# A class of models for Bayesian predictive inference 

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In a Bayesian framework, to make predictions on a sequence $X_{1}, X_{2}, \ldots$ of random observations, the inferrer needs to assign the predictive distributions $\sigma_{n}(\cdot)=P\left(X_{n+1} \in \cdot \mid X_{1}, \ldots, X_{n}\right)$. In this paper, we propose to assign $\sigma_{n}$ directly, without passing through the usual prior/posterior scheme. One main advantage is that no prior probability has to be assessed. The data sequence $\left(X_{n}\right)$ is assumed to be conditionally identically distributed (c.i.d.) in the sense of (Ann. Probab. 32 (2004) 2029-2052). To realize this programme, a class $\Sigma$ of predictive distributions is introduced and investigated. Such a $\Sigma$ is rich enough to model various real situations and $\left(X_{n}\right)$ is actually c.i.d. if $\sigma_{n}$ belongs to $\Sigma$. Furthermore, when a new observation $X_{n+1}$ becomes available, $\sigma_{n+1}$ can be obtained by a simple recursive update of $\sigma_{n}$. If $\mu$ is the a.s. weak limit of $\sigma_{n}$, conditions for $\mu$ to be a.s. discrete are provided as well.

Keywords: Bayesian nonparametrics; conditional identity in distribution; exchangeability; predictive distribution; random probability measure; sequential predictions; strategy

## 1. Introduction

The object of this paper is Bayesian predictive inference for a sequence of random observations. Let ( $X_{n}: n \geq 1$ ) be a sequence of random variables with values in a measurable space ( $S, \mathcal{B}$ ). Assuming that $\left(X_{1}, \ldots, X_{n}\right)=x$, for some $n \geq 1$ and $x \in S^{n}$, the problem consists of predicting $X_{n+1}$ based on the observed data $x$. In a Bayesian framework, this means to assess the predictive distribution, say

$$
\sigma_{n}(x)(B)=P\left(X_{n+1} \in B \mid\left(X_{1}, \ldots, X_{n}\right)=x\right) \quad \text { for all } B \in \mathcal{B} .
$$

To address this problem, the $X_{n}$ can be taken to be the coordinate random variables on $S^{\infty}$. Accordingly, in the sequel, we let

$$
X_{n}\left(s_{1}, \ldots, s_{n}, \ldots\right)=s_{n}
$$

for each $n \geq 1$ and each $\left(s_{1}, \ldots, s_{n}, \ldots\right) \in S^{\infty}$. Also, to avoid needless technicalities, $S$ is assumed to be a Borel subset of a Polish space and $\mathcal{B}$ the Borel $\sigma$-field on $S$.

Let $\mathcal{P}$ denote the collection of all probability measures on $\mathcal{B}$. Following Dubins and Savage [15], a strategy is a sequence

$$
\sigma=\left(\sigma_{0}, \sigma_{1}, \ldots\right)
$$

such that

- $\sigma_{0} \in \mathcal{P}$ and $\sigma_{n}=\left\{\sigma_{n}(x): x \in S^{n}\right\}$ is a collection of elements of $\mathcal{P}$;
- The map $x \mapsto \sigma_{n}(x)(B)$ is $\mathcal{B}^{n}$-measurable for fixed $n \geq 1$ and $B \in \mathcal{B}$.

Here, $\sigma_{0}$ should be regarded as the marginal distribution of $X_{1}$ and $\sigma_{n}(x)$ as the conditional distribution of $X_{n+1}$ given that $\left(X_{1}, \ldots, X_{n}\right)=x$.

According to the Ionescu-Tulcea theorem, for any strategy $\sigma$, there is a unique probability measure $P$ on $\left(S^{\infty}, \mathcal{B}^{\infty}\right)$ satisfying

$$
\begin{aligned}
& P\left(X_{1} \in \cdot\right)=\sigma_{0} \quad \text { and } \quad P\left(X_{n+1} \in \cdot \mid\left(X_{1}, \ldots, X_{n}\right)=x\right)=\sigma_{n}(x) \\
& \quad \text { for all } n \geq 1 \text { and } P \text {-almost all } x \in S^{n} .
\end{aligned}
$$

Such a $P$ is denoted as $P_{\sigma}$ in the sequel.
To make predictions on the sequence $\left(X_{n}\right)$, a Bayesian inferrer needs precisely a strategy $\sigma$. The Ionescu-Tulcea theorem establishes that, for any strategy $\sigma$, the predictions based on $\sigma$ are consistent with a unique probability distribution for $\left(X_{n}\right)$.

### 1.1. The standard and non-standard approach for exchangeable data

The data sequence $\left(X_{n}\right)$ is usually assumed to be exchangeable. In that case, there are essentially two procedures for selecting a strategy $\sigma$. For definiteness, we call them the standard approach (SA) and the non-standard approach (NSA). The only reason for using these terms is that the first approach is much more popular than the second. Both approaches can be adopted to make Bayesian predictive inference and both lead to a full specification of the probability distribution of $\left(X_{n}\right)$.

According to SA, to obtain $\sigma$, the inferrer should:

- Select a prior $\pi$, namely, a probability measure on $\mathcal{P}$;
- Calculate the posterior of $\pi$ given that $\left(X_{1}, \ldots, X_{n}\right)=x$, say $\pi_{n}(x)$;
- Evaluate $\sigma$ as

$$
\sigma_{0}(B)=\int_{\mathcal{P}} p(B) \pi(d p) \quad \text { and } \quad \sigma_{n}(x)(B)=\int_{\mathcal{P}} p(B) \pi_{n}(x)(d p) \quad \text { for all } B \in \mathcal{B}
$$

To assess a prior $\pi$ is not an easy task. In addition, once $\pi$ is selected, it is also quite difficult to evaluate the posterior $\pi_{n}(x)$. Frequently, it happens that $\pi_{n}(x)$ cannot be written in closed form but only approximated numerically.

On the other hand, SA is not motivated by prediction alone. Another motivation, possibly the main one, is to make inference on other features of the data distribution, such as a mean, a quantile, a correlation, or more generally some random parameter (possibly, infinite dimensional). In all these cases, the posterior $\pi_{n}(x)$ is fundamental. In short, SA is a cornerstone of Bayesian inference, but, when prediction is the main target, is possibly quite involved.

Instead, NSA entails assigning $\sigma_{n}$ directly, without passing through $\pi$ and $\pi_{n}$. Merely, rather than choosing $\pi$ and then evaluating $\pi_{n}$ and $\sigma_{n}$, the inferrer just selects his/her predictive distribution $\sigma_{n}$. This procedure makes sense because of the Ionescu-Tulcea theorem. See [3,4,8,11,12,16,19,20,24]; see also $[17,25,26,29]$ and references therein.

