# Smoothing Equations for Large Pólya Urns 

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#### Abstract

Consider a balanced nontriangular two-color Pólya-Eggenberger urn process, assumed to be large, which means that the ratio $\sigma$ of the replacement matrix eigenvalues satisfies $1 / 2<\sigma<1$. The composition vector of both discrete-time and continuous-time models admits a drift which is carried by the principal direction of the replacement matrix. In the second principal direction, this random vector admits also an almost sure asymptotics and a real-valued limit random variable arises, named $W^{\text {DT }}$ in discrete time and $W^{\mathrm{CT}}$ in continuous time. The paper deals with the distributions of both $W$. Appearing as martingale limits, known to be nonnormal, these laws remain up to now rather mysterious. Exploiting the underlying tree structure of the urn process, we show that $W^{\mathrm{DT}}$ and $W^{\mathrm{CT}}$ are the unique solutions of two distributional systems in some suitable spaces of integrable probability measures. These systems are natural extensions of distributional equations that already appeared in famous algorithmical problems like Quicksort analysis. Existence and unicity of the solutions of the systems are obtained by means of contracting smoothing transforms. Via the equation systems, we find upper bounds for the moments of $W^{\mathrm{DT}}$ and $W^{\mathrm{CT}}$ and we show that the laws of $W^{\mathrm{DT}}$ and $W^{\mathrm{CT}}$ are moment determined. We also prove that $W^{\mathrm{DT}}$ is supported by the whole real line, its exponential moment generating series has an infinite radius of convergence and $W^{\mathrm{DT}}$ admits a continuous density ( $W^{\mathrm{CT}}$ was already known to have a density, infinitely differentiable on $\mathbb{R} \backslash\{0\}$ and not bounded at the origin).


Keywords Pólya urn model • Martingale • Multitype branching process • Smoothing transforms • Contraction method • Moment-determined probability distributions

[^0]
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## 1 Introduction

Pólya urns provide a rich model for many situations in algorithmics. Consider an urn that contains red and black balls. Start with a finite number of red and black balls as initial composition (possibly monochromatic). At each discrete time $n$, draw a ball at random, notice its color, put it back into the urn and add balls according to the following rule: if the drawn ball is red, add $a$ red balls and $b$ black balls; if the drawn ball is black, add $c$ red balls and $d$ black balls. The integers $a, b, c, d$ are assumed to be nonnegative. ${ }^{1}$ Thus, the replacement rule is described by the so-called replacement matrix

$$
R=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) .
$$

"Drawing a ball at random" means choosing uniformly among the balls contained in the urn. That is why this model is related to many situations in mathematics, algorithmics or theoretical physics where a uniform choice among objects determines the evolution of a process. See Johnson and Kotz's book [20], Mahmoud's book [26] or Flajolet et al. [18] for many examples.

In the present paper, the urn is assumed to be balanced, which means that the total number of balls added at each step is a constant

[^1]$$
S=a+b=c+d
$$

The composition vector of the urn at time $n$ is denoted by

$$
U^{\mathrm{DT}}(n)=\binom{\text { number of red balls at time } n}{\text { number of black balls at time } n} .
$$

Two main points of view are classically used on this random vector. The forward point of view consists in considering the composition vector sequence $\left(U^{\mathrm{DT}}(n)\right)_{n \in \mathbb{N}}$ as a $\mathbb{N}^{2}$-valued Markov chain. The information on the successive configurations is thus concentrated in a global object: the random process, giving access to probabilistic tools like martingales, embedding in continuous time, branching processes. A vast part of the literature on Pólya urns relies on such probability tools, dealing most often with natural extensions of the model to a random replacement matrix or to an arbitrary finite number of colors. The forward point of view is particularly efficient to get results on the asymptotics of the process. See for instance Janson's seminal paper [19] or [32] for an extensive state of the art on such methods.

Alternatively, a natural feature consists in using the recursive properties of the random structure through a divide-and-conquer principle. This is the backward point of view. Applied to generating functions, it is the base tool for analytic combinatorics methods, developed in Flajolet et al. papers [17,18]. Expressed in terms of the random process, the backward approach leads to dislocation equations on limit distributions that can already be found in a wide generality in Janson [19]; these equations are further developed in [13] for two-color urns and in [11,12] for the urn related to $m$-ary search trees as well.

In order to state our results and also the asymptotic theorems they are based on, we first give some notations that are made more complete in Sect. 2. The eigenvalues of the replacement matrix $R$ are $S$ and the integer

$$
m:=a-c=d-b
$$

and we denote by

$$
\sigma:=\frac{m}{S} \leq 1
$$

the ratio between these eigenvalues. The particular case $\sigma=1$ is the original Pólya urn (see Pólya [31]); this process has a specific well-known asymptotics with a random drift. In appendix, our Sect. 6 is devoted to gather results on this almost sure limit and on the asymptotic Dirichlet distribution as well. When $\sigma<1$, it is well known that the asymptotics of the process has two different behaviors, depending on the position of $\sigma$ with respect to the value $1 / 2$ (see Athreya and Karlin [4] for the original result, Janson [19] or [32] for the results below). Briefly said,
(i) when $\sigma<\frac{1}{2}$, the urn is called small and, except when $R$ is triangular, the composition vector is asymptotically Gaussian ${ }^{2}$ :

[^2]$$
\frac{U^{\mathrm{DT}}(n)-n v_{1}}{\sqrt{n}} \underset{n \rightarrow \infty}{\stackrel{\mathcal{D}}{\rightarrow}} \mathcal{G}\left(0, \Sigma^{2}\right)
$$
where $v_{1}$ is a suitable eigenvector of ${ }^{t} R$ relative to $S$ and $\mathcal{G}$ a centered Gaussian vector with covariance matrix $\Sigma^{2}$ that has a simple closed form;
(ii) when $\frac{1}{2}<\sigma<1$, the urn is called large and the composition vector has a quite different strong asymptotic form:
\[

$$
\begin{equation*}
U^{\mathrm{DT}}(n)=n v_{1}+n^{\sigma} W^{\mathrm{DT}} v_{2}+o\left(n^{\sigma}\right) \tag{1}
\end{equation*}
$$

\]

where $v_{1}, v_{2}$ are suitable (nonrandom) eigenvectors of ${ }^{t} R$ relative to the respective eigenvalues $S$ and $m, W^{\mathrm{DT}}$ is a real-valued random variable arising as the limit of a martingale, the little $o$ being almost sure and in any $L^{p}, p \geq 1$.

Classically, like for any Markov chain, one can embed the discrete-time process $\left(U^{\mathrm{DT}}(n)\right)_{n \in \mathbb{Z}_{\geq 0}}$ into continuous time. In the case of Pólya urns having a replacement matrix with nonnegative entries, this defines a two-type branching process

$$
\left(U^{\mathrm{CT}}(t)\right)_{t \in \mathbb{R}_{\geq 0}}
$$

A similar phase transition occurs when $t$ tends to infinity: for small urns, the process $U^{\mathrm{CT}}$ has a (random) almost sure drift and satisfies a gaussian central limit theorem (see Janson [19]). When the urn is large, the asymptotic behavior of the process, when $t$ tends to infinity, is given by

$$
U^{\mathrm{CT}}(t)=e^{S t} \xi v_{1}(1+o(1))+e^{m t} W^{\mathrm{CT}} v_{2}(1+o(1))
$$

where $\xi$ is Gamma-distributed, $W^{\mathrm{CT}}$ is a real-valued random variable arising as the limit of a martingale, the little $o$ is almost sure and in any $L^{p}, p \geq 1$, the basis $\left(v_{1}, v_{2}\right)$ of deterministic vectors being the same one as in (1). These asymptotic results are more detailed in Sect. 2. Because of the canonical link between $U^{\mathrm{DT}}$ and $U^{\mathrm{CT}}$ via stopping times, the two random variables $W^{\mathrm{DT}}$ and $W^{\mathrm{CT}}$ are related by the so-called martingale connexion as explained in Sect. 2.3. Consequently, any information about one distribution is of interest for the other one. All along the paper, the symbol $D T$ is used to qualify discrete-time objects while $C T$ will refer to continuous-time ones.

In this article, we are interested by large urns. More precisely, the attention is focused on the nonclassical distributions in $W^{\mathrm{DT}}$ and $W^{\mathrm{CT}}$ when the replacement matrix $R$ is not triangular (i.e. when $b c \neq 0$ ). For example, $W^{\mathrm{CT}}$ is not normally distributed, which can be seen on its exponential moment generating series that has a radius of convergence equal to zero, as shown in [13] (see Sect. 5 for more details). Because of the martingale connexion, this implies that $W^{\mathrm{DT}}$ is not normal either. Our main goal is to get descriptions of these laws (density, moments, tail, ...).

What is already known about $W^{\mathrm{DT}}$ or $W^{\mathrm{CT}}$ ? In [13], the Fourier transform of $W^{\mathrm{CT}}$ is "explicitely" calculated, in terms of the inverse of an abelian integral on the Fermat curve of degree $m$. The existence of a density with respect to the Lebesgue measure on $\mathbb{R}$ and the fact that $W^{\mathrm{CT}}$ is supported by the whole real line are deduced from this
closed form. Nevertheless, the order of magnitude of the moments and the question of the determination of the law by its moments remained open questions. The shape of the density was mysterious, too. The present paper answers to these questions in Sects. 5 and 3.5, respectively.

In the present text, we exploit the underlying tree structure of a Pólya urn. Governing both the backward and the forward points of view, it contains a richer structure than the plain composition vector process. Section 3 is devoted to highlighting this tree process and to derive decomposition properties on the laws of the composition vector at finite time. These decompositions directly lead to distributional fixed-point systems (15) and (18), respectively, satisfied by $W^{\mathrm{DT}}$ and $W^{\mathrm{CT}}$, as stated in Theorem 4 and Theorem 5.

With a slightly different approach, Knape and Neininger [22] start from the tree decomposition of the discrete Pólya urn and establish the fixed-point system (15) with the contraction method tools developed in Neininger and Rüschendorf [28]. This complementary point of view does not take advantage of the limit random variable $W^{\text {DT }}$ but applies for small and large urns together, allowing to find limit Gaussian distributions thus providing an alternative method to the embedding method used by Janson in [19].

Sometimes called fixed-point equations for the smoothing transform or just smoothing equations in the literature (Liu [25], Durrett and Liggett [15]), distributional equations of type

$$
\begin{equation*}
X \stackrel{\mathcal{L}}{=} \sum A_{i} X^{(i)} \tag{2}
\end{equation*}
$$

have given rise to considerable interest in, and literature on. For a survey, see Aldous and Bandyopadhyay [1]. In theoretical probability, they are of relevance in connexion with branching processes (like in Liu [23], Biggins and Kyprianou [8], Alsmeyer et al [2]) or with Mandelbrot cascades (Mandelbrot [27], Barral [6]). They occur in various areas of applied probability, and also on the occasion of famous problems arising in analysis of algorithms, like Quicksort (Rösler [33]). They are naturally linked with the analysis of recursive algorithms and data structures (Neininger and Rüschendorf [29], surveys in Rösler and Rüschendorf [34] or Neininger and Rüschendorf [30])

Most often, in Eq. (2), the $A_{i}$ are given random variables and the $X^{(i)}$ are independent copies of $X$, independent of the $A_{i}$ as well. Our System (18) with unknown real-valued random variables (or distributions) $X$ and $Y$ is the following:

$$
\left\{\begin{array}{l}
X \stackrel{\mathcal{L}}{=} U^{m}\left(\sum_{k=1}^{a+1} X^{(k)}+\sum_{k=a+2}^{S+1} Y^{(k)}\right) \\
Y \stackrel{\mathcal{L}}{=} U^{m}\left(\sum_{k=1}^{c} X^{(k)}+\sum_{k=c+1}^{S+1} Y^{(k)}\right)
\end{array}\right.
$$

where $U$ is uniform on [0,1], $X^{(k)}$ and $Y^{(k)}$ are respective copies of $X$ and $Y$, all being independent of each other and of $U$. Our System (15) for the discrete-time limit
$W^{\text {DT }}$, slightly more complicated, is essentially of the same type. These systems can be seen as natural generalizations of equations of type (2), as set out in Neininger and Rüschendorf [28]. Section 4 is devoted to the existence and the unicity of solutions of our systems by means of a contraction method (Theorems 7 and 8), leading to a characterization of $W^{\mathrm{DT}}$ and $W^{\text {CT }}$ distributions.

Finally, in Sect. 5, we take advantage of the fixed-point systems again to give accurate bounds on the moments of $W^{\mathrm{CT}}$ (Lemma 7). Using this lemma, we show that the law $W^{\mathrm{CT}}$ is determined by its moments and that the Laplace series of $W^{\mathrm{DT}}$ has an infinite radius of convergence (Corollary to Theorem 9).

