phi(b)$. We then have $\underline{\omega}(\omega, \mathbf{e}) = \underline{\omega}(\omega', \mathbf{e}')$ and $\underline{x} = \underline{x}'$. In other words, a and b describe two trajectories in the towers, starting from the same real number x and with the same transformations $\underline{\omega}$, but possibly following different floors \mathbf{e} and \mathbf{e}' . Proving that \mathbb{P} is supported by the diagonal of $\tilde{X} \times \tilde{X}$ consists in proving that for \mathbb{P} -almost every $(a, b) \in \tilde{X} \times \tilde{X}$, the two sequences of floors \mathbf{e} and \mathbf{e}' are actually equal.

The idea of the proof consists in taking two trajectories in \tilde{X} starting from the same real number x, with the same transformations, and proving that they end up in the same floor. Once these trajectories come at the same time in the same floor, they coincide in the future and in the past by injectivity.

We first work on the extension X. Let $k \in \mathbb{N}$ and consider the set of positive measure

$$C_k := \left\{ (\omega, e, x) \in X : x < \frac{1}{\beta^k}, \omega_0 = \omega_1 = \cdots = \omega_{k-1} = g \right\}.$$

The set C_k describes the points of X whose real component is close to 0, and on which the greedy transformation will be applied k times in a row. The label e does not take part in the definition of C_k , which means that the event C_k is measurable with respect to the factor \underline{X} : given two trajectories with the same real component and the same sequence of transformations, the event C_k happens almost surely (by ergodicity of $\tilde{\mathcal{K}}$ and because C_k is of positive measure) and at the same time for both trajectories. Knowing that C_k is realized actually gives some information about the label. We will prove that for large k and conditionally to C_k , both trajectories are in a floor of the greedy tower whose interval is large, with high probability.

Let A_k be the union of the floors of \mathcal{G} whose length is less than $1/\beta^k$. Let B_k be the complementary set of A_k in \mathcal{G} . Then the set C_k can be divided into three parts: $C_k \cap \mathcal{L}$, $C_k \cap A_k$, and $C_k \cap B_k$. The set $C_k \cap B_k$ is particularly interesting since a point in this set will end up in the greedy base E_g in at most k iterations: such a point will remain in the left part of its floor under the action of \mathcal{K} , and the length of the floor will eventually be larger than $1/\beta$, allowing the point to go to E_g in less than k iterations (see Figures 3 and 4). In the following, we prove that for large k and knowing C_k , the event $C_k \cap B_k$ happens with high probability.

We first introduce the following technical lemma.

LEMMA 22. The intervals of the lazy tower \mathcal{L} never contain 0.

Proof. This property holds by construction of the extension \mathcal{K} . Indeed, if an interval of a floor of \mathcal{L} contained 0, it would have to be the image by T_g of an interval whose lower bound is $1/\beta$. However, if the lower bound of an interval of \mathcal{L} equals $1/\beta$, we are in the case where the part of the floor satisfying $x \in [1/\beta, 2/\beta]$ is sent to the greedy base when applying T_g .

For any label $e \in \mathcal{E}$, we set $Y_e := \Omega \times \{e\}$, and we have $E_e = Y_e \times I_e$. We then set $Y := \bigsqcup_{e \in \mathcal{E}} Y_e$, which is represented vertically in Figure 2. Let us define the measure ν on Y by, for any $e \in \mathcal{E}$ with $e = (h, \omega_{-n}, \ldots, \omega_{-1})$,

$$\nu_{|Y_e} = \frac{1}{\beta^n} \Gamma_h m_p([\omega_{-n},\ldots,\omega_{-1}]_0^{n-1}) m_p,$$

where Γ_h is one of the two constants Γ_g and Γ_ℓ defined in Theorem 1. In this way, we have

$$\mu_{|E_e} = \nu_{|Y_e} \otimes \lambda_{|I_e}.$$

Finally, we denote $\pi_{X,Y}$ the projection from X onto Y.

Let $\varepsilon > 0$ be small enough, so that $1 - (1 + 2(\beta - 1))/(\beta - 1)\mu(\mathcal{G})\varepsilon > 0.5$. Let $N \in \mathbb{N}$ be such that the measure μ of the floors beyond the level N in \mathcal{L} is less than ε . By the previous lemma, there exists $k_0 \in \mathbb{N}$ such that no floor of \mathcal{L} up to the level N has an interval whose lower bound is less than $1/\beta^k$.