NSA is in line with de Finetti, Dubins and Savage, among others. Pitman's work is fundamental as well; see, for example, [27] and [28]. In fact, NSA is usually adopted (or at least implicit) in species sampling models; see [24].

Similarly, NSA is used in [19] to obtain a fast online Bayesian prediction. Suppose that $S=\mathbb{R}$ and $\sigma_{n}(x)$ admits a density, with respect to some fixed measure $\lambda$ on $\mathcal{B}$, for all $n \geq 0$ and $x \in S^{n}$. In
[19], the update of predictive distributions is given a nice characterization in terms of copulas. Such a characterization, in turn, allows for making Bayesian predictions through an useful recursive procedure. In a sense, the present paper fits into the framework of [19].

From our point of view, NSA has essentially two merits. Firstly, it requires the assignment of probabilities on observable facts only. The value of the next observation $X_{n+1}$ is actually observable, while $\pi$ and $\pi_{n}$ (being probabilities on $\mathcal{P}$ ) do not deal with observable facts. Secondly, as noted in [19], Section 6 , NSA is much more efficient than SA when prediction is the main goal. In this case, why select the prior $\pi$ explicitly? Rather than wondering about $\pi$, it seems reasonable to reflect on how $X_{n+1}$ is affected by $\left(X_{1}, \ldots, X_{n}\right)$.

Finally, NSA is even more appealing in a nonparametric framework, where selecting a prior with large support is usually difficult.

We discuss an example to make the above remarks clearer.
Example 1 (SA versus NSA). If $\left(X_{n}\right)$ is exchangeable, de Finetti's theorem yields

$$
P\left(X_{1} \in B_{1}, \ldots, X_{n} \in B_{n}\right)=\int_{\Theta} \prod_{i=1}^{n} P_{\theta}\left(B_{i}\right) \pi(d \theta)
$$

for some parameter space $\Theta$, some prior $\pi$ on $\Theta$, and some statistical model

$$
M=\left\{P_{\theta}: \theta \in \Theta\right\} \quad \text { where } P_{\theta} \in \mathcal{P} \text { for each } \theta
$$

In the parametric case, $\Theta$ is a Borel subset of $\mathbb{R}^{k}$ and $M$ is dominated and smooth. In the nonparametric case, $\Theta$ is infinite-dimensional, typically $\Theta=\mathcal{P}$. In both cases, SA entails selecting $\pi$, evaluating the posterior $\pi_{n}$ and calculating $\sigma$ as

$$
\sigma_{n}(x)(B)=\int_{\Theta} P_{\theta}(B) \pi_{n}(x)(d \theta)
$$

In turn, NSA entails selecting $\sigma$ directly, without passing through $\pi$ and $\pi_{n}$.
In our opinion, SA may be unsuitable for prediction even in the parametric framework. Not only it is hard to choose $\pi$, but to evaluate $\pi_{n}$ may be difficult as well. On the contrary, NSA usually takes the available information on the data into account more effectively. In fact, in various practical situations, arguing in terms of strategies is simpler than arguing in terms of priors. An obvious example are Polya urns, where the strategy $\sigma$ is naturally determined by the sampling scheme, while the prior $\pi$ is not. The merits of NSA increase further in the nonparametric framework. In that case, if prediction is the main goal, to assess a prior $\pi$ and evaluate the posterior $\pi_{n}$ is really too expensive.

One more remark is in order. Because of exchangeability, the probability distribution $P$ of $\left(X_{n}\right)$ can be written as above for some $M$ and $\pi$. However, by assigning a strategy $\sigma$, the inferrer identifies $P=$ $P_{\sigma}$, not the pair $(M, \pi)$. In a sense, when applying NSA, the "model uncertainty" about $\theta$ is integrated out by the choice of $\sigma$. This appears reasonable after all, as when making predictions, the relevant object is $\sigma$ not $(M, \pi)$. An intriguing problem, pioneered by Diaconis, Ylvisaker and Freedman, is to give conditions on $\sigma$ implying that the statistical model $M$ underlying $P_{\sigma}$ has a given form, for instance $M$ is an exponential family. Such a problem, however, is not investigated in this paper. See [13,14,16,30] and references therein.

### 1.2. Conditionally identically distributed data

If $\left(X_{n}\right)$ is assumed to be exchangeable, however, NSA has a gap. Given an arbitrary strategy $\sigma$, the Ionescu-Tulcea theorem does not grant exchangeability of $\left(X_{n}\right)$ under $P_{\sigma}$. Therefore, for NSA to ap-
ply, one should first characterize those strategies $\sigma$ which make $\left(X_{n}\right)$ exchangeable under $P_{\sigma}$. A nice characterization is [16], Theorem 3.1. However, the conditions on $\sigma$ for making ( $X_{n}$ ) exchangeable are quite hard to check in real problems. This is possibly one of the reasons why NSA has not yet been developed. Another reason is the lack of constructive procedures for determining $\sigma$. It is precisely this lack which makes SA necessary for prediction, even if analytically more involved.

An obvious way to bypass the gap mentioned in the above paragraph is to weaken the exchangeability assumption. One option is to assume ( $X_{n}$ ) to be conditionally identically distributed (c.i.d.), namely

$$
P\left(X_{k} \in \cdot \mid \mathcal{F}_{n}\right)=P\left(X_{n+1} \in \cdot \mid \mathcal{F}_{n}\right) \quad \text { a.s. for all } k>n \geq 0
$$

where $\mathcal{F}_{n}=\sigma\left(X_{1}, \ldots, X_{n}\right)$ and $\mathcal{F}_{0}$ is the trivial $\sigma$-field.
Roughly speaking, the above condition means that, at each time $n \geq 0$, the future observations $\left(X_{k}: k>n\right)$ are identically distributed given the past $\mathcal{F}_{n}$. Such a condition is actually weaker than exchangeability. Indeed, $\left(X_{n}\right)$ is exchangeable if and only if it is stationary and c.i.d.

We refer to Section 2.1 for more information on c.i.d. sequences. Here, we just mention three reasons for taking c.i.d. data into account.

- It is not hard to characterize the strategies $\sigma$ which make $\left(X_{n}\right)$ c.i.d. under $P_{\sigma}$; see Theorem 3. Therefore, unlike the exchangeable case, NSA can be easily implemented.
- The asymptotic behavior of c.i.d. sequences is analogous to that of exchangeable ones.
- A number of meaningful strategies cannot be used if $\left(X_{n}\right)$ is assumed to be exchangeable, but are available if $\left(X_{n}\right)$ is only required to be c.i.d. See the examples in Sections 4-6.