## 2 Two-Color Pólya Urn: Definition and Asymptotics

### 2.1 Notations and Asymptotics in Discrete Time

Consider a two-color Pólya-Eggenberger urn random process. We adopt notations of the introduction: the replacement matrix $R=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$ is assumed to have nonnegative entries, the integers $S$ as balance and $m$ as second smallest eigenvalue. We assume $R$ to be nontriangular, i.e. that $b c \neq 0$; this implies that $m \leq S-1$. Moreover, the paper deals with large urns which means that the ratio $\sigma=m / S$ is assumed to satisfy

$$
\sigma>\frac{1}{2}
$$

We denote by $v_{1}$ and $v_{2}$ the vectors

$$
\begin{equation*}
v_{1}=\frac{S}{(b+c)}\binom{c}{b} \quad \text { and } \quad v_{2}=\frac{S}{(b+c)}\binom{1}{-1} \tag{3}
\end{equation*}
$$

they are eigenvectors of the matrix ${ }^{t} R$, respectively, associated with the eigenvalues $S$ and $m$. Let also $\left(u_{1}, u_{2}\right)$ be the dual basis

$$
\begin{equation*}
u_{1}(x, y)=\frac{1}{S}(x+y) \quad \text { and } \quad u_{2}(x, y)=\frac{1}{S}(b x-c y) ; \tag{4}
\end{equation*}
$$

$u_{1}$ and $u_{2}$ are eigenforms of ${ }^{t} R$, respectively, associated with the eigenvalues $S$ and $m$.
When the urn contains $\alpha$ white balls and $\beta$ black balls at (discrete) time 0 , the composition vector at time $n \in \mathbb{N}$ is denoted by

$$
U_{(\alpha, \beta)}^{\mathrm{DT}}(n) .
$$

Since the urn is assumed to be large, the asymptotics of its composition vector is given by the following result.

Theorem 1 (Asymptotics of discrete-time process, [19,32])
Let $\left(U_{(\alpha, \beta)}^{\mathrm{DT}}(n)\right)_{n \in \mathbb{N}}$ be a large Pólya urn discrete-time process. Then, when $n$ tends to infinity,

$$
\begin{equation*}
U_{(\alpha, \beta)}^{\mathrm{DT}}(n)=n v_{1}+n^{\sigma} W_{(\alpha, \beta)}^{\mathrm{DT}} v_{2}+o\left(n^{\sigma}\right) \tag{5}
\end{equation*}
$$

where $v_{1}$ and $v_{2}$ are the nonrandom vectors defined by (3), $W_{(\alpha, \beta)}^{\mathrm{DT}}$ is the real-valued random variable defined by

$$
\begin{equation*}
W_{(\alpha, \beta)}^{\mathrm{DT}}:=\lim _{n \rightarrow+\infty} \frac{1}{n^{\sigma}} u_{2}\left(U_{(\alpha, \beta)}^{\mathrm{DT}}(n)\right) \tag{6}
\end{equation*}
$$

$u_{2}$ being defined in (4), and where o() means almost surely and in any $L^{p}, p \geq 1$.

A proof of this result can be found in Janson [19] by means of embedding in continuous time, under an irreducibility assumption. Another proof, which is valid in any case and that remains in discrete time, is also given in [32]. The present paper is focused on the distribution of $W_{(\alpha, \beta)}^{\mathrm{DT}}$ which appears in both proofs as the limit of a bounded martingale. One remarkable fact that does not occur for small urns (i.e. when $\sigma \leq 1 / 2)$ is that the distribution of $W_{(\alpha, \beta)}^{\mathrm{DT}}$ actually depends on the initial composition vector $(\alpha, \beta)$. For example, its expectation turns out to be

$$
\begin{equation*}
\mathbb{E}\left(W_{(\alpha, \beta)}^{\mathrm{DT}}\right)=\frac{\Gamma\left(\frac{\alpha+\beta}{S}\right)}{\Gamma\left(\frac{\alpha+\beta}{S}+\sigma\right)} \frac{b \alpha-c \beta}{S} . \tag{7}
\end{equation*}
$$

This formula, explicitely stated in [13] can be shown by elementary means or using the convergent martingale

$$
\left(\frac{u_{2}\left(U_{(\alpha, \beta)}^{\mathrm{DT}}(n)\right)}{\prod_{0 \leq k \leq n-1}\left(1+\frac{\sigma}{k+\frac{\alpha+\beta}{S}}\right)}\right)_{n \in \mathbb{N}}
$$

For more developments about this discrete martingale which is the essential tool in the discrete method for proving Theorem 1, see [32].

The approach in analytic combinatorics makes easy to compute the probability generating function of the number of (say) red balls in the urn at finite time, by iteration of some suitable partial differential operator. The treatment of Pólya urns by analytic combinatorics, which is due to P. Flajolet and his co-authors, can be found in [17]. Figure 1 is the exact distribution of the (normalized) number of red balls after 300 drawings, centered around its expectation. The computations have been managed using Maple and concern the (large) urn with replacement matrix $R=\left(\begin{array}{cc}18 & 2 \\ 3 & 17\end{array}\right)$ and respective initial compositions $(1,0),(1,1)$ and $(0,1)$.

Some direct first observations can be made on these pictures. For example, one gets an illustration of the decomposition formula (12) which states that the distribution of $U_{(1,1)}$ is decomposed as a weighted convolution of $U_{(1,0)}$ 's and $U_{(0,1)}$ 's.


Fig. 1 Starting from initial composition $(\alpha, \beta)$, exact distribution of the number of red balls after $n=300$ drawings, centered around its mean and divided by $n^{\sigma}$ (Color figure online)

### 2.2 Embedding in Continuous Time

Classically, the discrete-time process is embedded in a continuous time multitype branching process; the idea of embedding discrete urn models in continuous-time branching processes goes back at least to Athreya and Karlin [4] and a description is given in Athreya and Ney [5], Section 9. The method has been revisited and developed by Janson [19], and we summarize hereunder the results obtained in [13].

We define the continuous-time Markov branching process

$$
\left(U_{(\alpha, \beta)}^{\mathrm{CT}}(t)\right)_{t \in \mathbb{R}_{\geq 0}}
$$

as being the embedded process of $\left(U_{(\alpha, \beta)}^{\mathrm{DT}}(n)\right)_{n \in \mathbb{N}}$. It starts from the same initial condition $U_{(\alpha, \beta)}^{\mathrm{CT}}(0)=U_{(\alpha, \beta)}^{\mathrm{DT}}(0)=(\alpha, \beta)$; at any moment, each ball is equipped with an $\mathcal{E} x p(1)$-distributed ${ }^{3}$ random clock, all the clocks being independent. When the clock of a red ball rings, $a$ red balls and $b$ black balls are added in the urn; when the ringing clock belongs to a black ball, one adds $c$ red balls and $d$ black balls, so that the replacement rules are the same as in the discrete-time urn process.

The important benefit of considering such a process comes from the independence of the subtrees in the branching process. In the continuous-time urn process, each ball reproduces independently from the other balls.

The asymptotics of this process is given by the following theorem.
Theorem 2 (Asymptotics of continuous-time process, [19,13])
Let $\left(U_{(\alpha, \beta)}^{\mathrm{CT}}(t)\right)_{t \geq 0}$ be a large Pólya urn continuous-time process. Then, when t tends to infinity,

$$
\begin{equation*}
U_{(\alpha, \beta)}^{\mathrm{CT}}(t)=e^{S t} \xi v_{1}(1+o(1))+e^{m t} W_{(\alpha, \beta)}^{\mathrm{CT}} v_{2}(1+o(1)), \tag{8}
\end{equation*}
$$

where $v_{1}, v_{2}, u_{1}, u_{2}$ are defined in (3) and (4), $\xi$ and $W_{(\alpha, \beta)}^{\mathrm{CT}}$ are real-valued random variables defined by

[^3]\[

$$
\begin{aligned}
\xi & =\lim _{t \rightarrow+\infty} e^{-S t} u_{1}\left(U^{\mathrm{CT}}(t)\right), \\
W_{(\alpha, \beta)}^{\mathrm{CT}} & =\lim _{t \rightarrow+\infty} e^{-m t} u_{2}\left(U_{(\alpha, \beta)}^{\mathrm{CT}}(t)\right),
\end{aligned}
$$
\]

all the convergences are almost sure and in any $\mathrm{L}^{p}$-space, $p \geq 1$. Furthermore, $\xi$ is Gamma $\left(\frac{\alpha+\beta}{S}\right)$ distributed.

Here again, the distribution of $W^{\mathrm{CT}}$ depends on the initial composition of the urn. For example, its expectation is

$$
\begin{equation*}
\mathbb{E}\left(W_{(\alpha, \beta)}^{\mathrm{CT}}\right)=\frac{b \alpha-c \beta}{S}, \tag{9}
\end{equation*}
$$

as can be seen from the continuous-time martingale

$$
\left(e^{-m t} u_{2}\left(U_{(\alpha, \beta)}^{\mathrm{CT}}(t)\right)\right)_{t \geq 0} .
$$

Some properties of $W^{\text {CT }}$ are already known. For example, it is supported by the whole real line $\mathbb{R}$ and admits a density. Moreover, this density is increasing on $\mathbb{R}_{<0}$, decreasing on $\mathbb{R}_{>0}$ and is not bounded in the neighborhood of the origin. Note that it is not an even function since $W^{\mathrm{CT}}$ is not centered. Finally, the characteristic function of $W^{\text {CT }}$ (i.e. its Fourier transform) is infinitely differentiable but not analytic at the origin: the domain of analyticity of $\mathbb{E} \exp \left(z W^{\mathrm{CT}}\right)$ is of the form $\mathbb{C} \backslash L_{+} \bigcup L_{-}$where $L_{+}$and $L_{-}$are half lines contained in $\mathbb{R}$, one of them being bordered at the origin. In particular, the exponential moment generating series of $W^{\mathrm{CT}}$ has a radius of convergence equal to zero, due to a ramification and a divergent series phenomenon as well. All these properties are shown in [13], based on the expression of this characteristic function in terms of the inverse of an abelian integral on the Fermat curve $x^{m}+y^{m}+z^{m}=0$.

### 2.3 Connexion Discrete Time/Continuous Time

As in any embedding into continuous time of a Markov chain, the discrete-time process and the continuous-time one are connected by

$$
\left(U^{\mathrm{CT}}\left(\tau_{n}\right)\right)_{n \in \mathbb{N}}=\left(U^{\mathrm{DT}}(n)\right)_{n \in \mathbb{N}}
$$

where

$$
0=\tau_{0}<\tau_{1}<\cdots<\tau_{n}<\cdots
$$

are the jumping times of the continuous process. These random times are independent of the positions $U^{\mathrm{CT}}\left(\tau_{n}\right)$. The embedding for urn processes is widely studied in Janson [19]. It is detailed in [13] in the special case of two-color Pólya urns. A dual formulation of this connexion is

$$
\left(U^{\mathrm{CT}}(t)\right)_{t \in \mathbb{R}_{\geq 0}}=\left(U^{\mathrm{DT}}(n(t))\right)_{t \in \mathbb{R}_{\geq 0}}
$$

where

$$
n(t):=\inf \left\{n \geq 0, \tau_{n} \geq t\right\}
$$

is the number of drawings in the urn before time $t$. After projection and normalization, these equalities provide the connexion between the limit variables $W_{(\alpha, \beta)}^{\mathrm{DT}}$ and $W_{(\alpha, \beta)}^{\mathrm{CT}}$ :

$$
\begin{equation*}
W_{(\alpha, \beta)}^{\mathrm{CT}} \stackrel{\mathcal{L}}{=} \xi^{\sigma} \cdot W_{(\alpha, \beta)}^{\mathrm{DT}} \tag{10}
\end{equation*}
$$

where $\xi$ and $W_{(\alpha, \beta)}^{\mathrm{DT}}$ are independent, $\xi$ being $\operatorname{Gamma}\left(\frac{\alpha+\beta}{S}\right)$ distributed. Note that the almost sure equality

$$
W_{(\alpha, \beta)}^{\mathrm{DT}}=\xi^{-\sigma} \cdot W_{(\alpha, \beta)}^{\mathrm{CT}}
$$

holds as well, the variables $\xi$ and $W_{(\alpha, \beta)}^{\mathrm{CT}}$ being, however, not independent.

## 3 Decomposition Properties

This section emphasizes the underlying tree structure of the urn process. This obvious vision is indeed the key in the following two decompositions: first, we reduce the study of $W_{(\alpha, \beta)}$ to the study of $W_{(1,0)}$ and $W_{(0,1)}$, called later on $X$ and $Y$ respectively, to lighten the notations and second, in Sect. 3.2, we exploit a "divide-and-conquer" property to deduce a system of fixed-point equations on $X$ and $Y$. The reasoning is detailed in discrete time. It is much more straightforward in continuous time, since the decomposition is contained inside the branching property. Detailed in [13], the continuous case is briefly recalled in Sect. 3.3.

The natural question "Is it possible to deduce the DT-system from the CT-system and conversely?" is partially addressed in Sect. 3.4.