As previously said, we have

$$\mu(C_k) = \mu(C_k \cap \mathcal{L}) + \mu(C_k \cap A_k) + \mu(C_k \cap B_k).$$
(12)

We want to estimate each term of the right-hand side of equation (12). We will prove that for large k, $C_k \cap \mathcal{L}$ and $C_k \cap A_k$ have a small measure, and that $C_k \cap B_k$ represents more than half of the measure of C_k .

LEMMA 23. For k large enough, we have

$$\mu(C_k \cap \mathcal{L}) \leqslant 2p^k \frac{\beta - 1}{\beta^k} \varepsilon.$$

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Proof. Let $k \ge k_0$. A floor of \mathcal{L} can intersect C_k only if the lower bound of its interval is less than $1/\beta^k$. The integer N is fixed such that every floor below the level N does not contain $1/\beta^k$. Therefore, only the floors beyond the level N can intersect C_k . Additionally, for any $e \in \mathcal{E}$, we have, by definition of the measure μ on each floor,

$$\mu\left(\left\{(\omega, e, x) \in E_e : \omega_0 = \dots = \omega_{k-1} = g, x < \frac{1}{\beta^k}\right\}\right) = p^k \mu\left(E_e \cap \left\{x < \frac{1}{\beta^k}\right\}\right).$$

We then have

$$\mu(C_k \cap \mathcal{L}) = p^k \sum_{e \in \mathcal{E}, e(0) = \ell, |e| \ge N+1} \nu(Y_e) \lambda \left(I_e \cap \left[0; \frac{1}{\beta^k}\right] \right).$$

For a floor E_e of the lazy tower such that $I_e \cap [0; 1/\beta^k] \neq \emptyset$, we have

$$\lambda(I_e) > \frac{1}{\beta - 1} - \frac{1}{\beta^k}.$$

For *k* large enough, we then have

$$\lambda(I_e) > \frac{1}{2(\beta-1)}.$$

It implies that

$$\lambda\left(I_e \cap \left[0; \frac{1}{\beta^k}\right]\right) \leqslant \frac{1}{\beta^k} \leqslant 2\frac{\beta - 1}{\beta^k}\lambda(I_e)$$

and hence

$$\mu(C_k \cap \mathcal{L}) \leqslant 2p^k \frac{\beta - 1}{\beta^k} \varepsilon.$$

LEMMA 24. For k large enough, we have

$$\mu(C_k \cap A_k) \leqslant \frac{p^k}{\beta^k} \varepsilon$$

Proof. Let $k \in \mathbb{N}$.

$$\mu(C_k \cap A_k) = p^k \sum_{e \in \mathcal{E}, E_e \subset A_k} \nu(Y_e) \lambda \left(I_e \cap \left[0; \frac{1}{\beta^k} \right] \right).$$

If E_e is a floor of A_k , then $\lambda(I_e \cap [0; 1/\beta^k]) \leq 1/\beta^k$. Therefore,

$$\mu(C_k \cap A_k) \leqslant \frac{p^k}{\beta^k} \nu(\pi_{X,Y}(A_k)).$$

We have

$$\bigcap_{k \ge 1} A_k = \emptyset$$

and hence

$$\pi_{X,Y}\bigg(\bigcap_{k\geqslant 1}A_k\bigg)=\varnothing.$$

Since the intersection of the A_k is decreasing, we have

$$\bigcap_{k\geqslant 1}\pi_{X,Y}(A_k)=\varnothing.$$

It implies that

$$\lim_{k\to\infty}\nu(\pi_{X,Y}(A_k))=0.$$

We can choose k large enough, so that $\nu(\pi_{X,Y}(A_k)) < \varepsilon$. Therefore,

$$\mu(C_k \cap A_k) \leqslant \frac{p^k}{\beta^k} \varepsilon.$$

PROPOSITION 25. For k large enough, we have

$$\frac{\mu(C_k \cap B_k)}{\mu(C_k)} > 0.5$$

Proof. Let $k \in \mathbb{N}$ be large enough such that Lemmas 23 and 24 are both satisfied. We have

$$\frac{\mu(C_k \cap B_k)}{\mu(C_k)} = 1 - \frac{\mu(C_k \cap A_k) + \mu(C_k \cap \mathcal{L})}{\mu(C_k)}.$$