To support the latter claim, we also note that conditional identity in distribution is a more appropriate assumption than exchangeability in some real problems. Examples occur in various fields, including clinical trials, generalized Polya urns, species sampling models and disease surveillance; see [1,2,5] and [10].

### 1.3. Further notation and conditions (a)-(b)

A kernel (or a random probability measure) on $(S, \mathcal{B})$ is a collection

$$
\alpha=\{\alpha(x): x \in S\}
$$

such that $\alpha(x) \in \mathcal{P}$ for each $x \in S$ and the map $x \mapsto \alpha(x)(B)$ is measurable for fixed $B \in \mathcal{B}$. Here, $\alpha(x)(B)$ denotes the value taken at $B$ by the probability measure $\alpha(x)$.

Let $\sigma_{0} \in \mathcal{P}$ and $\alpha$ a kernel on $(S, \mathcal{B})$. In the sequel, $\sigma_{0}$ and $\alpha$ are such that:
(a) $\sigma_{0}$ is a stationary distribution for $\alpha$, namely,

$$
\sigma_{0}(B)=\int \alpha(x)(B) \sigma_{0}(d x) \quad \text { for all } B \in \mathcal{B}
$$

(b) There is a set $A \in \mathcal{B}$ such that $\sigma_{0}(A)=1$ and

$$
\alpha(x)(B)=\int \alpha(z)(B) \alpha(x)(d z) \quad \text { for all } x \in A \text { and } B \in \mathcal{B} .
$$

Conditions (a)-(b) are not so unusual. For instance, they are satisfied whenever $\alpha$ is a regular conditional distribution for $\sigma_{0}$ given any sub- $\sigma$-field of $\mathcal{B}$; see Lemma 6. In particular, conditions (a)-(b)
trivially hold if

$$
A=S \quad \text { and } \quad \alpha(x)=\delta_{x} \quad \text { for all } x \in S
$$

where $\delta_{x}$ denotes the point mass at $x$.
Finally, if $x=\left(x_{1}, \ldots, x_{n}\right) \in S^{n}$ and $y \in S$, we write $(x, y)$ to denote

$$
(x, y)=\left(x_{1}, \ldots, x_{n}, y\right)
$$

In addition, for any strategy $\sigma$, we let

$$
S^{0}=\{\varnothing\}, \quad \sigma_{0}(\varnothing)=\sigma_{0}, \quad \sigma_{1}(\varnothing, y)=\sigma_{1}(y)
$$

### 1.4. Content of this paper

We aim to develop NSA for c.i.d. data. To this end, we introduce and investigate a class $\Sigma$ of strategies. Such a $\Sigma$ is rich enough to model various real situations and $\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$ for each $\sigma \in \Sigma$. Furthermore, when a new observation $X_{n+1}$ becomes available, $\sigma_{n+1}$ can be obtained from a simple recursive update of $\sigma_{n}$.

Each $\sigma \in \Sigma$ can be described as follows. Fix $\sigma_{0} \in \mathcal{P}$ and a kernel $\alpha$ on $(S, \mathcal{B})$ satisfying conditions (a)-(b). Also, for every $n \geq 0$, fix a measurable function $f_{n}: S^{n+2} \rightarrow[0,1]$ such that

$$
f_{n}(x, y, z)=f_{n}(x, z, y) \quad \text { for all } x \in S^{n} \text { and }(y, z) \in S^{2}
$$

Given $\sigma_{0}, \alpha$ and ( $f_{n}: n \geq 0$ ), a strategy $\sigma$ can be obtained via the recursive equation

$$
\sigma_{n+1}(x, y)(B)=\int \alpha(z)(B) f_{n}(x, y, z) \sigma_{n}(x)(d z)+\alpha(y)(B)\left\{1-\int f_{n}(x, y, z) \sigma_{n}(x)(d z)\right\}
$$

for all $n \geq 0, B \in \mathcal{B}, x \in S^{n}$ and $y \in S$.
We define $\Sigma$ as the collection of all the strategies $\sigma$ obtained as above.
The simplest example corresponds to

$$
f_{n}(x, y, z)=q_{n}(x),
$$

where $q_{n}: S^{n} \rightarrow[0,1]$ is any measurable map (with $q_{0}$ constant). In that case, the recursive equation reduces to

$$
\begin{equation*}
\sigma_{n+1}(x, y)=q_{n}(x) \sigma_{n}(x)+\left\{1-q_{n}(x)\right\} \alpha(y) \tag{1}
\end{equation*}
$$

for all $n \geq 0, x \in S^{n}$ and $y \in S$. (Here, for the sake of simplicity, we are assuming $A=S$ where $A$ is the set involved in condition (b).)

In this specific case, the updating rule is quite transparent: $\sigma_{n+1}(x, y)$ is just a convex combination of the previous predictive distribution $\sigma_{n}(x)$ and the kernel $\alpha$ evaluated in the last observation $y$. In addition, the weight $q_{n}(x)$ does not depend on $y$. Note also that $\sigma_{n+1}(x, y)$ can be written explicitly (and not only in recursive form) as

$$
\sigma_{n+1}(x, y)=\sigma_{0} \prod_{i=0}^{n} q_{i}+\alpha(y)\left(1-q_{n}\right)+\sum_{i=1}^{n} \alpha\left(x_{i}\right)\left(1-q_{i-1}\right) \prod_{j=i}^{n} q_{j}
$$

where $x=\left(x_{1}, \ldots, x_{n}\right) \in S^{n}, y \in S$ and $q_{i}$ is a shorthand notation to denote

$$
q_{i}=q_{i}\left(x_{1}, \ldots, x_{i}\right) .
$$

In general, specifying $f_{n}$ and $\alpha$ suitably, various meaningful strategies can be shown to be members of $\Sigma$. Some of these strategies are well known and some are new (in the sense that, to our knowledge, they have not been proposed to date). Examples of the former are the predictive distributions of Dirichlet sequences, species sampling sequences and generalized Polya urns. Examples of the latter are the strategies of Sections 5-6.

Another nice feature of $\Sigma$ is that it also includes diffuse strategies, and this fact may be useful to model real situations. We recall that a probability measure is diffuse if it vanishes on singletons, and a strategy $\sigma$ is diffuse if $\sigma_{n}(x)$ is diffuse for all $n \geq 0$ and $x \in S^{n}$.

In addition to introducing $\Sigma$, our main contributions are Theorems 4-5 and Theorems 18-20. The former state that $\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$ for each $\sigma \in \Sigma$, while the latter deal with the asymptotics of $\sigma_{n}$. A few words should be spent on Theorem 18.