### 3.1 Tree Structure in Discrete Time

In this section dealing with the discrete-time process, we skip the index $D T$ when no confusion is possible.

Let us make precise the tree structure of the urn process: a forest $\left(\mathcal{T}_{n}\right)$ grows at each drawing from the urn. At time 0 , the forest is reduced to $\alpha$ red nodes and $\beta$ black nodes, which are the roots of the forest trees. At time $n$, each leaf in the forest represents a ball in the urn. When a leaf is chosen (a ball is drawn), it becomes an internal node and gives birth to $(a+1)$ red leaves and $b$ black leaves, or $c$ red leaves and $(d+1)$ black leaves, according to the color of the chosen leaf.

The dynamics of the urn process was described saying "at each time $n$, a ball is uniformly chosen in the urn." It becomes "a leaf is uniformly chosen among the leaves of the forest." This forest therefore appears as a nonbinary colored generalization of a binary search tree.

For example, take the following urn with $R=\left(\begin{array}{ll}6 & 1 \\ 2 & 5\end{array}\right)$ as replacement matrix (it is a large urn) and start from $\alpha=3$ red balls and $\beta=2$ black balls. Below is a possible configuration after 3 drawings.


Initial red balls are numbered from 1 to $\alpha$ and initial black balls from $(\alpha+1)$ to $(\alpha+\beta)$. The following figure represents the forest coming from these initial balls.


For any $n \geq 0$ and $k \in\{1, \ldots, \alpha+\beta\}$, denote by $D_{k}(n)$ the number of leaves of the $k$-th tree in the forest at time $n$. Thus, at time $n$, the number of drawings in the $k$-th tree is $\frac{D_{k}(n)-1}{S}$. These numbers represent the time inside this $k$-th tree.

Remember that the balls of the whole urn are uniformly drawn at any time and notice that at each drawing in the $k$-th tree, $D_{k}(n)$ increases by $S$ : the random vector $D(n)=\left(D_{1}(n), \ldots, D_{\alpha+\beta}(n)\right)$ has exactly the same distribution as the composition vector at time $n$ of an $(\alpha+\beta)$-color Pólya urn process having $S I_{\alpha+\beta}$ as replacement matrix and $(1, \ldots, 1)$ as initial composition vector.

Gathering these arguments, the distribution of $U_{(\alpha, \beta)}(n)$ can be described the following way: consider simultaneously
(i) an original $(\alpha+\beta)$-color urn process $D=\left(D_{1}, \ldots, D_{\alpha+\beta}\right)$ having $S I_{\alpha+\beta}$ as matrix replacement and $(1, \ldots, 1)$ as initial condition;
(ii) for any $k \in\{1, \ldots, \alpha\}$, an urn process $U_{(1,0)}^{(k)}$ having $R$ as replacement matrix and $(1,0)$ as initial condition;
(iii) for any $k \in\{\alpha+1, \ldots, \alpha+\beta\}$, an urn process $U_{(0,1)}^{(k)}$ having $R$ as replacement matrix and $(0,1)$ as initial condition,
all these processes being independent of each other. Then, the process $U_{(\alpha, \beta)}=$ $\left(U_{(\alpha, \beta)}(n)\right)_{n}$ has the same distribution as the process defined by the sum of the $U_{(1,0)}^{(k)}$ and of the $U_{(0,1)}^{(k)}$ at respective times $\frac{D_{k}(n)-1}{S}$. In other words, for any $n \geq 0$,

$$
\begin{equation*}
U_{(\alpha, \beta)}(n) \stackrel{\mathcal{L}}{=} \sum_{k=1}^{\alpha} U_{(1,0)}^{(k)}\left(\frac{D_{k}(n)-1}{S}\right)+\sum_{k=\alpha+1}^{\alpha+\beta} U_{(0,1)}^{(k)}\left(\frac{D_{k}(n)-1}{S}\right) \tag{11}
\end{equation*}
$$

where the $U_{(1,0)}^{(k)}$ and the $U_{(0,1)}^{(k)}$ are respective copies of the random vector processes $U_{(1,0)}$ and $U_{(0,1)}$, all being independent of each other and of $D$.

The following claim is a direct consequence of Proposition 3 in Sect. 6.
Claim When $n$ goes off to infinity, $\frac{1}{n S}\left(D_{1}(n), \ldots, D_{\alpha+\beta}(n)\right)$ converges almost surely to a Dirichlet $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$ - distributed random vector, denoted by $Z=\left(Z_{1}, \ldots, Z_{\alpha+\beta}\right)$.

Notice that for any $k, D_{k}(n)$ tends almost surely to $+\infty$ when $n$ tends to infinity. Starting from Eq. (11), dividing by $n^{\sigma}$, taking the image by the second projection $u_{2}$ (notations of Sect. 2.1) and passing to the (almost sure) limit $n \rightarrow \infty$ thanks to Theorem 1, one obtains the following theorem.

Theorem 3 For any $(\alpha, \beta) \in \mathbb{N}^{2} \backslash(0,0)$, let $W_{(\alpha, \beta)}$ be the limit distribution of a large two-color discrete time Pólya urn process with ratio $\sigma$ and initial condition $(\alpha, \beta)$. Then,

$$
\begin{equation*}
W_{(\alpha, \beta)} \stackrel{\mathcal{L}}{=} \sum_{k=1}^{\alpha} Z_{k}^{\sigma} W_{(1,0)}^{(k)}+\sum_{k=\alpha+1}^{\alpha+\beta} Z_{k}^{\sigma} W_{(0,1)}^{(k)} \tag{12}
\end{equation*}
$$

where
(i) $Z=\left(Z_{1}, \ldots, Z_{\alpha+\beta}\right)$ is a Dirichlet-distributed random vector, with parameters $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$;
(ii) the $W_{(1,0)}^{(k)}$ and the $W_{(0,1)}^{(k)}$ are respective copies of $W_{(1,0)}$ and $W_{(0,1)}$, all being independent of each other and of $Z$.

Notice that any $Z_{k}$ is $\operatorname{Beta}\left(\frac{1}{S}, \frac{\alpha+\beta-1}{S}\right)$-distributed (see Sect. 6).

### 3.2 Discrete-Time Fixed-Point Equation

Theorem 3 shows that the limit distribution of a large urn process starting with any initial composition can be written as a function of two "elementary" particular laws, namely the laws of $W_{(1,0)}^{\mathrm{DT}}$ and $W_{(0,1)}^{\mathrm{DT}}$. The present section gives a characterization of these two distributions by means of a fixed-point equation.

Let $(U(n))_{n \geq 0}$ be a two-color Pólya urn process, with all the notations of Sect. 2.1. In order to simplify the notations, denote

$$
\left\{\begin{array}{l}
X:=W_{(1,0)}^{\mathrm{DT}}=\lim _{n \rightarrow+\infty} u_{2}\left(\frac{U_{(1,0)}(n)}{n^{\sigma}}\right)  \tag{13}\\
Y:=W_{(0,1)}^{\mathrm{DT}}=\lim _{n \rightarrow+\infty} u_{2}\left(\frac{U_{(0,1)}(n)}{n^{\sigma}}\right)
\end{array}\right.
$$

Focus now on the study of $U_{(1,0)}(n)$. At time 1 , the composition of the urn is deterministic: there are $(a+1)$ red balls and $b$ black balls. Exactly like in Sect. 3.1, the tree structure of the urn appears, with a forest starting from $(a+1)$ red balls and $b$ black balls. In the same example with replacement matrix $R=\left(\begin{array}{ll}6 & 1 \\ 2 & 5\end{array}\right)$, this fact is illustrated by the following figure:


For any $n \geq 1$, denote by $J_{k}(n)$ the number of leaves at time $n$ of the $k$-th subtree. Then, at time $n$, the number of drawings in the $k$-th subtree is $\frac{J_{k}(n)-1}{S}$ so that, as in Sect. 3.1, one gets the equation in distribution (a scalar version of this equation can be found in Knape and Neininger [22], Eq. (2))

$$
\begin{equation*}
U_{(1,0)}(n) \stackrel{\mathcal{L}}{=} \sum_{k=1}^{a+1} U_{(1,0)}^{(k)}\left(\frac{J_{k}(n)-1}{S}\right)+\sum_{k=a+2}^{S+1} U_{(0,1)}^{(k)}\left(\frac{J_{k}(n)-1}{S}\right) \tag{14}
\end{equation*}
$$

where the $U_{(1,0)}^{(k)}$ and the $U_{(0,1)}^{(k)}$ are respective copies of the random vector processes $U_{(1,0)}$ and $U_{(0,1)}$, all being independent of each other and of the $J_{k}$ 's. Besides, the random vector $\left(J_{1}(n), \ldots, J_{S+1}(n)\right)$ is exactly distributed like the composition vector at time $(n-1)$ of an $(S+1)$-color Pólya urn process having $S I_{S+1}$ as replacement matrix and $(1, \ldots, 1)$ as initial composition vector, so that, by Proposition 3 in Sect. 6,

$$
\frac{1}{n S}\left(J_{1}(n), \ldots, J_{S+1}(n)\right) \underset{n \rightarrow \infty}{\longrightarrow} V=\left(V_{1}, \ldots, V_{S+1}\right)
$$

almost surely, the random vector $V$ being $\operatorname{Dirichlet}\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$-distributed. Like in Sect. 3.1, divide Eq. (14) by $n^{\sigma}$, take the image by the second projection $u_{2}$ and pass to the limit $n \rightarrow \infty$ using Theorem 1. This leads to the following theorem.

Theorem 4 As defined just above by (13), let $X$ and $Y$ be the elementary limit laws of a large two-color discrete time Pólya urn process with replacement matrix $\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$, balance $S=a+b=c+d$ and ratio $\sigma>\frac{1}{2}$. Then, $X$ and $Y$ satisfy the distributional equations system

$$
\left\{\begin{array}{l}
X \xlongequal[=]{\mathcal{L}} \sum_{k=1}^{a+1} V_{k}^{\sigma} X^{(k)}+\sum_{k=a+2}^{S+1} V_{k}^{\sigma} Y^{(k)}  \tag{15}\\
Y \xlongequal{\mathcal{L}} \sum_{k=1}^{c} V_{k}^{\sigma} X^{(k)}+\sum_{k=c+1}^{S+1} V_{k}^{\sigma} Y^{(k)}
\end{array}\right.
$$

where
(i) $V=\left(V_{1}, \ldots, V_{S+1}\right)$ is a Dirichlet-distributed random vector, with parameters $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$;
(ii) the $X^{(k)}$ and the $Y^{(k)}$ are respective copies of $X$ and $Y$, all being independent of each other and of $V$.

Notice that any $V_{k}$ is distributed like a random variable $U^{S}, U$ being uniformly distributed on $[0,1]$. Equivalently, $V_{k}^{\sigma}$ is distributed like $U^{m}$ (notations of Sect. 2.1).

Remark 1 The system of equations (15) is the same as (41)-(42) in Knape and Neininger [22], which is not surprising since the rescaled number of red balls, considered in [22], actually tends to $\frac{S}{b+c} W_{(1,0)}^{\mathrm{DT}}$ via Theorem 1 . Thus, starting from the composition vector at finite time, there are two paths leading to (41)-(42) in [22] or to (15) above: either stay in the canonical basis of $\mathbb{R}^{2}$ or focus on the eigendirection $v_{2}$. In the first case, a contraction method is used in [22] to prove the convergence of the rescaled first coordinate of the composition vector, whereas in the second case, the convergence is given by a martingale limit and the contraction method is only used (cf. Sect. 4.2) to prove the unicity of the solution of (15).

### 3.3 Decomposition Properties in Continuous Time

Remember that $\left(U^{\mathrm{CT}}(t)\right)_{t}$ is a continuous-time branching process. Thanks to the branching property, the decomposition properties of this process are somehow automatic. First,

$$
U_{(\alpha, \beta)}^{\mathrm{CT}}(t)=[\alpha] U_{(1,0)}^{\mathrm{CT}}(t)+[\beta] U_{(0,1)}^{\mathrm{CT}}(t),
$$

where the notation $[n] X$ means the sum of $n$ independent random variables having the same distribution as $X$. Consequently, passing to the limit when $t \rightarrow+\infty$ after normalization and projection yields

$$
\begin{equation*}
W_{(\alpha, \beta)}^{\mathrm{CT}}=[\alpha] W_{(1,0)}^{\mathrm{CT}}+[\beta] W_{(0,1)}^{\mathrm{CT}} . \tag{16}
\end{equation*}
$$

This convolution formula expresses how the limit law $W^{\text {CT }}$ is decomposed in terms of elementary limit laws $W_{(1,0)}^{\mathrm{CT}}$ and $W_{(0,1)}^{\mathrm{CT}}$. It corresponds to the discrete-time decomposition shown in Theorem 3.