From Lemmas 23 and 24, we then have

$$\frac{\mu(C_k \cap B_k)}{\mu(C_k)} \ge 1 - \frac{(p^k/\beta^k)\varepsilon + 2p^k((\beta - 1)/\beta^k)\varepsilon}{\mu(C_k)}$$
$$\ge 1 - \frac{(p^k/\beta^k)\varepsilon + 2p^k((\beta - 1)/\beta^k)\varepsilon}{\mu(C_k \cap \mathcal{G})}.$$

However,

$$\mu(C_k \cap \mathcal{G}) = p^k \sum_{e \in \mathcal{E}, e(0) = g} \nu(Y_e) \lambda \left(I_e \cap \left[0; \frac{1}{\beta^k} \right] \right)$$

For any floor E_e of \mathcal{G} , we have the minoration $\lambda(I_e \cap [0, 1/\beta^k]) \ge \lambda(I_e)(\beta - 1)/\beta^k$ and hence

$$\frac{\mu(C_k \cap B_k)}{\mu(C_k)} \ge 1 - \frac{(p^k/\beta^k)\varepsilon + 2p^k((\beta - 1)/\beta^k)\varepsilon}{p^k((\beta - 1)/\beta^k)\mu(\mathcal{G})} = 1 - \frac{1 + 2(\beta - 1)}{(\beta - 1)\mu(\mathcal{G})}\varepsilon.$$

The real ε is such that $1 - (1 + 2(\beta - 1))/(\beta - 1)\mu(\mathcal{G})\varepsilon > 0.5$, which gives the result.

All these computations are also valid on the extension \tilde{X} . Indeed, $\pi_{\tilde{X},X}^{-1}(C_k) \in \pi_{\tilde{X},X}^{-1}(\mathcal{B})$, which implies $\tilde{\mu}(\pi_{\tilde{X},X}^{-1}(C_k)) = \mu(C_k)$. In the following, we will still denote C_k instead of $\pi_{\tilde{X},X}^{-1}(C_k)$.

Since $\tilde{\mathcal{K}}$ is ergodic and $\mu(C_k) > 0$, almost every trajectory in \tilde{X} encounters C_k with frequency $\mu(C_k)$. Therefore, for \mathbb{P} -almost every couple of trajectories in \tilde{X} , both trajectories simultaneously encounter C_k (since C_k is measurable with respect to the factor \underline{X}) with frequency $\mu(C_k)$ (since \mathbb{P} is a self-joining of μ). In other words, we have, for \mathbb{P} -almost every $(a, b) \in \tilde{X} \times \tilde{X}$, for any $n \in \mathbb{N} \cup \{0\}$:

• $\tilde{\mathcal{K}}^n(a) \in C_k$ if and only if $\tilde{\mathcal{K}}^n(b) \in C_k$;

• $(1/N) \sum_{i=1}^{N} \mathbb{1}_{C_k}(\tilde{\mathcal{K}}^i(a)) = (1/N) \sum_{i=1}^{N} \mathbb{1}_{C_k}(\tilde{\mathcal{K}}^i(b)) \to \mu(C_k).$

The ergodicity of $\tilde{\mathcal{K}}$ and Proposition 25 also imply that among the occurrences of C_k , the event $C_k \cap B_k$ occurs with a frequency greater that 0.5 for the orbits of \mathbb{P} -almost every $(a, b) \in \tilde{X} \times \tilde{X}$ under $\tilde{\mathcal{K}}$. Therefore, for \mathbb{P} -almost every $(a, b) \in \tilde{X} \times \tilde{X}$, there exists $n_0 \in \mathbb{N} \cup \{0\}$ such that $\tilde{\mathcal{K}}^{n_0}(a) \in C_k \cap B_k$ and $\tilde{\mathcal{K}}^{n_0}(b) \in C_k \cap B_k$. Once the two trajectories are in $C_k \cap B_k$ simultaneously, they will end up in the greedy base in at most kiterations (note that even if one of the two points goes in the base before the other, it will 'wait' for the second one to join it. Since its real component is small enough, the point is in the atom P_0). In conclusion, for \mathbb{P} -almost every $(a, b) \in \tilde{X} \times \tilde{X}$, if we note $a = (\omega, \mathbf{e}, x)$ and $b = (\omega, \mathbf{e}', x)$, there exists an integer $M = M(\underline{\omega}, \underline{x}, \mathbf{e}, \mathbf{e}')$ such that $e_M = e'_M = (g)$.