Let $X_{1}^{*}, X_{2}^{*}, \ldots$ denote the (finite or infinite) sequence of distinct values corresponding to the observations $X_{1}, X_{2}, \ldots$ If ( $X_{n}$ ) is c.i.d. under $P_{\sigma}$, where $\sigma$ is any strategy (possibly not belonging to $\Sigma$ ), there is a random probability measure $\mu$ on $(S, \mathcal{B})$ such that

$$
\sigma_{n}(B) \xrightarrow{\text { a.s. }} \mu(B) \quad \text { for every fixed } B \in \mathcal{B}
$$

where "a.s." stands for " $P_{\sigma}$-a.s."; see Section 2.1. Theorem 18 states that

$$
\mu \stackrel{\text { a.s. }}{=} \sum_{k} W_{k} \delta_{X_{k}^{*}},
$$

for some random weights $W_{k} \geq 0$ such that $\sum_{k} W_{k}=1$, if and only if

$$
\lim _{n} P_{\sigma}\left(X_{n} \neq X_{i} \text { for each } i<n\right)=0
$$

Furthermore, $W_{k}$ admits the representation

$$
W_{k} \stackrel{\text { a.s. }}{=} \lim _{n} \frac{1}{n} \sum_{i=1}^{n} 1_{\left\{X_{i}=X_{k}^{*}\right\}} .
$$

By applying Theorem 18 to $\sigma \in \Sigma$, it is not hard to give conditions on $f_{n}$ and $\alpha$ implying that $\mu$ is a.s. discrete. Conditions for $X_{1}^{*}, X_{2}^{*}, \ldots$ to be i.i.d. and independent of the weights $W_{1}, W_{2}, \ldots$ are given as well.

It is worth noting that Theorem 18 holds true for any strategy $\sigma$ which makes $\left(X_{n}\right)$ c.i.d. Hence, Theorem 18 extends a known fact concerning exchangeability to all c.i.d. sequences; see, for example, [23].

In addition to the results quoted above, other main contributions of this paper are the examples included in Sections 4-6. In our opinion, these examples should support the fact that $\Sigma$ is rich enough to cover a wide range of problems.

## 2. Preliminaries

### 2.1. Conditional identity in distribution

C.i.d. sequences have been introduced in [5] and [21] and then investigated in various papers; see, for example, $[1-4,6,7,10,18]$. Here, we just recall a few basic facts.

Let $\left(\mathcal{G}_{n}: n \geq 0\right)$ be a filtration and $\left(Y_{n}: n \geq 1\right)$ a sequence of $S$-valued random variables. Then, $\left(Y_{n}\right)$ is c.i.d. with respect to $\left(\mathcal{G}_{n}\right)$ if it is adapted to $\left(\mathcal{G}_{n}\right)$ and

$$
P\left(Y_{k} \in \cdot \mid \mathcal{G}_{n}\right)=P\left(Y_{n+1} \in \cdot \mid \mathcal{G}_{n}\right) \quad \text { a.s. for all } k>n \geq 0 .
$$

When $\left(\mathcal{G}_{n}\right)$ is the canonical filtration of $\left(Y_{n}\right)$, the filtration is not mentioned at all and $\left(Y_{n}\right)$ is just called c.i.d. From a result in [21], $\left(Y_{n}\right)$ is exchangeable if and only if it is stationary and c.i.d.

Let $\left(Y_{n}\right)$ be c.i.d., $\mathcal{G}_{n}=\sigma\left(Y_{1}, \ldots, Y_{n}\right)$, and

$$
\mu_{n}=\frac{1}{n} \sum_{i=1}^{n} \delta_{Y_{i}}
$$

the empirical measure. In a sense, the asymptotic behavior of $\left(Y_{n}\right)$ is similar to that of an exchangeable sequence. This claim can be supported by two facts.

First, there is a random probability measure $\mu$ on $(S, \mathcal{B})$ satisfying

$$
\mu_{n}(B) \xrightarrow{\text { a.s. }} \mu(B) \quad \text { for every fixed } B \in \mathcal{B} \text {. }
$$

As a consequence, for fixed $n \geq 0$ and $B \in \mathcal{B}$, one obtains

$$
\begin{aligned}
E\left\{\mu(B) \mid \mathcal{G}_{n}\right\} & =\lim _{m} E\left\{\mu_{m}(B) \mid \mathcal{G}_{n}\right\} \\
& =\lim _{m} \frac{1}{m} \sum_{i=n+1}^{m} P\left(Y_{i} \in B \mid \mathcal{G}_{n}\right)=P\left(Y_{n+1} \in B \mid \mathcal{G}_{n}\right) \quad \text { a.s. }
\end{aligned}
$$

Thus, as in the exchangeable case, the predictive distribution $P\left(Y_{n+1} \in \cdot \mid \mathcal{G}_{n}\right)$ can be written as $E\left\{\mu(\cdot) \mid \mathcal{G}_{n}\right\}$, where $\mu$ is the a.s. weak limit of the empirical measures $\mu_{n}$. In particular, for each $B \in \mathcal{B}$, the martingale convergence theorem implies

$$
\begin{equation*}
P\left(Y_{n+1} \in B \mid \mathcal{G}_{n}\right)=E\left\{\mu(B) \mid \mathcal{G}_{n}\right\} \xrightarrow{\text { a.s. }} \mu(B) . \tag{2}
\end{equation*}
$$

Second, $\left(Y_{n}\right)$ is asymptotically exchangeable, in the sense that

$$
\left(Y_{n}, Y_{n+1}, \ldots\right) \rightarrow\left(Z_{1}, Z_{2}, \ldots\right) \quad \text { in distribution, as } n \rightarrow \infty,
$$

where $\left(Z_{n}\right)$ is an exchangeable sequence. Moreover, $\left(Z_{n}\right)$ is directed by $\mu$, namely

$$
P\left(Z_{1} \in B_{1}, \ldots, Z_{k} \in B_{k}\right)=E\left\{\prod_{i=1}^{k} \mu\left(B_{i}\right)\right\}
$$

for all $k \geq 1$ and $B_{1}, \ldots, B_{k} \in \mathcal{B}$.
The role played by $\mu$ is not as crucial as in the exchangeable case, since the probability distribution of $\left(Y_{n}\right)$ is not completely determined by $\mu$; see Example 17. Nevertheless, $\mu$ is a meaningful
random parameter for $\left(Y_{n}\right)$. In fact, $\mu(B)$ is the long run frequency of the events $\left\{Y_{n} \in B\right\}$. Similarly, because of (2), $\mu(B)$ can be regarded as the asymptotically optimal predictor of the event $\{$ the next observation belongs to $B\}$. And finally, $\mu$ is the directing measure of the exchangeable limit sequence $\left(Z_{n}\right)$.