Now start from one red ball or from one black ball, and apply again the branching property at the first splitting time. As before, define $X^{\mathrm{CT}}$ and $Y^{\mathrm{CT}}$ by

$$
\left\{\begin{array}{l}
X^{\mathrm{CT}}:=W_{(1,0)}^{\mathrm{CT}}=\lim _{t \rightarrow+\infty} e^{-m t} u_{2}\left(U_{(1,0)}^{\mathrm{CT}}(t)\right),  \tag{17}\\
Y^{\mathrm{CT}}:=W_{(0,1)}^{\mathrm{CT}}=\lim _{t \rightarrow+\infty} e^{-m t} u_{2}\left(U_{(0,1)}^{\mathrm{CT}}(t)\right) .
\end{array}\right.
$$

Then, with the above Theorem 2, one gets the following result.
Theorem 5 ([13,19]) Let $X=X^{\mathrm{CT}}$ and $Y=Y^{\mathrm{CT}}$ be the elementary limit laws of a large two-color continuous time Pólya urn process with replacement matrix $\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$, balance $S=a+b=c+d$ and ratio $\sigma>\frac{1}{2}$, as defined just above by (17). Then, $X$ and $Y$ satisfy the distributional equations system

$$
\left\{\begin{array}{l}
X \stackrel{\mathcal{L}}{=} U^{m}\left(\sum_{k=1}^{a+1} X^{(k)}+\sum_{k=a+2}^{S+1} Y^{(k)}\right)  \tag{18}\\
Y \stackrel{\mathcal{L}}{=} U^{m}\left(\sum_{k=1}^{c} X^{(k)}+\sum_{k=c+1}^{S+1} Y^{(k)}\right),
\end{array}\right.
$$

where $U$ is uniform on $[0,1]$, where $X, X^{(k)}$ and $Y, Y^{(k)}$ are respective copies of $X^{\mathrm{CT}}$ and $Y^{\mathrm{CT}}$, all being independent of each other and of $U$.
Remark 2 As mentioned above, it is shown in [13] that $X^{\mathrm{CT}}$ (and $Y^{\mathrm{CT}}$ ) admit densities. The proof is based on the computation of the Fourier transform of $X^{\mathrm{CT}}$ in terms of the inverse of an abelian integral on a Fermat curve. This method is specific to twocolor urn processes. Theorems 4 and 5 give a new way of proving this fact by means of techniques that can be adapted from Liu's method (see [24] for example). This alternative method provides a perspective (adressed in a forthcoming paper): it can be applied to show that the limit laws of $d$-color large urns admit densities as well.

### 3.4 Connexion Between Continuous-Time and Discrete-Time Systems

In Sect. 2.3, we described the connexion between the limit laws of large urns in discrete and continuous time, called the martingale connexion. It was seen as a consequence of the embedding into continuous time of the initial discrete-time Markov chain defining the urn process. In this paragraph, we show how one can deduce the solutions of the continuous-time system (18) from the solutions of the discrete-time system (15).

Proposition 1 Let $X$ and $Y$ be solutions of the distributional System (15) and let $\xi$ be a Gamma-distributed random variable with parameter $\frac{1}{S}$, independent of $X$ and $Y$. Then, $\xi^{\sigma} X$ and $\xi^{\sigma} Y$ are solutions of the distributional System (18).

The assertion of Proposition 1 is a consequence of the following lemma which is an elementary result in probability theory.

Lemma 1 Consider the two following distributional equations with unknown independent real-valued random variables $X, X_{1}, \ldots, X_{S+1}$.

1- Equation D:

$$
X \xlongequal{\mathcal{L}} \sum_{1 \leq k \leq S+1} V_{k}^{\sigma} X_{k}
$$

where $\left(V_{1}, \ldots, V_{S+1}\right)$ is a Dirichlet-distributed random vector with parameter $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$, independent of $\left(X_{1}, \ldots, X_{S+1}\right)$.

2- Equation C:

$$
X \stackrel{\mathcal{L}}{=} V^{\sigma} \sum_{1 \leq k \leq S+1} X_{k}
$$

where $V$ is a Beta-distributed random variable with parameter $\left(\frac{1}{S}, 1\right)$ (in other words, $V^{1 / S}$ is uniformly distributed on $[0,1]$ ), independent of $\left(X_{1}, \ldots, X_{S+1}\right)$.

Let $V, \xi_{1}, \ldots, \xi_{S+1}$ be independent random variables, the $\xi_{k}$ 's being $\operatorname{Gamma}\left(\frac{1}{S}\right)$ distributed and $V$ being Beta $\left(\frac{1}{S}, 1\right)$-distributed. Denote

$$
\xi:=V \sum_{1 \leq j \leq S+1} \xi_{j}
$$

and, for any $k \in\{1, \ldots, S+1\}$,

$$
V_{k}:=\frac{\xi_{k}}{\sum_{1 \leq j \leq S+1} \xi_{j}}
$$

Then,
(i) the random variable $\xi$ is Gamma $\left(\frac{1}{S}\right)$-distributed;
(ii) the random vector $\left(V_{1}, \ldots, V_{S+1}\right)$ is independent of $\xi$ and Dirichlet-distributed with parameter $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$;
(iii) if $X, X_{1}, \ldots, X_{S+1}$ satisfy Equation $D$, if $\left(X_{1}, \ldots, X_{S+1}\right)$ is independent of $\left(V, \xi_{1}, \ldots, \xi_{S+1}\right)$ and if $X$ is independent of $\xi$, then $\xi^{\sigma} X, \xi_{1}^{\sigma} X_{1}, \ldots, \xi_{S+1}^{\sigma} X_{S+1}$ satisfy Equation C.

Proof (i) This can be seen for example by computation of moments (Beta and Gamma distributions are moment determined): the $p$-th moment of a Gamma $(\alpha)$ distribution is $\frac{\Gamma(\alpha+p)}{\Gamma(\alpha)}$ and the $p$-th moment of a $\operatorname{Beta}(\alpha, \beta)$ distribution is $\frac{\Gamma(\alpha+p) \Gamma(\alpha+\beta)}{\Gamma(\alpha+\beta+p) \Gamma(\alpha)}$ where $\Gamma$ is Euler Gamma function. Moreover, the sum of independent $\operatorname{Gamma}\left(\alpha_{1}\right), \ldots, \operatorname{Gamma}\left(\alpha_{d}\right)$-distributed random variables is $\operatorname{Gamma}\left(\alpha_{1}+\right.$ $\cdots+\alpha_{d}$ )-distributed. Assertion (i) is a direct consequence of these facts.
(ii) Classically, $\left(V_{1}, \ldots, V_{S+1}\right)$ is Dirichlet distributed and independent of the sum $\sum_{1 \leq j \leq S+1} \xi_{j}$. For a proof of this result, see for example Chaumont and Yor [10].

Since $V$ is independent of the $\xi_{k}$ 's, the random variable $\xi$ is also independent of $\left(V_{1}, \ldots, V_{S+1}\right)$.
(iii) Suppose that $X, X_{1}, \ldots, X_{S+1}$ satisfy Equation D. Then, $X \xlongequal{\mathcal{L}} \sum_{1 \leq k \leq S+1}$ $V_{k}^{\sigma} X_{k}$. Multiplying the equality by the random variable $\xi^{\sigma}$ leads to the distributive relation

$$
\xi^{\sigma} X \stackrel{\mathcal{L}}{=}\left(\frac{\xi}{\sum_{1 \leq j \leq S+1} \xi_{j}}\right)^{\sigma} \sum_{1 \leq k \leq S+1} \xi_{k}^{\sigma} X_{k}=V^{\sigma} \sum_{1 \leq k \leq S+1} \xi_{k}^{\sigma} X_{k}
$$

which makes the proof complete.

### 3.5 Densities

As shown in [13], the law of $W_{(\alpha, \beta)}^{\mathrm{CT}}$ turns out to be absolutely continuous with regard to Lebesgue measure on $\mathbb{R}$. In this section, the same property is deduced for $W_{(1,0)}^{\mathrm{DT}}$ and $W_{(0,1)}^{\mathrm{DT}}$ from the fixed-point equation (15). The decomposition property (12) then implies that any $W_{(\alpha, \beta)}^{\mathrm{DT}}$ also admits a density. The observations made in Sect. 2.1 on Fig. 1 can be seen as a first approximation of the shape of the density of $W_{(\alpha, \beta)}^{\mathrm{DT}}$.

The method we use is widely inspired from Q. Liu papers [24] and [25]. See also $[11,12]$ for an argumentation of the same vein for complex-valued probability measures. Applied to fixed-point equation (18), this method provides a second proof for the absolute continuity of $W_{(\alpha, \beta)}^{\mathrm{CT}}$; details are left to the reader.

For the whole Sect. 3.5, we denote

$$
X:=W_{(1,0)}^{\mathrm{DT}} \text { and } Y:=W_{(0,1)}^{\mathrm{DT}} .
$$

Let $\varphi_{X}$ and $\varphi_{Y}$ be the Fourier transforms of $X$ and $Y$ : for any $t \in \mathbb{R}$,

$$
\varphi_{X}(t)=\mathbb{E}\left(e^{i t X}\right) \text { and } \varphi_{Y}(t)=\mathbb{E}\left(e^{i t Y}\right) .
$$

Theorem 6 As defined just above, let $X$ and $Y$ be the elementary limit distributions of a large two-color discrete time Pólya urn process with replacement matrix $R=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$. As in the whole paper, let $S=a+b=c+d$ and $m=a-c=d-b \in] \frac{S}{2}, S[$ both eigenvalues of $R$. Then,
(i) the support of $X$ and $Y$ is the whole real line $\mathbb{R}$;
(ii) for any $\rho \in] 0, \frac{a+1}{m}[$, there exists $C>0$ such that for any $t \in \mathbb{R} \backslash\{0\}$,

$$
\left|\varphi_{X}(t)\right| \leq \frac{C}{|t|^{\rho}}
$$

(iii) for any $\rho \in] 0, \frac{d+1}{m}[$, there exists $C>0$ such that for any $t \in \mathbb{R} \backslash\{0\}$,

$$
\left|\varphi_{Y}(t)\right| \leq \frac{C}{|t|^{\rho}}
$$

(iv) $X$ and $Y$ are absolutely continuous with regard to Lebesgue's measure. Their densities are bounded and continuous on $\mathbb{R}$.

Corollary 1 For any $(\alpha, \beta)$, the distribution of $W_{(\alpha, \beta)}^{\mathrm{DT}}$ admits a bounded and continuous density. Its support is the whole real line $\mathbb{R}$.

Corollary 1 is an immediate consequence of Theorem 6 and of the decomposition property (12). The proof of Theorem 6, that follows Liu's method, shows successively that the distributions of $X$ and $Y$ are supported by the whole real line, that the characteristic functions $\varphi_{X}$ and $\varphi_{Y}$ reach the value 1 only at the origin, that they tend to zero at $\pm \infty$ and finally that they are bounded above, in a neighborhood of infinity, by a suitable power function so that a Fourier inversion theorem can apply, revealing the absolute continuity.

We first show a couple of lemmas, Lemma 2 being the first item of Theorem 6.
Lemma 2 The support of both $X$ and $Y$ is the whole real line $\mathbb{R}$.
Proof of Lemma 2 We denote by $\operatorname{Supp}(X)$ and $\operatorname{Supp}(Y)$ the supports of $X$ and $Y$.
Let $x \in \operatorname{Supp}(X)$ and $y \in \operatorname{Supp}(Y)$. Because of fixed-point equation (15), for any $v=\left(v_{1}, \ldots, v_{S+1}\right)$ and $w=\left(w_{1}, \ldots, w_{S+1}\right)$ in $[0,1]$ such that $\sum_{1 \leq k \leq S+1} v_{k}=$ $\sum_{1 \leq k \leq S+1} w_{k}=1$,

$$
\begin{equation*}
\left(x \sum_{k=1}^{a+1} v_{k}^{\sigma}+y \sum_{k=a+2}^{S+1} v_{k}^{\sigma}, x \sum_{k=1}^{c} w_{k}^{\sigma}+y \sum_{k=c+1}^{S+1} w_{k}^{\sigma}\right) \in \operatorname{Supp}(X) \times \operatorname{Supp}(Y) . \tag{19}
\end{equation*}
$$

Apply (19) for $v=w=(t, 0, \ldots, 0,1-t)$ where $t \in[0,1]$; since the segment $[x, y]$ can be written $[x, y]=\left\{t^{\sigma} x+(1-t)^{\sigma} y, t \in[0,1]\right\}$, this entails that $[x, y] \subseteq$ $\operatorname{Supp}(X) \cap \operatorname{Supp}(Y)$. Thus, for any real numbers $x$ and $y$,

$$
\begin{equation*}
(x, y) \in \operatorname{Supp}(X) \times \operatorname{Supp}(Y) \Longleftrightarrow[x, y] \subseteq \operatorname{Supp}(X) \cap \operatorname{Supp}(Y) \tag{20}
\end{equation*}
$$

As a first consequence, $\operatorname{Supp}(X)=\operatorname{Supp}(Y)=:$ Supp. We now proceed in three steps.
Step 1: there exists $\varepsilon>0$ such that $[-\varepsilon, \varepsilon] \subseteq$ Supp.
It suffices to take $x>0$ and $y<0$, respectively, belonging to $\operatorname{Supp}(X)$ and $\operatorname{Supp}(Y)$ —which is possible since $\mathbb{E} X>0$ and $\mathbb{E} Y<0$-and to apply (20), taking $\varepsilon=\min \{x,-y\}$.