Let $\delta > 0$. There exists an integer M_{δ} such that

$$\mathbb{P}(M > M_{\delta}) < \delta.$$

Therefore, we have $\mathbb{P}(e_{M_{\delta}} = e'_{M_{\delta}}) \ge 1 - \delta$. Since \mathbb{P} is $\tilde{\mathcal{K}} \times \tilde{\mathcal{K}}$ -invariant, we deduce that for any integer $k \in \mathbb{Z}$ and any $\delta > 0$,

$$\mathbb{P}(e_k = e'_k) \ge 1 - \delta$$

and hence $\mathbf{e} = \mathbf{e}'$. We proved that \mathbb{P} is supported by the diagonal of $\tilde{X} \times \tilde{X}$, which implies that ϕ is an isomorphism. Therefore, $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ is a natural extension of the initial system, which concludes the proof of Theorem 21.

3.3. *Bernoullicity.* We know that the natural extension of the initial random system is ergodic. In fact, we have the following stronger result.

THEOREM 26. The natural extension of the random β -transformation is isomorphic to a Bernoulli shift.

To prove this theorem, we introduce a generating partition of the two towers. We then prove that this partition is weak-Bernoulli, by an argument on countable state space Markov chains of Ito, Murata, and Totoki [9], which implies the theorem. In the following, we denote $\tilde{\mu}_g := \tilde{\mu}(\cdot|\tilde{E}_g)$, and by $\tilde{\mathcal{K}}_g$ the induced transformation of $\tilde{\mathcal{K}}$ on the base \tilde{E}_g .

LEMMA 27. Let (X, μ, T) be an ergodic dynamical system, and A a measurable set of X with positive measure. Denote by $(\tilde{X}, \tilde{\mu}, \tilde{T})$ a natural extension of the system, and by π a factor map from \tilde{X} onto X. Then the system induced by \tilde{T} on $\pi^{-1}(A)$ is a natural extension of the system induced by T on A.

Proof. We start by introducing the necessary notation. The natural extension of the system (X, μ, T) can be described by a two-sided shift \tilde{T} on the set

$$\tilde{X} := \{ \mathbf{x} = (x_n)_{n \in \mathbb{Z}} \in X^{\mathbb{Z}}, \text{ for all } n \in \mathbb{Z}, T(x_n) = x_{n+1} \},\$$

with the factor map π given by $\pi((x_n)_{n \in \mathbb{Z}}) = x_0$, and the measure $\tilde{\mu}$ defined by

$$\tilde{\mu}([A_k, A_{k+1}, \dots, A_{n+k}]_k^{n+k}) = \mu\left(\bigcap_{i=0}^n T^{-i}(A_{k+i})\right),$$

where A_k, \ldots, A_{n+k} are measurable sets of X and $[A_k, A_{k+1}, \ldots, A_{n+k}]_k^{n+k} := \{\mathbf{x} \in \tilde{X} :$ for $k \leq i \leq n+k, x_i \in A_i\}$.

Denote by $\tilde{A} := \pi^{-1}(A) = \{(x_n) \in \tilde{X}, x_0 \in A\}$ and $(\tilde{A}, \tilde{\mu}_{\tilde{A}}, \tilde{T}_{\tilde{A}})$ the system induced by \tilde{T} on \tilde{A} . The transformation $\tilde{T}_{\tilde{A}}$ is defined by $\tilde{T}_{\tilde{A}}(\mathbf{x}) = \tilde{T}^{r_A(x_0)}(\mathbf{x})$ for every $\mathbf{x} \in \tilde{A}$, where $r_A(x_0)$ is the return time of x_0 in A, which is the same as the return time of \mathbf{x} in \tilde{A} .