### 2.2. Stationarity, reversibility and characterizations

We first recall some definitions. Let $\tau \in \mathcal{P}$ and $\alpha=\{\alpha(x): x \in S\}$ a kernel on $(S, \mathcal{B})$. Then:

- $\tau$ is a stationary distribution for $\alpha$ if

$$
\int \alpha(x)(B) \tau(d x)=\tau(B) \quad \text { for all } B \in \mathcal{B}
$$

- $\alpha$ is reversible with respect to $\tau$ if

$$
\int_{A} \alpha(x)(B) \tau(d x)=\int_{B} \alpha(x)(A) \tau(d x) \quad \text { for all } A, B \in \mathcal{B}
$$

- $\alpha$ is a regular conditional distribution for $\tau$ given $\mathcal{G}$, where $\mathcal{G} \subset \mathcal{B}$ is a sub- $\sigma$-field, if $x \mapsto \alpha(x)(B)$ is $\mathcal{G}$-measurable and

$$
\int_{A} \alpha(x)(B) \tau(d x)=\tau(A \cap B) \quad \text { for all } A \in \mathcal{G} \text { and } B \in \mathcal{B}
$$

Since $S$ is nice (it is in fact a Borel subset of a Polish space), for any sub- $\sigma$-field $\mathcal{G} \subset \mathcal{B}$, there exists a $\tau$-a.s. unique regular conditional distribution for $\tau$ given $\mathcal{G}$; see e.g. [22], page 107. Note also that reversibility implies stationarity (just take $A=S$ ) but not conversely. In addition, $\tau$ is a stationary distribution for $\alpha$ provided $\alpha$ is a regular conditional distribution for $\tau$ (take $A=S$ again).

We next characterize exchangeable and c.i.d. sequences in terms of strategies.
Theorem 2 ([16], Theorem 3.1). For any strategy $\sigma,\left(X_{n}\right)$ is exchangeable under $P_{\sigma}$ if and only if
(i) The kernel $\left\{\sigma_{n+1}(x, y): y \in S\right\}$ is reversible with respect to $\sigma_{n}(x)$ for all $n \geq 0$ and $P_{\sigma}$-almost all $x \in S^{n}$;
(ii) $\sigma_{n}(x)=\sigma_{n}(f(x))$ for all $n \geq 2$, all permutations $f$ on $S^{n}$ and $P_{\sigma}$-almost all $x \in S^{n}$.

To deal with the c.i.d. case, it suffices to drop condition (ii) and to replace "reversible" with "stationary" in condition (i).

Theorem 3 ([6], Theorem 3.1). For any strategy $\sigma,\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$ if and only if
(i*) The kernel $\left\{\sigma_{n+1}(x, y): y \in S\right\}$ has stationary distribution $\sigma_{n}(x)$ for all $n \geq 0$ and $P_{\sigma}$-almost all $x \in S^{n}$.

An obvious consequence of Theorem 3 is that $\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$ whenever $\left\{\sigma_{n+1}(x, y): y \in S\right\}$ has stationary distribution $\sigma_{n}(x)$ for all $n \geq 0$ and all $x \in C^{n}$, where $C \in \mathcal{B}$ is any set with $\sigma_{0}(C)=1$.

Theorem 3 also suggests how to assess a c.i.d. sequence stepwise. First, select $\sigma_{0} \in \mathcal{P}$, the marginal distribution of $X_{1}$. Then, choose a kernel $\left\{\sigma_{1}(y): y \in S\right\}$ with stationary distribution $\sigma_{0}$, where $\sigma_{1}(y)$ is the conditional distribution of $X_{2}$ given $X_{1}=y$. Next, for each $x \in S$, select a kernel $\left\{\sigma_{2}(x, y): y \in S\right\}$ with stationary distribution $\sigma_{1}(x)$, where $\sigma_{2}(x, y)$ is the conditional distribution of $X_{3}$ given $X_{1}=x$ and $X_{2}=y$. And so on. In other terms, for getting a c.i.d. sequence, it is sufficient to assign a kernel with a given stationary distribution at each step.

## 3. A sequential updating rule

Our starting point is the following simple fact.
Theorem 4. Let $\tau \in \mathcal{P}$ and $f: S^{2} \rightarrow[0,1]$ a measurable symmetric function. Fix a kernel $\alpha=\{\alpha(x)$ : $x \in S\}$ on $(S, \mathcal{B})$ and define

$$
\beta(x)(B)=\int \alpha(z)(B) f(x, z) \tau(d z)+\alpha(x)(B) \int(1-f(x, z)) \tau(d z)
$$

for all $x \in S$ and $B \in \mathcal{B}$. Then, $\beta=\{\beta(x): x \in S\}$ is a kernel on $(S, \mathcal{B})$. Moreover:

- If $\tau$ is stationary for $\alpha$, then $\tau$ is stationary for $\beta$;
- If $\alpha(x)=\delta_{x}$ for all $x \in S$, then $\beta$ is reversible with respect to $\tau$.

Proof. Let $\phi(x)=\int f(x, z) \tau(d z)$. If $\phi(x)=0$, then $\beta(x)$ is clearly a probability measure on $\mathcal{B}$. If $\phi(x) \in(0,1]$,

$$
\beta(x)(B)=\phi(x) \frac{\int \alpha(z)(B) f(x, z) \tau(d z)}{\phi(x)}+(1-\phi(x)) \alpha(x)(B)
$$

Hence, $\beta(x) \in \mathcal{P}$ for all $x \in S$. Further, for fixed $B \in \mathcal{B}$, the map $x \mapsto \beta(x)(B)$ is measurable because of Fubini's theorem. Thus, $\beta$ is a kernel on $(S, \mathcal{B})$.

Next, suppose $\tau$ is stationary for $\alpha$. Since $f(x, z)=f(z, x)$, one obtains

$$
\begin{aligned}
\int \beta(x)(B) \tau(d x)= & \iint \alpha(z)(B) f(x, z) \tau(d z) \tau(d x) \\
& +\int \alpha(x)(B) \tau(d x)-\int \alpha(x)(B) \phi(x) \tau(d x) \\
= & \int \alpha(z)(B) \int f(z, x) \tau(d x) \tau(d z)+\tau(B)-\int \alpha(x)(B) \phi(x) \tau(d x) \\
= & \int \alpha(z)(B) \phi(z) \tau(d z)+\tau(B)-\int \alpha(x)(B) \phi(x) \tau(d x)=\tau(B)
\end{aligned}
$$

for all $B \in \mathcal{B}$. Thus, $\tau$ is stationary for $\beta$.
Finally, if $\alpha(x)=\delta_{x}$, then

$$
\begin{aligned}
\int_{A} \beta(x)(B) \tau(d x)= & \iint 1_{A}(x) 1_{B}(z) f(x, z) \tau(d z) \tau(d x) \\
& +\int 1_{A}(x) 1_{B}(x) \tau(d x)-\int 1_{A}(x) 1_{B}(x) \phi(x) \tau(d x)
\end{aligned}
$$

for all $A, B \in \mathcal{B}$. It follows that

$$
\begin{aligned}
& \int_{A} \beta(x)(B) \tau(d x)-\int_{B} \beta(x)(A) \tau(d x) \\
& \quad=\iint 1_{A}(x) 1_{B}(z) f(x, z) \tau(d z) \tau(d x)-\iint 1_{B}(x) 1_{A}(z) f(x, z) \tau(d z) \tau(d x) \\
& \quad=\int 1_{B}(z) \int 1_{A}(x) f(z, x) \tau(d x) \tau(d z)-\int 1_{B}(x) \int 1_{A}(z) f(x, z) \tau(d z) \tau(d x)=0 .
\end{aligned}
$$