Step 2: there exists $\eta>0$ such that, for any $z \in \mathbb{R}$,

$$
z \in \operatorname{Supp} \Longrightarrow(1+\eta) z \in \operatorname{Supp}
$$

Indeed, apply Formula (19) for $v=\left(\frac{1}{2}, \frac{1}{2}, 0, \ldots, 0\right)$. Then, $2^{1-\sigma} z \in \operatorname{Supp}$ as soon as $z \in$ Supp. It suffices to take $\eta=2^{1-\sigma}-1$.

Step 3: the images of $[-\varepsilon, \varepsilon]$ by the iterates of the homothetic transformation $z \mapsto$ $(1+\eta) z$ fill the whole real line.

Lemma 3 For any $t \neq 0,\left|\varphi_{X}(t)\right|<1$ and $\left|\varphi_{Y}(t)\right|<1$.
Proof of Lemma 3 Of course, $\left|\varphi_{X}(t)\right| \leq 1$ for any real number $t$. Assume that $t_{0} \in \mathbb{R}$ satisfies $\left|\varphi_{X}\left(t_{0}\right)\right|=1$. Let $\theta_{0} \in \mathbb{R}$ such that $\mathbb{E}\left(e^{i t_{0} X}\right)=e^{i \theta_{0}}$. Then, almost surely, $e^{i t_{0} X}=e^{i \theta_{0}}$ which is possible only if $t_{0}=0\left(\right.$ and $\theta_{0} \in 2 \pi \mathbb{Z}$ ) since $\operatorname{Supp}(X)=\mathbb{R}$. Same proof for $Y$.

Lemma $4 \lim _{t \rightarrow \pm \infty} \varphi_{X}(t)=0$ and $\lim _{t \rightarrow \pm \infty} \varphi_{Y}(t)=0$
Proof of Lemma 4 This proof and the remainder of the argumentation on the absolute continuity rely on the following equalities that are consequences of the fixed-point equation (15): for any $t \in \mathbb{R}$,

$$
\left\{\begin{array}{l}
\varphi_{X}(t)=\mathbb{E}\left(\varphi_{X}\left(t V_{1}^{\sigma}\right) \cdots \varphi_{X}\left(t V_{a+1}^{\sigma}\right) \varphi_{Y}\left(t V_{a+2}^{\sigma}\right) \cdots \varphi_{Y}\left(t V_{S+1}^{\sigma}\right)\right)  \tag{21}\\
\varphi_{Y}(t)=\mathbb{E}\left(\varphi_{X}\left(t V_{1}^{\sigma}\right) \cdots \varphi_{X}\left(t V_{c}^{\sigma}\right) \varphi_{Y}\left(t V_{c+1}^{\sigma}\right) \cdots \varphi_{Y}\left(t V_{S+1}^{\sigma}\right)\right)
\end{array}\right.
$$

where $V=\left(V_{1}, \ldots, V_{S+1}\right)$ is a Dirichlet-distributed random vector with parameters $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$. These relations are obtained from fixed-point equation (15) by conditioning with respect to $V$. In particular, since all $V_{k}$ are not zero with positive probability, Fatou's Lemma together with Eq. (21) imply that

$$
\left\{\begin{array}{l}
\lim \sup _{t \rightarrow \pm \infty}\left|\varphi_{X}(t)\right| \leq\left(\lim \sup _{t \rightarrow \pm \infty}\left|\varphi_{X}(t)\right|\right)^{a+1} \\
\lim \sup _{t \rightarrow \pm \infty}\left|\varphi_{Y}(t)\right| \leq\left(\lim \sup _{t \rightarrow \pm \infty}\left|\varphi_{Y}(t)\right|\right)^{d+1}
\end{array}\right.
$$

Consequently, since $a \geq 1$ and $d \geq 1$ (the urn is assumed to be large), $\lim \sup _{t \rightarrow \pm \infty}\left|\varphi_{X}(t)\right| \in\{0,1\}$ and the same holds for $\varphi_{Y}$. It remains to show that $\lim \sup _{t \rightarrow \pm \infty}\left|\varphi_{X}(t)\right|=1$ is impossible to get the result for $X$ (a same argument for $Y$ applies as well). The first Eq. (21) implies that

$$
\left|\varphi_{X}(t)\right| \leq \mathbb{E}\left|\varphi_{X}\left(V_{1}^{\sigma} t\right)\right| .
$$

The end of the proof relies on the following idea: denoting by $U$ the uniform distribution on [0, 1], since $V_{1}^{\sigma} \stackrel{\mathcal{L}}{=} U^{m} \leq 1$ almost surely and $E V_{k}^{\sigma}=\frac{1}{m+1}<1$ for any $k$, iterating this last inequality leads to $\lim \sup _{t \rightarrow \pm \infty}\left|\varphi_{X}(t)\right|=0$ which implies the final result. The details, that are rather technical, can be almost literally adapted from Liu's proof of a result of the same kind. See [25], Lemma 3.1, page 93. One can also refer to [12] for a similar argument in a slightly different context.

Lemma 5 For any $\left.\rho^{\prime} \in\right] 0, \frac{1}{m}\left[\right.$, when $t$ tends to $\pm \infty, \varphi_{X}(t) \in O\left(t^{-\rho^{\prime}}\right)$ and $\varphi_{Y}(t)$ $\in O\left(t^{-\rho^{\prime}}\right)$.

Proof of Lemma 5 Let $\varepsilon>0$. Let $T>0$ such that $\left|\varphi_{X}(t)\right| \leq \varepsilon$ and $\left|\varphi_{Y}(t)\right| \leq \varepsilon$ as soon as $|t| \geq T$; the existence of $T$ is guaranteed by Lemma 4. Then, because of (21), for any $t \in \mathbb{R}$,

$$
\left|\varphi_{X}(t)\right| \leq \varepsilon^{S} \mathbb{E}\left|\varphi_{X}\left(V_{S+1}^{\sigma} t\right)\right|+\sum_{k=1}^{S} \mathbb{P}\left(V_{k}^{\sigma}|t| \leq T\right)
$$

Since any $V_{k}^{\sigma} \stackrel{\mathcal{L}}{=} U^{m}$ where $U$ denotes the uniform distribution on [0, 1], this leads to

$$
\left|\varphi_{X}(t)\right| \leq \varepsilon^{S} \mathbb{E}\left|\varphi_{X}\left(U^{m} t\right)\right|+S\left(\frac{T}{|t|}\right)^{\frac{1}{m}}
$$

for any $t \in \mathbb{R} \backslash\{0\}$. Now, for any $\rho \in] 0,1 / m\left[, \mathbb{E}\left(U^{-m \rho}\right)<\infty\right.$ and the former inequality implies that there exists a positive constant $C$ such that for any nonzero $t$,

$$
\left|\varphi_{X}(t)\right| \leq \varepsilon^{S} \mathbb{E}\left|\varphi_{X}\left(U^{m} t\right)\right|+C\left(\frac{1}{|t|}\right)^{\rho}
$$

Thus, the random variables $X$ and $U$ satisfy the assumptions of the Gronwall type Lemma shown in [25], Lemme 3.2 page 93 . Using iterations of the former inequality, one gets for any $n \geq 1$,

$$
\left|\varphi_{X}(t)\right| \leq \varepsilon^{n s} \mathbb{E}\left|\varphi_{X}\left(U_{1}^{m} \ldots U_{n}^{m} t\right)\right|+C|t|^{-\rho} \sum_{k=0}^{n-1}\left(\varepsilon^{s} \mathbb{E}\left(U^{-m \rho}\right)\right)^{k}
$$

which entails that $\left|\varphi_{X}(t)\right| \leq C|t|^{-\rho} /\left(1-\varepsilon^{S} \mathbb{E}\left(U^{-m \rho}\right)\right)$ as soon as $\varepsilon$ is chosen in order that $1-\varepsilon^{S_{\mathbb{E}}}\left(U^{-m \rho}\right)>0$. This implies the result. A same argument is used for $\varphi_{Y}$.

End of the proof of Theorem 6 Let $\rho \in] 0, \frac{a+1}{m}\left[\right.$ and let $\rho^{\prime}=\frac{\rho}{a+1}$. Let $\kappa>0$ such that $\varphi_{X}(t) \leq \kappa|t|^{-\rho^{\prime}}$ for any $t \neq 0$; the existence of $\kappa$ is due to Lemma 5. Applying (21), one gets the successive inequalities

$$
\left|\varphi_{X}(t)\right| \leq \mathbb{E}\left(\prod_{k=1}^{a+1}\left|\varphi_{X}\left(V_{k}^{\sigma} t\right)\right|\right) \leq \frac{\kappa^{a+1}}{|t|^{\rho}} \mathbb{E}\left(\prod_{k=1}^{a+1} V_{k}^{-\sigma \rho^{\prime}}\right)
$$

as soon as the last expectation is defined. Since the random vector $V=\left(V_{1}, \ldots, V_{S+1}\right)$ is Dirichlet distributed with parameters $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$, this expectation can be computed from the Stieltjes transform of $V$ (see the Appendix for the general form of joint moments that can be extended to nonreal powers):

$$
\mathbb{E}\left(\prod_{k=1}^{a+1} V_{k}^{-\sigma \rho^{\prime}}\right)=\frac{\Gamma\left(1+\frac{1}{S}\right)}{\Gamma\left(1+\frac{1}{S}-(a+1) \sigma \rho^{\prime}\right)}\left(\frac{\Gamma\left(\frac{1}{S}-\sigma \rho^{\prime}\right)}{\Gamma\left(\frac{1}{S}\right)}\right)^{a+1}
$$

is finite since $\sigma \rho^{\prime}<\frac{\sigma}{m}=\frac{1}{S}$ and $1+\frac{1}{S}-(a+1) \sigma \rho^{\prime}>1-\frac{a}{S}>0$. Note that $a<S=a+b$ because the urn assumed to be is nontriangular. This proves (ii). The same argument is used for the similar result on $\varphi_{Y}$ (iii).

Since $\frac{a+1}{m}=\frac{a+1}{a-c}>1$, item (ii) implies that the Fourier transform $\varphi_{X}$ of the probability measure of $X$ is integrable. This implies that $X$ admits a bounded continuous function as density. The same result holds for $Y$.

Remark 3 For the continuous-time urn process, as shown in [13], the limit random variables $W_{(1,0)}^{\mathrm{CT}}$ and $W_{(0,1)}^{\mathrm{CT}}$ admit densities as well. These functions have been shown to be infinitely differentiable outside 0 , monotonic on $\mathbb{R}_{<0}$ and $\mathbb{R}_{>0}$, but not bounded around the origin. The different behaviors of $W^{\mathrm{CT}}$ and $W^{\mathrm{DT}}$ have to be related to the martingale connexion (10): when the process starts with one ball, the density of $\xi^{\sigma}$ is not bounded at 0 since $\xi$ is Gamma distributed, with parameter $\frac{1}{S}$.

## 4 Smoothing Transforms

This section is devoted to the existence and the unicity of solutions of the distributional systems (15) and (18). Notice that existence and unicity of solutions of the discrete-time system (15) could be deduced from the general result in Neininger and Rüschendorf [28]; nevertheless, we give hereunder a rapid and autonomous proof of Theorem 7, in order to make explicit the contraction method in the case of large Pólya urns. The same approach is developed in Knape and Neininger [22] in the case of $d$-color Pólya urns. The proof is also reminiscent of the one in Fill and Kapur [16].

When $A$ is a real number, let $\mathcal{M}_{2}(A)$ be the space of probability distributions on $\mathbb{R}$ that have $A$ as expectation and a finite second moment, endowed with a complete metric space structure by the Wasserstein distance. Note first that when $X$ and $Y$ are solutions of (15) or (18) that have, respectively, $B$ and $C$ as expectations, then $c B+b C=0$ (elementary computation). In Theorems 7 and 8 , we prove that when $B$ and $C$ are two real numbers that satisfy $c B+b C=0$, the systems (15) and (18) both have a unique solution in the product metric space $\mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C)$. To do so, we use the Banach contraction method.

Since $(\mathbb{E} X, \mathbb{E} Y)$ is proportional to $(b,-c)$ in both continuous-time and discretetime urn processes (Formulae (7) and (9)), this result shows that the systems (15) and (18) characterize the limit distributions $W_{(1,0)}^{\mathrm{DT}}$ and $W_{(0,1)}^{\mathrm{DT}}$ on one hand, $W_{(1,0)}^{\mathrm{CT}}$ and $W_{(0,1)}^{\mathrm{CT}}$ on the other hand.