We denote by $(A, \mu(\cdot|A), T_A)$ the system induced by T on A. The natural extension of this induced system can also be described by a two-sided shift \underline{T} on the set $\underline{A} = \{\mathbf{x} = (x_n)_{n \in \mathbb{Z}} \in A^{\mathbb{Z}}, \text{ for all } n \in \mathbb{Z}, T_A(x_n) = x_{n+1}\}$, with the measure $\underline{\mu}$ defined in the same way as $\tilde{\mu}$.

We want to prove that the two systems $(\tilde{A}, \tilde{\mu}_{\tilde{A}}, \tilde{T}_{\tilde{A}})$ and $(\underline{A}, \underline{\mu}, \underline{T})$ are isomorphic. We define the map φ as

$$\varphi: \frac{\tilde{A} \to \underline{A}}{\mathbf{x} \mapsto ((\tilde{T}^n_{\tilde{A}}(\mathbf{x}))_0)_{n \in \mathbb{Z}}}$$

The map φ consists in keeping only the terms of the sequence x that belong to A.

We clearly have that $\varphi \circ \tilde{T}_{\tilde{A}} = \underline{T} \circ \varphi$ and that $\underline{\mu}$ is the pushforward measure $\tilde{\mu}_{\tilde{A}} \circ \varphi^{-1}$: the map φ is a factor map from $(\tilde{A}, \tilde{\mu}_{\tilde{A}}, \tilde{T}_{\tilde{A}})$ to $(\underline{A}, \mu, \underline{T})$.

We now prove that φ is an isomorphism. Let $\mathbf{y} = (y_n)_{n \in \mathbb{Z}} \in \underline{A}$. Let $n \in \mathbb{Z}$. From y_n , we can construct the finite sequence of the iterates of y_n :

$$s(y_n) := (y_n, T(y_n), \dots, T^{r_A(y_n)-1}(y_n)),$$

stopping just before $y_{n+1} = T^{r_A(y_n)}(y_n)$. Then the sequence **x** obtained by concatenating the sequences $s(y_n)$ for all $n \in \mathbb{Z}$ is, by construction, a sequence of \tilde{A} such that $\varphi(\mathbf{x}) = \mathbf{y}$.

Let us prove that φ is one-to-one. Let $\mathbf{w}, \mathbf{x} \in \tilde{A}$ such that $\varphi(\mathbf{w}) = \varphi(\mathbf{x})$. Then $w_0 = x_0$, which implies that for any $n \in \mathbb{N}$, $T^n(w_0) = T^n(x_0)$, that is, $x_n = w_n$. Let $j_0 := \max\{j < 0, \tilde{T}^j(\mathbf{w}) \in \tilde{A}\}$ and $j_1 := \max\{j < 0, \tilde{T}^j(\mathbf{x}) \in \tilde{A}\}$. The existence of the integers j_0 and j_1 is assured by Poincaré's recurrence theorem applied to \tilde{T}^{-1} . Note that $\tilde{T}^{j_0}(\mathbf{w}) = \tilde{T}_{\tilde{A}}^{-1}(\mathbf{w})$ and $\tilde{T}^{j_1}(\mathbf{x}) = \tilde{T}_{\tilde{A}}^{-1}(\mathbf{x})$. Since $\varphi(\mathbf{w}) = \varphi(\mathbf{x})$, we get that $\tilde{T}^{j_0}(\mathbf{w}) = \tilde{T}^{j_1}(\mathbf{x})$, that is, $w_{j_0} = x_{j_1}$. The integer $-j_0$ is the return time of w_{j_0} in A, which implies that $j_0 = j_1$, and $w_n = x_n$ for $n \ge j_0$. By induction, we prove that $\mathbf{w} = \mathbf{x}$, and φ is an isomorphism, which concludes the proof.

COROLLARY 28. The system $(\tilde{E}_g, \tilde{\mu}_g, \tilde{\mathcal{K}}_g)$ is a natural extension of the system $(E_g, \mu_g, \mathcal{K}_g)$.