Thus, $\beta$ is reversible with respect to $\tau$.

Heuristically, in the special case $\alpha(x)=\delta_{x}$, the idea underlying $\beta$ reminds of Metropolis' algorithm. Starting from a state $x$, one first selects a new state $z$ according to $\tau$, and then goes to $z$ or remains in $x$ with probabilities $f(x, z)$ and $1-f(x, z)$, respectively. This naive idea can be adapted to an arbitrary kernel $\alpha$ as follows. First, select $z$ according to $\tau$. Then, the new state $y$ is drawn from $\alpha(z)$ with probability $f(x, z)$, or from $\alpha(x)$ with probability $1-f(x, z)$. From our point of view, however, what is meaningful is that this idea provides a simple updating procedure.

Next, as in Section 1.4, fix $\sigma_{0} \in \mathcal{P}$, a kernel $\alpha$ on $(S, \mathcal{B})$ and a sequence of measurable functions $f_{n}: S^{n+2} \rightarrow[0,1]$ such that

$$
f_{n}(x, y, z)=f_{n}(x, z, y) \quad \text { for all } n \geq 0, x \in S^{n} \text { and }(y, z) \in S^{2} .
$$

The kernel $\alpha$ is assumed to satisfy conditions (a)-(b) of Section 1.3. We recall that (a)-(b) are automatically true if $\alpha$ is a regular conditional distribution for $\sigma_{0}$ given any sub- $\sigma$-field $\mathcal{G} \subset \mathcal{B}$; see Lemma 6 below. In particular, conditions (a)-(b) hold if $\alpha(x)=\delta_{x}$ for all $x \in S$.

Given $\sigma_{0}, \alpha$ and ( $f_{n}: n \geq 0$ ), define a strategy $\sigma$ according to

$$
\sigma_{n+1}(x, y)(B)=\int \alpha(z)(B) f_{n}(x, y, z) \sigma_{n}(x)(d z)+\alpha(y)(B)\left\{1-\int f_{n}(x, y, z) \sigma_{n}(x)(d z)\right\}
$$

for all $n \geq 0, x \in S^{n}, y \in S$ and $B \in \mathcal{B}$.
Note that, when a new observation $y$ becomes available, $\sigma_{n+1}(x, y)$ is just a recursive update of $\sigma_{n}(x)$.

Let $\Sigma$ denote the collection of all the strategies $\sigma$ obtained in this way, for $\sigma_{0}, \alpha$ and ( $f_{n}: n \geq 0$ ) varying. Each $\sigma \in \Sigma$ makes ( $X_{n}$ ) c.i.d.

Theorem 5. Let $\sigma \in \Sigma$. Then, $\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$. Moreover, if $\alpha(x)=\delta_{x}$ for all $x \in S$, then

$$
\begin{equation*}
P_{\sigma}\left[\left(X_{n+1}, X_{n+2}\right) \in \cdot \mid \mathcal{F}_{n}\right]=P_{\sigma}\left[\left(X_{n+2}, X_{n+1}\right) \in \cdot \mid \mathcal{F}_{n}\right] \quad \text { a.s. } \tag{3}
\end{equation*}
$$

for all $n \geq 0$, where $\mathcal{F}_{0}$ is the trivial $\sigma$-field and $\mathcal{F}_{n}=\sigma\left(X_{1}, \ldots, X_{n}\right)$.

Proof. We show that there is $C \in \mathcal{B}$ such that $\sigma_{0}(C)=1$ and $\left\{\sigma_{n+1}(x, y): y \in S\right\}$ has stationary distribution $\sigma_{n}(x)$ for all $n \geq 0$ and all $x \in C^{n}$. By the remark after Theorem 3, this implies that ( $X_{n}$ ) is c.i.d. under $P_{\sigma}$.

Let $A \in \mathcal{B}$ be the set involved in condition (b). Define

$$
A_{0}=A \quad \text { and } \quad A_{n+1}=\left\{x \in A_{n}: \alpha(x)\left(A_{n}\right)=1\right\} \quad \text { for all } n \geq 0 .
$$

If $\sigma_{0}\left(A_{n}\right)=1$ for some $n \geq 0$, condition (a) yields

$$
\int \alpha(x)\left(A_{n}\right) \sigma_{0}(d x)=\sigma_{0}\left(A_{n}\right)=1
$$

which in turn implies $\sigma_{0}\left(A_{n+1}\right)=1$. Since $\sigma_{0}\left(A_{0}\right)=\sigma_{0}(A)=1$, by induction, one obtains $\sigma_{0}\left(A_{n}\right)=1$ for each $n \geq 0$. Let

$$
C=\bigcap_{n=0}^{\infty} A_{n} .
$$

If $x \in C$, then $\alpha(x)\left(A_{n}\right)=1$ for all $n$, so that $\alpha(x)(C)=1$. Also, $C \subset A$ and $\sigma_{0}(C)=1$. To summarize, $C$ satisfies

$$
\sigma_{0}(C)=1, \quad \alpha(x)(C)=1 \quad \text { and } \quad \int \alpha(z)(B) \alpha(x)(d z)=\alpha(x)(B) \quad \text { for all } x \in C \text { and } B \in \mathcal{B}
$$

Next, if $\sigma_{n}(x)(C)=1$ for some $n \geq 0$ and all $x \in C^{n}$, then

$$
\begin{aligned}
\sigma_{n+1}(x, y)(C) & =\int_{C} \alpha(z)(C) f_{n}(x, y, z) \sigma_{n}(x)(d z)+\alpha(y)(C)\left\{1-\int f_{n}(x, y, z) \sigma_{n}(x)(d z)\right\} \\
& =\int_{C} f_{n}(x, y, z) \sigma_{n}(x)(d z)+1-\int f_{n}(x, y, z) \sigma_{n}(x)(d z)=1 \quad \text { for all }(x, y) \in C^{n+1}
\end{aligned}
$$