### 4.1 The Wasserstein Distance

Let $A \in \mathbb{R}$. The Wasserstein distance on $\mathcal{M}_{2}(A)$ is defined as follows:

$$
d_{W}\left(\mu_{1}, \mu_{2}\right)=\min _{\left(X_{1}, X_{2}\right)}\left(\mathbb{E}\left(X_{1}-X_{2}\right)^{2}\right)^{1 / 2}
$$

where the minimum is taken over random vectors $\left(X_{1}, X_{2}\right)$ on $\mathbb{R}^{2}$ having respective marginal distributions $\mu_{1}$ and $\mu_{2}$; the minimum is attained by the Kantorovich-

Rubinstein Theorem. With this distance, $\mathcal{M}_{2}(A)$ is a complete metric space (see for instance Dudley [14]).

Let $(B, C) \in \mathbb{R}^{2}$. The product space $\mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C)$ is equipped with the product metric, defined (for example) by the distance

$$
d\left(\left(\mu_{1}, \nu_{1}\right),\left(\mu_{2}, \nu_{2}\right)\right)=\max \left\{d_{W}\left(\mu_{1}, \mu_{2}\right), d_{W}\left(\nu_{1}, \nu_{2}\right)\right\} .
$$

Of course, this product remains a complete metric space.

### 4.2 Contraction Method in Discrete Time

Let us recall the fixed-point system (15) satisfied by ( $X^{\mathrm{DT}}, Y^{\mathrm{DT}}$ ), the elementary limits of a large two-color discrete-time Pólya urn process:

$$
\left\{\begin{array}{l}
X \stackrel{\mathcal{L}}{=} \sum_{k=1}^{a+1} V_{k}^{\sigma} X^{(k)}+\sum_{k=a+2}^{S+1} V_{k}^{\sigma} Y^{(k)} \\
Y \stackrel{\mathcal{L}}{=} \sum_{k=1}^{c} V_{k}^{\sigma} X^{(k)}+\sum_{k=c+1}^{S+1} V_{k}^{\sigma} Y^{(k)} .
\end{array}\right.
$$

Let $\mathcal{M}_{2}$ be the space of square-integrable probability measures on $\mathbb{R}$. When $(B, C) \in \mathbb{R}^{2}$, let $K_{1}$ be the function defined on $\mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C)$ by:

$$
\begin{aligned}
K_{1}: \mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C) & \longrightarrow \mathcal{M}_{2} \\
(\mu, v) & \longmapsto \mathcal{L}\left(\sum_{k=1}^{a+1} V_{k}^{\sigma} X^{(k)}+\sum_{k=a+2}^{S+1} V_{k}^{\sigma} Y^{(k)}\right)
\end{aligned}
$$

where $X^{(1)}, \ldots, X^{(a+1)}$ are $\mu$-distributed random variables, $Y^{(a+2)}, \ldots, Y^{(S+1)}$ are $v$-distributed random variables, $V=\left(V_{1}, \ldots, V_{S+1}\right)$ is a Dirichlet-distributed random vector with parameter $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$, the $X^{(k)}, Y^{(k)}$ and $V$ being all independent of each other. Similarly, let $K_{2}$ be defined by

$$
\begin{aligned}
K_{2}: \mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C) & \longrightarrow \mathcal{M}_{2} \\
(\mu, \nu) & \longmapsto \mathcal{L}\left(\sum_{k=1}^{c} V_{k}^{\sigma} X^{(k)}+\sum_{k=c+1}^{S+1} V_{k}^{\sigma} Y^{(k)}\right) .
\end{aligned}
$$

A simple computation shows that if $(\mu, v) \in \mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C)$, then

$$
\mathbb{E} K_{1}(\mu, v)=\frac{(a+1) B+b C}{m+1}
$$

and

$$
\mathbb{E} K_{2}(\mu, v)=\frac{c B+(d+1) C}{m+1}
$$

so that, since $m=a-c=d-b$, the relation $c B+b C=0$ is a sufficient and necessary condition for the product function $\left(K_{1}, K_{2}\right)$ to range $\mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C)$ into itself.

Lemma 6 Let $B$ and $C$ be real numbers that satisfy $c B+b C=0$. Then, the smoothing transform

$$
\begin{aligned}
K: \mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C) & \longrightarrow \mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C) \\
(\mu, \nu) & \longmapsto\left(K_{1}(\mu, v), K_{2}(\mu, \nu)\right)
\end{aligned}
$$

is $\sqrt{\frac{S+1}{2 m+1}}$-Lipschitz. In particular, it is a contraction.
Theorem 7 (i) When $B$ and $C$ are real numbers that satisfy $c B+b C=0$, System (15) has a unique solution in $\mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C)$.
(ii) The pair ( $X^{\mathrm{DT}}, Y^{\mathrm{DT}}$ ) is the unique solution of the distributional System (15) having $\left(\frac{\Gamma\left(\frac{1}{S}\right)}{\Gamma\left(\frac{m+1}{S}\right)} \frac{b}{S},-\frac{\Gamma\left(\frac{1}{S}\right)}{\Gamma\left(\frac{m+1}{S}\right)} \frac{c}{S}\right)$ as expectation and a finite second moment.

Theorem 7 is a direct consequence of Lemma 6 and of Banach's fixed-point theorem.
Proof of Lemma 6 Let $\left(\mu_{1}, \nu_{1}\right)$ and $\left(\mu_{2}, \nu_{2}\right)$ in $\mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C)$. Let $V=\left(V_{1}, \ldots, V_{S+1}\right)$ be a Dirichlet-distributed random vector with parameter $\left(\frac{1}{S}, \ldots, \frac{1}{S}\right)$. Let $X_{1}^{(1)}, \ldots, X_{1}^{(a+1)}$ be $\mu_{1}$-distributed random variables, $Y_{1}^{(a+2)}, \ldots$, $Y_{1}^{(S+1)}$ be $\nu_{1}$-distributed random variables, $X_{2}^{(1)}, \ldots, X_{1}^{(c)}$ be $\mu_{2}$-distributed random variables and $Y_{2}^{(c+1)}, \ldots, Y_{2}^{(S+1)}$ be $\nu_{2}$-distributed random variables, all of them being independent and independent of $V$. Then,

$$
\begin{aligned}
d_{W}\left(K_{1}\left(\mu_{1}, v_{1}\right), K_{1}\left(\mu_{2}, \nu_{2}\right)\right)^{2} \leq & \| \sum_{k=1}^{a+1} V_{k}^{\sigma}\left(X_{1}^{(k)}-X_{2}^{(k)}\right) \\
& +\sum_{k=a+2}^{S+1} V_{k}^{\sigma}\left(Y_{1}^{(k)}-Y_{2}^{(k)}\right) \|_{2}^{2} \\
= & \operatorname{Var}
\end{aligned} \begin{aligned}
& \sum_{k=1}^{a+1} V_{k}^{\sigma}\left(X_{1}^{(k)}-X_{2}^{(k)}\right) \\
& \left.+\sum_{k=a+2}^{S+1} V_{k}^{\sigma}\left(Y_{1}^{(k)}-Y_{2}^{(k)}\right)\right]
\end{aligned}
$$

$$
\begin{aligned}
=\mathbb{E} \operatorname{Var}( & \sum_{k=1}^{a+1} V_{k}^{\sigma}\left(X_{1}^{(k)}-X_{2}^{(k)}\right) \\
& \left.+\sum_{k=a+2}^{S+1} V_{k}^{\sigma}\left(Y_{1}^{(k)}-Y_{2}^{(k)}\right) \mid V\right) \\
& +\operatorname{Var} \mathbb{E}\left(\sum_{k=1}^{a+1} V_{k}^{\sigma}\left(X_{1}^{(k)}-X_{2}^{(k)}\right)\right. \\
& \left.+\sum_{k=a+2}^{S+1} V_{k}^{\sigma}\left(Y_{1}^{(k)}-Y_{2}^{(k)}\right) \mid V\right)
\end{aligned}
$$

thanks to the law of total variance. Since $V=\left(V_{1}, \ldots, V_{S+1}\right)$ is independent of the $X_{j}^{(k)}$ and of the $Y_{j}^{(k)}$, one gets

$$
\begin{aligned}
d_{W}\left(K_{1}\left(\mu_{1}, \nu_{1}\right), K_{1}\left(\mu_{2}, \nu_{2}\right)\right)^{2} \leq & \sum_{k=1}^{a+1} \mathbb{E} V_{k}^{2 \sigma} \operatorname{Var}\left(X_{1}^{(k)}-X_{2}^{(k)}\right) \\
& +\sum_{k=a+2}^{S+1} \mathbb{E} V_{k}^{2 \sigma} \operatorname{Var}\left(Y_{1}^{(k)}-Y_{2}^{(k)}\right) \\
\leq & \operatorname{Var}\left(X_{1}^{(1)}-X_{2}^{(1)}\right) \sum_{k=1}^{a+1} \mathbb{E} V_{k}^{2 \sigma} \\
& +\operatorname{Var}\left(Y_{1}^{(1)}-Y_{2}^{(1)}\right) \sum_{k=a+2}^{S+1} \mathbb{E} V_{k}^{2 \sigma} \\
= & \frac{a+1}{2 m+1}\left\|X_{1}^{(1)}-X_{2}^{(1)}\right\|_{2}^{2}+\frac{b}{2 m+1}\left\|Y_{1}^{(1)}-Y_{2}^{(1)}\right\|_{2}^{2} .
\end{aligned}
$$

Since the inequality holds for any random variables $X_{1}^{(1)}, X_{2}^{(1)}, Y_{1}^{(1)}$ and $Y_{2}^{(1)}$ having respective distributions $\mu_{1}, \mu_{2}, \nu_{1}$ and $\nu_{2}$, this leads to

$$
\begin{aligned}
d_{W}\left(K_{1}\left(\mu_{1}, \nu_{1}\right), K_{1}\left(\mu_{2}, \nu_{2}\right)\right)^{2} & \leq \frac{a+1}{2 m+1} d_{W}\left(\mu_{1}, \mu_{2}\right)^{2}+\frac{b}{2 m+1} d_{W}\left(\nu_{1}, \nu_{2}\right)^{2} \\
& \leq \frac{S+1}{2 m+1} d\left(\left(\mu_{1}, \nu_{1}\right),\left(\mu_{2}, \nu_{2}\right)\right)^{2}
\end{aligned}
$$

A very similar computation shows that

$$
d_{W}\left(K_{2}\left(\mu_{1}, \nu_{1}\right), K_{2}\left(\mu_{2}, \nu_{2}\right)\right)^{2} \leq \frac{S+1}{2 m+1} d\left(\left(\mu_{1}, \nu_{1}\right),\left(\mu_{2}, \nu_{2}\right)\right)^{2}
$$

so that, finally,

$$
d\left(K\left(\mu_{1}, \nu_{1}\right), K\left(\mu_{2}, \nu_{2}\right)\right)^{2} \leq \frac{S+1}{2 m+1} d\left(\left(\mu_{1}, \nu_{1}\right),\left(\mu_{2}, \nu_{2}\right)\right)^{2}
$$

making the proof complete. Note that the assumption $\sigma=\frac{m}{S}>\frac{1}{2}$ guarantees that the Lipschitz constant is in $] 0,1[$.

### 4.3 Contraction Method in Continuous Time

In continuous time, the laws of $X^{\mathrm{CT}}$ and $Y^{\mathrm{CT}}$ are solutions of the following system (cf. (18)):

$$
\left\{\begin{array}{l}
X \stackrel{\mathcal{L}}{=} U^{m}\left(\sum_{k=1}^{a+1} X^{(k)}+\sum_{k=a+2}^{S+1} Y^{(k)}\right) \\
Y \stackrel{\mathcal{L}}{=} U^{m}\left(\sum_{k=1}^{c} X^{(k)}+\sum_{k=c+1}^{S+1} Y^{(k)}\right),
\end{array}\right.
$$

The following theorem, which is the continuous-time version of Theorem 7, can be proved by adapting the arguments of Theorem 7. It is still in the same vein as Neininger and Rüschendorf [28] or Knape and Neininger [22]. Details are left to the reader.

Theorem 8 (i) When $B$ and $C$ are real numbers that satisfy $c B+b C=0$, System (18) has a unique solution in $\mathcal{M}_{2}(B) \times \mathcal{M}_{2}(C)$.
(ii) The pair $\left(X^{\mathrm{CT}}, Y^{\mathrm{CT}}\right)$ is the unique solution of the distributional system (18) having $\left(\frac{b}{S},-\frac{c}{S}\right)$ as expectation and a finite second moment.