PROPOSITION 29. The partition $\tilde{\mathcal{P}} := \pi^{-1}(\mathcal{P})$ is a generator of the system

$$(\tilde{E}_g, \tilde{\mu}_g, \tilde{\mathcal{K}}_g).$$

Proof. μ_g -almost every point $a \in E_g$ is associated to a unique sequence $(C_n^a)_{n \ge 0}$ of $\mathscr{C}^{\mathbb{N} \cup \{0\}}$ such that for any $n \in \mathbb{N} \cup \{0\}$, $\mathcal{K}_g^n(a) \in E_{C_n^a}$: the sequence $(C_n^a)_{n \ge 0}$ can be

viewed as the \mathcal{P} -name of the point *a*. Denote by $\phi_{\mathcal{P}}$ the application defined almost everywhere on E_g by $\phi_{\mathcal{P}}(a) = (C_n^a)_{n \ge 0}$. Denote by ν the image measure of μ_g by $\phi_{\mathcal{P}}$ and by σ the left shift on $\mathscr{C}^{\mathbb{N} \cup \{0\}}$. The application $\phi_{\mathcal{P}}$ is a factor map from the system $(E_g, \mu_g, \mathcal{K}_g)$ to the one-sided Bernoulli shift $(\mathscr{C}^{\mathbb{N} \cup \{0\}}, \nu, \sigma)$. From Proposition 10, the partition \mathcal{P} is a generator of the system $(E_g, \mu_g, \mathcal{K}_g)$. Therefore, the application $\phi_{\mathcal{P}}$ is an isomorphism.

We can construct a natural extension of the system $(\mathscr{C}^{\mathbb{N}\cup\{0\}}, \nu, \sigma)$ the same way we did for <u>*K*</u>. Doing so, we construct the system $(\mathscr{C}^{\mathbb{Z}}, \tilde{\nu}, \sigma)$, which is a two-sided Bernoulli shift on $\mathscr{C}^{\mathbb{Z}}$. Since the natural extensions of a system are isomorphic to one another, we deduce from the previous lemma that the systems $(\mathscr{C}^{\mathbb{Z}}, \tilde{\nu}, \sigma)$ and $(\tilde{E}_g, \tilde{\mu}_G, \tilde{\mathcal{K}}_g)$ are isomorphic, and thus conclude the proof.

From this partition $\tilde{\mathcal{P}}$, we construct the set

$$\mathcal{P}_{\tilde{X}} := \{ \tilde{\mathcal{K}}^k(\tilde{E}_C), C \in \mathscr{C}, 0 \leq k \leq |C| - 1 \}.$$

This family is obtained by 'unfolding' the partition $\tilde{\mathcal{P}}$ on the two towers: each path *C* of \mathscr{C} is not only described by the first atom \tilde{E}_C , but by the sequence of successive images of this atom by $\tilde{\mathcal{K}}$.

LEMMA 30. The set $\mathcal{P}_{\tilde{X}}$ is a partition of \tilde{X} .

Proof. Two distinct sets of $\mathcal{P}_{\tilde{X}}$ are clearly disjoint.

However, let $(\omega, \mathbf{e}, x) \in \tilde{X}$. Almost surely, there exists $k \in \mathbb{N} \cup \{0\}$ such that $\mathbf{e}_{-k} = (g)$ (applying Poincaré recurrence theorem to $\tilde{\mathcal{K}}^{-1}$). Let $k_0 := \min\{k \in \mathbb{N} \cup \{0\}, \mathbf{e}_{-k} = (g)\}$. The integer $-k_0$ corresponds to the last visit of the greedy base in the past orbit of (ω, \mathbf{e}, x) . Almost surely, there exists a unique $C \in \mathscr{C}$ such that $\tilde{\mathcal{K}}^{-k_0}(\omega, \mathbf{e}, x) \in \tilde{E}_C$. Since k_0 is minimal, we have $0 \leq k_0 \leq |C| - 1$, and $(\omega, \mathbf{e}, x) \in \tilde{\mathcal{K}}^{k_0}(\tilde{E}_C)$.

PROPOSITION 31. The partition $\mathcal{P}_{\tilde{X}}$ is a generator of the system $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$.

Proof. The goal is to prove that almost every point in \tilde{X} is uniquely determined by its $\mathcal{P}_{\tilde{X}}$ -name.