Arguing by induction again, $\sigma_{0}(C)=1$ implies

$$
\sigma_{n}(x)(C)=1 \quad \text { for all } n \geq 0 \text { and all } x \in C^{n}
$$

Finally, fix $(x, y) \in C^{n+1}$. Since $\sigma_{n}(x)(C)=1$,

$$
\begin{aligned}
\int \alpha(v)(B) \sigma_{n+1}(x, y)(d v)= & \int_{C} \int \alpha(v)(B) \alpha(z)(d v) f_{n}(x, y, z) \sigma_{n}(x)(d z) \\
& +\left\{1-\int f_{n}(x, y, z) \sigma_{n}(x)(d z)\right\} \int \alpha(v)(B) \alpha(y)(d v) \\
= & \int_{C} \alpha(z)(B) f_{n}(x, y, z) \sigma_{n}(x)(d z) \\
& +\left\{1-\int f_{n}(x, y, z) \sigma_{n}(x)(d z)\right\} \alpha(y)(B) \\
= & \sigma_{n+1}(x, y)(B) \quad \text { for all } B \in \mathcal{B} .
\end{aligned}
$$

Therefore, $\sigma_{n+1}(x, y)$ is a stationary distribution for the kernel $\alpha$. By Theorem $4, \sigma_{n+1}(x, y)$ is still stationary for the kernel $\left\{\sigma_{n+2}(x, y, z): z \in S\right\}$.

This concludes the proof that $\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$. To conclude the proof of the whole theorem, suppose $\alpha(x)=\delta_{x}$ for all $x \in S$. Then, condition (3) is a direct consequence of Theorem 4 and the following well-known fact. Let $X$ and $Z$ be $S$-valued random variables, $\tau$ the probability distribution of $X$, and $\gamma=\{\gamma(x): x \in S\}$ a regular version of the conditional distribution of $Z$ given $X$. Then,

$$
(X, Z) \sim(Z, X) \quad \Leftrightarrow \quad \gamma \text { is reversible with respect to } \tau .
$$

Condition (3) is stronger than the c.i.d. condition. As an example, (3) implies

$$
\left(X_{i}, X_{j}\right) \sim\left(X_{j}, X_{i}\right) \quad \text { for all } i \neq j
$$

and this may fail for an arbitrary c.i.d. sequence; see e.g. [6], Example 3. Therefore, when $\alpha(x)=\delta_{x}$, the updating procedure of this section yields a special type of c.i.d. sequences. On the other hand, just because of condition (3), Theorem 2 implies

$$
\left(X_{n}\right) \text { is exchangeable under } P_{\sigma} \quad \Leftrightarrow \quad \sigma_{n}(x)=\sigma_{n}(f(x))
$$

for all $n \geq 2$, all permutations $f$ on $S^{n}$ and $P_{\sigma}$-almost all $x \in S^{n}$. In other terms, if $\sigma \in \Sigma$ and $\alpha(x)=$ $\delta_{x}$, the exchangeability of ( $X_{n}$ ) can be easily characterized.

Finally, we turn to conditions (a)-(b). The next result is helpful for finding a kernel $\alpha$ satisfying (a)-(b).

Lemma 6. If $\alpha=\{\alpha(x): x \in S\}$ is a regular conditional distribution for $\sigma_{0}$ given a sub- $\sigma$-field $\mathcal{G} \subset \mathcal{B}$, then $\alpha$ satisfies conditions (a)-(b).

Proof. Condition (a) (that is, $\sigma_{0}$ stationary for $\alpha$ ) has been already noted in Section 2.2. In turn, the proof of (b) essentially agrees with that of [9], Lemma 10, but we report it for completeness. Let $\mathcal{G}_{0}$ be the $\sigma$-field over $S$ generated by the maps $z \mapsto \alpha(z)(B)$ for all $B \in \mathcal{B}$. Then, $\alpha$ is also a regular conditional distribution for $\sigma_{0}$ given $\mathcal{G}_{0}$. In addition, since $\mathcal{B}$ is countably generated, $\mathcal{G}_{0}$ is countably generated as well. Hence, there is $A \in \mathcal{B}$ such that $\sigma_{0}(A)=1$ and

$$
\alpha(x)(B)=\delta_{x}(B) \quad \text { for all } x \in A \text { and } B \in \mathcal{G}_{0} .
$$

Fix $x \in A$ and $B \in \mathcal{B}$. Since the map $z \mapsto \alpha(z)(B)$ is $\mathcal{G}_{0}$-measurable, one obtains

$$
\int \alpha(z)(B) \alpha(x)(d z)=\int \alpha(z)(B) \delta_{x}(d z)=\alpha(x)(B) .
$$

## 4. Examples: Discrete strategies

From now on, we fix $\sigma_{0} \in \mathcal{P}$ and a sequence

$$
q_{n}: S^{n} \rightarrow[0,1], \quad n \geq 0
$$

of measurable functions (with $q_{0}$ constant).
Moreover, in this section, we let

$$
\begin{aligned}
\alpha(x) & =\delta_{x} \quad \text { for all } x \in S \quad \text { and } \\
f_{n}(x, y, z) & =q_{n}(x) \quad \text { for all } x \in S^{n} \text { and }(y, z) \in S^{2} .
\end{aligned}
$$

With this choice of $f_{n}$, the calculation of $\sigma_{n}(x)$ is straightforward. Writing

$$
x=\left(x_{1}, \ldots, x_{n}\right) \quad \text { and } \quad q_{i}=q_{i}\left(x_{1}, \ldots, x_{i}\right),
$$

one obtains

$$
\begin{equation*}
\sigma_{n}(x)=\sigma_{0} \prod_{i=0}^{n-1} q_{i}+\delta_{x_{n}}\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1} \delta_{x_{i}}\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j} . \tag{4}
\end{equation*}
$$

The strategy (4) is connected to Beta-GOS processes, as meant in [1]. If $\sigma_{0}$ is diffuse, the $q_{i}$ have the following interpretation. Let $x=\left(x_{1}, \ldots, x_{n}\right)$. Since $\sigma_{0}\left(\left\{x_{1}, \ldots, x_{n}\right\}\right)=0$ and $\delta_{x_{i}}\left(\left\{x_{1}, \ldots, x_{n}\right\}\right)=1$ for $i \leq n$, it follows that
$P_{\sigma}\left(X_{n+1}=X_{i}\right.$ for some $\left.i \leq n \mid\left(X_{1}, \ldots, X_{n}\right)=x\right)=\sigma_{n}(x)\left(\left\{x_{1}, \ldots, x_{n}\right\}\right)$

$$
=\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1}\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j}=1-\prod_{i=0}^{n-1} q_{i} .
$$

More importantly, by specifying the $q_{i}$ suitably, a lot of meaningful predictive distributions can be obtained from (4).