## 5 Moments

This section is devoted to the asymptotics of the moments of the limit variables $W^{\text {DT }}$ and $W^{\text {CT }}$. We shall see that they are big but not too much. Observe first that the connexion (10) allows us to study only one of the two cases among discrete or continuous case. We chose to focus on the continuous case, since the fixed-point equation system is slightly easier to deal with. Let us recall here system (18).

$$
\left\{\begin{array}{l}
X \stackrel{\mathcal{L}}{=} U^{m}\left(\sum_{k=1}^{a+1} X^{(k)}+\sum_{k=a+2}^{S+1} Y^{(k)}\right) \\
Y \stackrel{\mathcal{L}}{=} U^{m}\left(\sum_{k=1}^{c} X^{(k)}+\sum_{k=c+1}^{S+1} Y^{(k)}\right),
\end{array}\right.
$$

where $U$ is uniform on $[0,1]$, where $X, X^{(k)}$ and $Y, Y^{(k)}$ are respective copies of $X^{\mathrm{CT}}$ and $Y^{\mathrm{CT}}$, all being independent of each other and of $U$.

Up to now, what is known about the size of these moments is contained in [13] where it is proved that the radius of convergence of the Laplace series of a nontrivial square-integrable solution of (18) is equal to zero. Consequently, by the Hadamard formula for the radius of convergence,

$$
\limsup _{p}\left(\frac{\mathbb{E}|X|^{p}}{p!}\right)^{\frac{1}{p}}=+\infty
$$

In other words, for any constant $C$, for any integer $p_{0} \geq 1$, there exists an integer $p \geq p_{0}$ such that

$$
C^{p} \leq \frac{\mathbb{E}|X|^{p}}{p!} .
$$

The following lemma gives an upperbound for $\frac{\mathbb{E}|X|^{p}}{p!}$. It is the argument leading to Theorem 9 where it is proved that the law of $X$ is determined by its moments.

Lemma 7 If $X$ and $Y$ are square-integrable solutions of (18), they admit absolute moments of all orders $p \geq 1$ and the sequences $\left(\frac{\mathbb{E}|X|^{p}}{p!\log ^{p} p}\right)^{\frac{1}{p}}$ and $\left(\frac{\mathbb{E}|Y|^{p}}{p!\log ^{p} p}\right)^{\frac{1}{p}}$ are bounded.

Proof Since $X$ and $Y$ are square-integrable solutions of (18), their respective expectations $B$ and $C$ satisfy $c B+b C=0$. We may assume that $B C \neq 0$ (otherwise $X$ and $Y$ are trivial). By Theorem 8, the random variables $\frac{b X}{S B}$ and $\frac{-c Y}{S C}$ are, respectively, distributed like $X^{\mathrm{CT}}$ and $Y^{\mathrm{CT}}$ so that they admit moments of all orders.

Let $\varphi(p):=\log ^{p}(p+2)$ and define

$$
u_{p}:=\frac{\mathbb{E}|X|^{p}}{p!\varphi(p)} \text { and } v_{p}:=\frac{\mathbb{E}|Y|^{p}}{p!\varphi(p)}
$$

We show by induction that, for all integer $p \geq 1$, the positive numbers $\left(\frac{\mathbb{E}|X|^{p}}{p!\varphi(p)}\right)^{\frac{1}{p}}$ and $\left(\frac{\mathbb{E}|Y|^{p}}{p!\varphi(p)}\right)^{\frac{1}{p}}$ define bounded sequences. Notice that a similar technique is used in Kahane and Peyrière [21]. Take the power $p$ in the first equation, notice that $\mathbb{E} U^{m p}=\frac{1}{m p+1}$, and isolate the two extreme terms. One gets (remember $S+1=a+1+b$ )

$$
\begin{aligned}
\mathbb{E}|X|^{p} \leq & \frac{1}{m p+1}\left((a+1) \mathbb{E}|X|^{p}+b \mathbb{E}|Y|^{p}\right. \\
& \left.+\sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\
p_{j} \leq p-1}} \frac{p!}{p_{1}!\cdots p_{S+1}!} \mathbb{E}|X|^{p_{1}} \cdots \mathbb{E}|X|^{p_{a+1}} \mathbb{E}|Y|^{p_{a+2}} \cdots \mathbb{E}|Y|^{p_{S+1}}\right)
\end{aligned}
$$

or also

$$
\begin{aligned}
& (m p-a) \mathbb{E}|X|^{p} \leq b \mathbb{E}|Y|^{p} \\
& \quad+\sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\
p_{j} \leq p-1}} \frac{p!}{p_{1}!\cdots p_{S+1}!} \mathbb{E}|X|^{p_{1}} \cdots \mathbb{E}|X|^{p_{a+1}} \mathbb{E}|Y|^{p_{a+2}} \cdots \mathbb{E}|Y|^{p_{S+1}}
\end{aligned}
$$

An analog inequality holds for $\mathbb{E}|Y|^{p}$, leading to the system

Since the eigenvalues of $R$ are $S$ and $m$, and since $\frac{m}{S}>\frac{1}{2}$ (we consider a large urn), then, for all integer $p \geq 2$, the matrix ( $m p I_{2}-R$ ) is invertible. Moreover, its inverse is given by

$$
\left(m p I_{2}-R\right)^{-1}=\frac{1}{(m p-d)(m p-a)-b c}\left(\begin{array}{cc}
m p-d & b \\
c & m p-a
\end{array}\right) .
$$

For all integer $p \geq 2, m p-a>S-a=b \geq 0, m p-d>c \geq 0$ and $(m p-d)(m p-$ $a)-b c>0$, and the inverse of $\left(m p I_{2}-R\right)$ has thus nonnegative coefficients. We can thus rewrite System (22) as

$$
\left(m p I_{2}-R\right)\binom{u_{p}}{v_{p}} \leq t_{p-1},
$$

where the inequality has to be read coefficient by coefficient, and where the two coefficients of vector $t_{p-1}$ are the two sums arisen in System (22). Therefore, since the inverse of $\left(m p I_{2}-R\right)$ has nonnegative coefficients, we get

$$
\binom{u_{p}}{v_{p}} \leq\left(m p I_{2}-R\right)^{-1} t_{p-1} .
$$

Let $p_{0}$ be the smallest positive integer such that for any $p \geq p_{0}$,

$$
\frac{m(p-1)}{(m p-a)(m p-d)-b c}(1+8 \log (p+2))^{S+1} \leq 1
$$

Such a $p_{0}$ exists since the left-hand side goes to 0 when $p$ goes to $+\infty$. Denote

$$
A:=\max _{1 \leq q \leq p_{0}}\left\{\left(u_{q}\right)^{\frac{1}{q}},\left(v_{q}\right)^{\frac{1}{q}}\right\} .
$$

Assume by induction on $p \geq p_{0}+1$ that for every $q \leq p-1,\left(u_{q}\right)^{\frac{1}{q}} \leq A$ and $\left(v_{q}\right)^{\frac{1}{q}} \leq A$. Then,

$$
\left\{\begin{array}{l}
(m p-a) u_{p} \leq b v_{p}+A^{p} \sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\
p_{j} \leq p-1}} \frac{\varphi\left(p_{1}\right) \cdots \varphi\left(p_{S+1}\right)}{\varphi(p)} \\
(m p-d) v_{p} \leq c u_{p}+A^{p} \sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\
p_{j} \leq p-1}} \frac{\varphi\left(p_{1}\right) \cdots \varphi\left(p_{S+1}\right)}{\varphi(p)}
\end{array}\right.
$$

Let

$$
\begin{equation*}
\Phi(p):=\sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\ p_{j} \leq p-1}} \frac{\varphi\left(p_{1}\right) \cdots \varphi\left(p_{S+1}\right)}{\varphi(p)} \tag{23}
\end{equation*}
$$

so that

$$
\left\{\begin{array}{l}
(m p-a) u_{p} \leq b v_{p}+A^{p} \Phi(p) \\
(m p-d) v_{p} \leq c u_{p}+A^{p} \Phi(p)
\end{array}\right.
$$

which implies

$$
u_{p} \leq \frac{m(p-1)}{(m p-a)(m p-d)-b c} A^{p} \Phi(p)
$$

and the same inequality for $v_{p}$ as well. Admit for a while the following lemma.
Lemma 8 For every $p \geq 2, \Phi(p) \leq(1+8 \log (p+2))^{S+1}$.
Consequently,

$$
u_{p} \leq \frac{m(p-1)}{(m p-a)(m p-d)-b c} A^{p}(1+8 \log (p+2))^{S+1}
$$

By definition of $p_{0}$, this implies that $\left(u_{p}\right)^{\frac{1}{p}} \leq A$ and the recurrence holds.
Proof of Lemma 8 The definitions of $\varphi$ and $\Phi$ imply directly that

$$
\begin{aligned}
\Phi(p) & =\sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\
p_{j} \leq p-1}} \frac{\log ^{p_{1}}\left(p_{1}+2\right) \cdots \log ^{p_{S+1}}\left(p_{S+1}+2\right)}{\log ^{p}(p+2)} \\
& =\sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\
p_{j} \leq p-1}}\left(1+\frac{\log \left(1-\frac{p-p_{1}}{p+2}\right)}{\log (p+2)}\right)^{p_{1}} \cdots\left(1+\frac{\log \left(1-\frac{p-p_{S+1}}{p+2}\right)}{\log (p+2)}\right)^{p_{S+1}} .
\end{aligned}
$$

Using $\log (1-u) \leq-u$ for all $u<1$ leads to

$$
\Phi(p) \leq \sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\ p_{j} \leq p-1}}\left(1-\frac{p-p_{1}}{(p+2) \log (p+2)}\right)^{p_{1}} \cdots\left(1-\frac{p-p_{S+1}}{(p+2) \log (p+2)}\right)^{p_{S+1}}
$$

which can be written with an exponential to get, using again $\log (1-u) \leq-u$ :

$$
\Phi(p) \leq \sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\ p_{j} \leq p-1}} \exp \left\{-\frac{p^{2}}{(p+2) \log (p+2)} \sum_{j=1}^{S+1} \frac{p_{j}}{p}\left(1-\frac{p_{j}}{p}\right)\right\}
$$

Let $\psi_{p}(x):=\exp \left(-\frac{p^{2}}{(p+2) \log (p+2)} x(1-x)\right)$, so that

$$
\begin{aligned}
\Phi(p) & \leq \sum_{\substack{p_{1}+\cdots+p_{S+1}=p \\
p_{j} \leq p-1}} \psi_{p}\left(\frac{p_{1}}{p}\right) \cdots \psi_{p}\left(\frac{p_{S+1}}{p}\right) \\
& \leq \sum_{\substack{0 \leq p_{1}, \ldots, p_{S+1} \leq p-1}} \psi_{p}\left(\frac{p_{1}}{p}\right) \cdots \psi_{p}\left(\frac{p_{S+1}}{p}\right) \\
& =\left(\sum_{k=0}^{p-1} \psi_{p}\left(\frac{k}{p}\right)\right)^{S+1} .
\end{aligned}
$$

Elementary calculations lead then to

$$
\sum_{k=0}^{p-1} \psi_{p}\left(\frac{k}{p}\right) \leq 1+p \int_{0}^{1} \psi_{p}(t) d t
$$

and for any $\alpha>0$

$$
\int_{0}^{1} \exp (-\alpha x(1-x)) \mathrm{d} t \leq \frac{4}{\alpha}
$$

so that

$$
\int_{0}^{1} \psi_{p}(t) \mathrm{d} t \leq 4 \frac{(p+2) \log (p+2)}{p^{2}}
$$

and the lemma holds.
The upperbound on the moments obtained in Lemma 7 leads to the following theorem.

Theorem 9 (i) Let $X$ and $Y$ be square-integrable solutions of (18). Then $X$ and $Y$ admit absolute moments of all orders $p \geq 1$ and the probability distributions of $|X|,|Y|, X$ and $Y$ are determined by their moments.
(ii) Let $X$ and $Y$ be square-integrable solutions of (15). Then their Laplace series have a radius of convergence equal to $+\infty$ (and their laws are thus determined by their moments).

Proof (i) By Lemma 7, if $X$ and $Y$ are square-integrable solutions of (18), they admit moments of all orders and, when $p$ is large enough,

$$
\begin{equation*}
\left(\mathbb{E}|X|^{p}\right)^{-\frac{1}{p}} \geq C \frac{(p!)^{-\frac{1}{p}}}{\log p} \tag{24}
\end{equation*}
$$

Besides, by Stirling's formula, when $p$ tends to infinity,

$$
\frac{(p!)^{-\frac{1}{p}}}{\log p} \sim \frac{e}{p \log p}
$$

which is the general term of a Bertrand divergent series. The Carleman's criterion applies, implying that $X$ and $Y$ are moment determined.
(ii) If $X$ and $Y$ are square-integrable solutions of (15) and if $\xi$ is an independent Gamma $\left(\frac{1}{S}\right)$-distributed random variable, then thanks to Proposition $1, \xi^{\sigma} X$ and $\xi^{\sigma} Y$ are square-integrable solutions of (18) so that they both satisfy

$$
\frac{\mathbb{E}\left|\xi^{\sigma} X\right|^{p}}{p!} \leq C^{p} \log ^{p} p
$$

Consequently, for another constant $D$

$$
\frac{\mathbb{E}|X|^{p}}{p!} \leq D^{p} \frac{\log ^{p} p}{\Gamma\left(\sigma p+\frac{1}{S}\right)}
$$

so that the radius of convergence of the Laplace series of $X$ and $Y$ are infinite.
Corollary 2 (i) For any initial composition ( $\alpha, \beta$ ), the limit law $W_{(\alpha, \beta)}^{\mathrm{CT}}$ of a large Pólya urn process is determined by its moments.
(ii) For any initial composition $(\alpha, \beta)$, the Laplace series of $W_{(\alpha, \beta)}^{\mathrm{DT}}$ has an infinite radius of convergence.