Let $(\omega, \mathbf{e}, x) \in \tilde{X}$. Almost surely, we can define the integer k_0 as in the previous proof:

$$k_0 := \min\{k \in \mathbb{N} \cup \{0\}, \mathbf{e}_{-k} = (g)\}.$$

We note $a := \tilde{\mathcal{K}}^{-k_0}(e, x) \in \tilde{E}_g$. Almost surely, *a* is characterized by its $\tilde{\mathcal{P}}$ -name $(C_n)_{n \in \mathbb{Z}}$. To each path, C_n corresponds to a unique sequence of length $|C_n|$ of atoms of $\mathcal{P}_{\tilde{X}}$: the sequence $(\tilde{\mathcal{K}}^k(\tilde{E}_{C_n}), 0 \leq k \leq |C_n| - 1)$. Therefore, *a* is characterized by the unique sequence $(P_n)_{n \in \mathbb{Z}}$ of elements of $\mathcal{P}_{\tilde{X}}$ such that for any $n \in \mathbb{Z}$,

$$\tilde{\mathcal{K}}^n(a) \in P_n.$$

We deduce that the point (ω, \mathbf{e}, x) is uniquely determined by the sequence $(P_{n+k})_{n \in \mathbb{Z}}$.

We then define the isomorphism φ , which associates to almost every point in \tilde{X} its $\mathcal{P}_{\tilde{X}}$ -name. We denote by η the image measure of $\tilde{\mu}$ by φ and by σ the left shift on $\mathcal{P}_{\tilde{X}}^{\mathbb{Z}}$. Then φ is a factor map from $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ to $(\mathcal{P}_{\tilde{Y}}^{\mathbb{Z}}, \eta, \sigma)$.

The measure η is the Markov measure generated by the invariant distribution $(\tilde{\mu}(P))_{P \in \mathcal{P}_{\tilde{v}}}$ and the following transition probabilities.

- From an atom $\tilde{\mathcal{K}}^k(\tilde{E}_C)$ with $C \in \mathscr{C}$ and $0 \leq k \leq |C| 2$, the process goes to the atom $\tilde{\mathcal{K}}^{k+1}(\tilde{E}_C)$ with probability 1.
- *K*⁻¹ (*E_C*) with probability 1.
 From an atom *K*^{|C|-1}(*Ẽ_C*) with *C* ∈ *C*, the process can go to *Ẽ_{C'}* for any *C'* ∈ *C* with probability μ_g(*E_{C'}*).

Therefore, we can identify the natural extension $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ with a Markov shift on $\mathcal{P}_{\tilde{X}}^{\mathbb{Z}}$. This Markov shift is irreducible and recurrent, since $\tilde{\mathcal{K}}$ is ergodic. Moreover, it is aperiodic (the atom $[g]_0 \times \{(g)\} \times [0, 1/\beta[$ can appear twice in a row, for example). Proposition 2 on pp. 579 of [9] implies that the partition $\mathcal{P}_{\tilde{X}}$ is weak-Bernoulli, and the system $(\tilde{X}, \tilde{\mu}, \tilde{\mathcal{K}})$ is isomorphic to a Bernoulli shift. This concludes the proof of Theorem 26.

- 3.4. Open questions.
- (1) Can we generalize the extensions of this paper to any $\beta > 2$? We must take into account the different branches of the greedy and lazy transformations in the construction of the natural extension, especially the branches that induce a return to a base.
- (2) It seems natural to apply this kind of construction to other dynamical systems, either deterministic or random ones. It seems that it can be generalized to random systems with piecewise linear maps, even when the branches have different slopes. Indeed, with the construction of the towers, we can keep track of the exact branch that is being applied at each step.

However, systems based on continued fraction expansions should pose many more difficulties. In [10], the authors highlight a bijection between the β parameter of the non-integer base expansions and the λ parameter of λ -continued fractions. The first question would be to identify the analog of the greedy and lazy transformations in this case, and then to closely study this bijection.

(3) The description of the extension (X, μ, K) provides a better understanding of the random β-transformation. The extension shows off some nice renewal times at each passage to the greedy base, which provided a proof of the Bernoullicity of the natural extension. We can hope to get limit theorems as well thanks to this property. Indeed, in [18], Young describes a certain class of dynamical systems with renewal properties on reference sets. Young shows that the tail distribution of the return time to this reference set is directly linked to the convergence to equilibrium and limit theorems. In our setup, an estimation of the length of paths of C should help in this direction. Could we then obtain a central limit theorem for the digit sequences in base β following a fixed sequence ω? Or an average on ω?

(4) Our construction heavily relies on the independent choice of the transformations at each step. Could we adapt some of these constructions to more general stationary measures on Ω ? For example, Markov measures?

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