Proof ( $i$ ) For elementary initial compositions $(1,0)$ or $(0,1)$, the result is a direct consequence of Theorems 7, 8 and 9. For a general initial composition $(\alpha, \beta)$ in continuous time, notice that decomposition Formula (16) implies that

$$
\left\|W_{(\alpha, \beta)}^{\mathrm{CT}}\right\|_{p} \leq \alpha\left\|W_{(1,0)}^{\mathrm{CT}}\right\|_{p}+\beta\left\|W_{(0,1)}^{\mathrm{CT}}\right\|_{p} .
$$

Since $W_{(1,0)}^{\mathrm{CT}}$ and $W_{(0,1)}^{\mathrm{CT}}$ satisfy (24), $W_{(\alpha, \beta)}^{\mathrm{CT}}$ satisfies Carleman's criterion; it is thus determined by its moments.
(ii) is a direct consequence of Theorem 9 and of decomposition Formula (12).

Remark 4 We have shown throughout the present paper that the random variable $W^{\text {DT }}$ is smoother than $W^{\text {CT }}$ :

- its Laplace series is convergent,
- its Fourier transform is $L^{1}$, and
- its density is continuous and bounded,
which is not true concerning $W^{\mathrm{CT}}$.

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## 6 Appendix: Pólya urns and Dirichlet Distribution

In this section, we deal with results that belong to the "folklore": they are not new neither very difficult, but are nowhere properly gathered, to the best of our knowledge. Proposition 3 goes back to Athreya [3] with different names and a different proof. It is partially given in Blackwell and Kendall [9] for $S=1$ and starting from one ball of each color. The moment method is evocated in Johnson and Kotz book [20]. We detail here a proof to make our paper self-contained.

### 6.1 Dirichlet Distributions

This section gathers some well-known facts on Dirichlet distributions. Besides, we fix notations we use in the sequel.

Let $d \geq 2$ be a natural integer. Let $\Sigma$ be the $(d-1)$-dimensional simplex

$$
\Sigma=\left\{\left(x_{1}, \ldots x_{d}\right) \in[0,1]^{d}, \sum_{k=1}^{d} x_{k}=1\right\}
$$

The following formula is a generalization of the definition of Euler's Bêta function: let $\left(v_{1}, \ldots, v_{d}\right)$ be positive real numbers. Then,

$$
\begin{equation*}
\int_{\Sigma}\left[\prod_{k=1}^{d} x_{k}^{v_{k}-1}\right] d \Sigma\left(x_{1}, \ldots, x_{d}\right)=\frac{\Gamma\left(v_{1}\right) \cdots \Gamma\left(v_{d}\right)}{\Gamma\left(v_{1}+\cdots+v_{d}\right)} \tag{25}
\end{equation*}
$$

where $d \Sigma$ denotes the positive measure on the simplex $\Sigma$, defined by

$$
\begin{aligned}
& f\left(x_{1}, \ldots, x_{d}\right) d \Sigma\left(x_{1}, \ldots, x_{d}\right) \\
& \quad=f\left(x_{1}, \ldots, x_{d-1}, 1-\sum_{k=1}^{d-1} x_{k}\right) \mathbb{1}_{\left\{x \in[0,1]^{d-1}, \sum_{k=1}^{d-1} x_{k} \leq 1\right\}} d x_{1} \cdots d x_{d-1}
\end{aligned}
$$

for any continuous function $f$ defined on $\Sigma$.

By means of this formula, one defines usually the Dirichlet distribution with parameters $\left(v_{1}, \ldots, v_{d}\right)$, denoted by Dirichlet $\left(v_{1}, \ldots, v_{d}\right)$, whose density on $\Sigma$ is given by

$$
\frac{\Gamma\left(v_{1}+\cdots+v_{d}\right)}{\Gamma\left(v_{1}\right) \cdots \Gamma\left(v_{d}\right)}\left[\prod_{k=1}^{d} x_{k}^{\nu_{k}-1}\right] d \Sigma\left(x_{1}, \ldots, x_{d}\right)
$$

In particular, if $D=\left(D_{1}, \ldots, D_{d}\right)$ is a $d$-dimensional random vector which is Dirichlet distributed with parameters $\left(\nu_{1}, \ldots, v_{d}\right)$, then, for any $p=\left(p_{1}, \ldots, p_{d}\right) \in \mathbb{N}^{d}$, the (joint) moment of order $p$ of $D$ is

$$
\mathbb{E}\left(D^{p}\right)=\mathbb{E}\left(D_{1}^{p_{1}} \cdots D_{d}^{p_{d}}\right)=\frac{\Gamma(v)}{\Gamma(v+|p|)} \prod_{k=1}^{d} \frac{\Gamma\left(v_{k}+p_{k}\right)}{\Gamma\left(v_{k}\right)}
$$

where $v=\sum_{k=1}^{d} v_{k}$ and $|p|=\sum_{k=1}^{d} p_{k}$.
Finally, a computation of same kind shows that the [0, 1]-valued random variable $D_{k}$, which is the $k$-th marginal distribution of $D$, is $\operatorname{Beta}\left(v_{k}, v-v_{k}\right)$-distributed i.e. admits the density

$$
\frac{1}{B\left(v_{k}, v-v_{k}\right)} t^{v_{k}-1}(1-t)^{\nu-v_{k}-1} \mathbb{1}_{[0,1]} d t
$$

An alternative description of a Dirichlet distribution can be made by considering a sequence $\left(G_{1}, \ldots, G_{d}\right)$ of Gamma-distributed random variables:

Proposition 2 (cf. Bertoin [7, p. 63]) If $\xi_{1}, \ldots, \xi_{d}$ are d independent Gammadistributed random variables of respective parameters $\left(\nu_{1}, v\right), \ldots,\left(v_{d}, \nu\right)$, if $\xi=$ $\sum_{i=1}^{d} \xi_{i}$, then $\xi$ is Gamma- $\left(\nu_{1}+\cdots+v_{d}, v\right)$ distributed, and the random vector $\left(\frac{\xi_{1}}{\xi}, \ldots, \frac{\xi_{d}}{\xi}\right)$ is Dirichlet $\left(v_{1}, \ldots, v_{d}\right)$-distributed and independent of $\xi$.

### 6.2 Original/Diagonal Pólya Urns

Proposition 3 Let $d \geq 2$ and $S \geq 1$ be integers. Let also $\left(\alpha_{1}, \ldots, \alpha_{d}\right) \in \mathbb{N}^{d} \backslash\{0\}$.
Let $\left(P_{n}\right)_{n \geq 0}$ be the d-color Pólya urn random process having $S I_{d}$ as replacement matrix and $\left(\alpha_{1}, \ldots, \alpha_{d}\right)$ as initial composition. Then, almost surely and in any $\mathrm{L}^{t}$, $t \geq 1$,

$$
\frac{P_{n}}{n S} \underset{n \rightarrow \infty}{\longrightarrow} V
$$

where $V$ is a d-dimensional Dirichlet-distributed random vector, with parameters $\left(\frac{\alpha_{1}}{S}, \ldots, \frac{\alpha_{d}}{S}\right)$.

Remark 5 For any $k \in\{1, \ldots, d\}$, the $k$-th coordinate of $V$ is $\operatorname{Beta}\left(\frac{\alpha_{k}}{S}, \sum_{j \neq k} \frac{\alpha_{j}}{S}\right)$ distributed.

Proof We give here a short autonomous proof. Denote $\alpha=\sum_{k=1}^{d} \alpha_{k} \geq 1$. Conditional expectation at time $n+1$ writes

$$
\mathbb{E}\left(P_{n+1} \mid \mathcal{F}_{n}\right)=\frac{\alpha+(n+1) S}{\alpha+n S} P_{n}
$$

so that $\left(\frac{P_{n}}{\alpha+n S}\right)_{n \geq 0}$ is a $[0,1]^{d}$-valued convergent martingale with mean $\left(\alpha_{1} / \alpha, \ldots\right.$, $\left.\alpha_{d} / \alpha\right)$; let $V$ be its limit. If $f$ is any function defined on $\mathbb{R}^{d}$,

$$
\mathbb{E}\left(f\left(P_{n+1}\right) \mid \mathcal{F}_{n}\right)=\left(I+\frac{\Phi}{\alpha+n S}\right)(f)\left(P_{n}\right)
$$

where

$$
\Phi(f)(v)=\sum_{k=1}^{d} v_{k}\left[f\left(v+S e_{k}\right)-f(v)\right]
$$

( $e_{k}$ is the $k$-th vector in $\mathbb{R}^{d}$ canonical basis and $v=\sum_{k=1}^{d} v_{k} e_{k}$ ). In particular, as can be straightforwardly checked, if $p=\left(p_{1}, \ldots, p_{d}\right) \in \mathbb{N}^{d}$ and $|p|=\sum_{k=1}^{d} p_{k}$, the function

$$
\Gamma_{p}(v)=\prod_{k=1}^{d} \frac{\Gamma\left(\frac{v_{k}}{S}+p_{k}\right)}{\Gamma\left(\frac{v_{k}}{S}\right)},
$$

defined on $\mathbb{R}^{d}$, is an eigenfunction of the operator $\Phi$, associated with the eigenvalue $|p| S$. Consequently, after a direct induction, for any $p \in \mathbb{N}^{d}$,

$$
\mathbb{E}\left(\Gamma_{p}\left(P_{n}\right)\right)=\frac{\Gamma\left(\frac{\alpha}{S}+n+|p|\right)}{\Gamma\left(\frac{\alpha}{S}+n\right)} \cdot \frac{\Gamma\left(\frac{\alpha}{S}\right)}{\Gamma\left(\frac{\alpha}{S}+|p|\right)} \cdot \Gamma_{p}\left(P_{0}\right)
$$

so that, when $n$ tends to infinity, by Stirling's formula,

$$
\mathbb{E}\left(\Gamma_{p}\left(P_{n}\right)\right)=n^{|p|} \cdot \frac{\Gamma\left(\frac{\alpha}{S}\right)}{\Gamma\left(\frac{\alpha}{S}+|p|\right)} \cdot \Gamma_{p}\left(P_{0}\right) \cdot\left(1+O\left(\frac{1}{n}\right)\right) .
$$

Besides, expanding real polynomials $X^{p}=X_{1}^{p_{1}} \cdots X_{d}^{p_{d}}$ in the basis $\left(\Gamma_{p}\right)_{p \in \mathbb{N}^{d}}$, one gets formulae

$$
X^{p}=S^{|p|} \Gamma_{p}+\sum_{\substack{k \in \mathbb{N}^{d} \\|k| \leq|p|-1}} a_{p, k} \Gamma_{k}(X)
$$

where the $a_{p, k}$ are rational numbers. Consequently, when $n$ tends to infinity, one gets the asymptotics

$$
\mathbb{E}\left(\frac{P_{n}}{\alpha+n S}\right)^{p}=\frac{\Gamma\left(\frac{\alpha}{S}\right)}{\Gamma\left(\frac{\alpha}{S}+|p|\right)} \Gamma_{p}\left(P_{0}\right)\left(1+O\left(\frac{1}{n}\right)\right) .
$$

which implies that, for any $p \in \mathbb{N}^{d}$,

$$
\begin{equation*}
\mathbb{E}\left(V^{p}\right)=\frac{\Gamma\left(\frac{\alpha}{S}\right)}{\Gamma\left(\frac{\alpha}{S}+|p|\right)} \prod_{k=1}^{d} \frac{\Gamma\left(\frac{\alpha_{k}}{S}+p_{k}\right)}{\Gamma\left(\frac{\alpha_{k}}{S}\right)} \tag{26}
\end{equation*}
$$

Note that this proves the convergence of the martingale in $\mathrm{L}^{t}$ for all $t \geq 1$. Since a Dirichlet distribution is determined by its moments (because its support is compact), this shows that the law of $V$ is a Dirichlet distribution with parameters $\left(\frac{\alpha_{1}}{S}, \ldots, \frac{\alpha_{d}}{S}\right)$.

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[^1]:    ${ }^{1}$ One admits classically negative values for $a$ and $d$, together with arithmetical conditions on $c$ and $b$. Nevertheless, the paper deals with so-called large urns, for which this never happens.

[^2]:    ${ }^{2}$ The case $\sigma=1 / 2$ is similar to this one, the normalization being $\sqrt{n \log n}$ instead of $\sqrt{n}$.

[^3]:    3 For any positive real $a, \mathcal{E} x p(a)$ denotes the exponential distribution with parameter $a$.