Example 7 (Exponential smoothing). If $q_{i}=q$ for all $i \geq 0$, where $q \in[0,1]$ is any constant, formula (4) reduces to

$$
\sigma_{n}(x)=q^{n} \sigma_{0}+(1-q) \sum_{i=1}^{n} q^{n-i} \delta_{x_{i}}
$$

see also [2]. Roughly speaking, this choice of $\sigma$ makes sense when the inferrer has only vague opinions on the dependence structure of the data, and yet he/she feels that the weight of the $i$-th observation $x_{i}$ should be a decreasing function of $n-i$. Note that $\sigma_{n}(x)$ is not invariant under permutations of $x$, so that $\left(X_{n}\right)$ fails to be exchangeable under $P_{\sigma}$. Yet, $\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$ because of Theorem 5.

Example 8 (Dirichlet sequences). If $q_{i}=\frac{i+c}{i+1+c}$ for some constant $c>0$, formula (4) yields

$$
\sigma_{n}(x)=\frac{c \sigma_{0}+\sum_{i=1}^{n} \delta_{x_{i}}}{n+c}
$$

These are the predictive distributions of a Dirichlet sequence. In this case, $\left(X_{n}\right)$ is exchangeable under $P_{\sigma}$.

Example 9 (Latent variables). Suppose $q_{i}$ of the form

$$
q_{i}=q_{i}\left(x_{1}, \ldots, x_{i} ; \lambda_{1}, \ldots, \lambda_{i}\right)
$$

where $\lambda_{1}, \ldots, \lambda_{i}$ take values in a Borel set $T$ of some Polish space.
To cover this situation, fix a Borel probability measure $\sigma_{0}^{*}$ on $S \times T$ such that

$$
\sigma_{0}^{*}(B \times T)=\sigma_{0}(B) \quad \text { for all } B \in \mathcal{B},
$$

and define

$$
\sigma_{n}^{*}\left[\left(x_{1}, \lambda_{1}\right), \ldots,\left(x_{n}, \lambda_{n}\right)\right]=\sigma_{0}^{*} \prod_{i=0}^{n-1} q_{i}+\delta_{\left(x_{n}, \lambda_{n}\right)}\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1} \delta_{\left(x_{i}, \lambda_{i}\right)}\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j}
$$

Marginalizing $\sigma_{n}^{*}$, one obtains

$$
\sigma_{n}^{*}\left[\left(x_{1}, \lambda_{1}\right), \ldots,\left(x_{n}, \lambda_{n}\right)\right](B \times T)=\sigma_{n}(x)(B) \quad \text { for all } B \in \mathcal{B}
$$

where $\sigma_{n}(x)$ is given by (4). Also, up to replacing $S$ with $S \times T$, Theorem 5 applies to the strategy $\sigma^{*}$. More precisely, let $P_{\sigma^{*}}$ be the probability measure on the Borel sets of $(S \times T)^{\infty}$ induced by $\sigma^{*}$ and let $\Lambda_{n}$ be the $n$-th coordinate random variable on $T^{\infty}$. Then, the sequence ( $X_{n}, \Lambda_{n}$ ) is c.i.d. under $P_{\sigma^{*}}$. In other terms, $\left(X_{n}\right)$ is c.i.d. (under $P_{\sigma^{*}}$ ) even if $q_{i}$ depends on the latent variables $\lambda_{1}, \ldots, \lambda_{i}$.

A last remark, motivated by the next Example 10, is as follows. The above argument still applies if $\lambda_{1}$ is a known constant and

$$
q_{i}=q_{i}\left(x_{1}, \ldots, x_{i} ; \lambda_{1}, \ldots, \lambda_{i}, \lambda_{i+1}\right) .
$$

In fact, since $\lambda_{1}$ is constant, $q_{0}=q_{0}\left(\lambda_{1}\right)$ is constant as well. Thus, it suffices to replace $\left(x_{n}, \lambda_{n}\right)$ with $\left(x_{n}, \lambda_{n+1}\right)$, namely, to define $\sigma_{n}^{*}$ as

$$
\sigma_{n}^{*}\left[\left(x_{1}, \lambda_{2}\right), \ldots,\left(x_{n}, \lambda_{n+1}\right)\right]=\sigma_{0}^{*} \prod_{i=0}^{n-1} q_{i}+\delta_{\left(x_{n}, \lambda_{n+1}\right)}\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1} \delta_{\left(x_{i}, \lambda_{i+1}\right)}\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j} .
$$

Arguing as above, the sequence $\left(X_{n}, \Lambda_{n+1}\right)$ is c.i.d. under $P_{\sigma^{*}}$ and

$$
\sigma_{n}^{*}\left[\left(x_{1}, \lambda_{2}\right), \ldots,\left(x_{n}, \lambda_{n+1}\right)\right](B \times T)=\sigma_{n}(x)(B) \quad \text { for all } B \in \mathcal{B}
$$

where $\sigma_{n}(x)$ is given by (4).
Example 10 (Generalized Polya urns). An urn contains $a>0$ white balls and $b>0$ black balls. At each time $n \geq 1$, one ball is taken out and then replaced together with $D_{n}$ more balls of the same color. In the classical scheme, $D_{n}=d$ for all $n$ where $d \geq 0$ is a fixed constant. Here, instead, $\left(D_{n}\right)$ is any sequence of non-negative random variables.

Let $Y_{n}$ be the indicator of the event $\{$ white ball at time $n\}$. Following [5], Example 1.3, it is natural to let

$$
P\left(Y_{n+1}=1 \mid Y_{1}, \ldots, Y_{n}, D_{1}, \ldots, D_{n}\right)=\frac{a+\sum_{i=1}^{n} D_{i} Y_{i}}{a+b+\sum_{i=1}^{n} D_{i}} \quad \text { a.s. }
$$

Assuming $D_{1}$ constant, this is a special case of Example 9. Take in fact $S=\{0,1\}, T=[0, \infty)$, and $\sigma_{0}^{*}$ a Borel probability on $S \times T$ such that

$$
\sigma_{0}^{*}(\{1\} \times[0, \infty))=\frac{a}{a+b}
$$

Then, it suffices to let

$$
q_{i}\left(x_{1}, \ldots, x_{i} ; \lambda_{1}, \ldots, \lambda_{i}, \lambda_{i+1}\right)=\frac{a+b+\sum_{j=1}^{i} \lambda_{j}}{a+b+\sum_{j=1}^{i+1} \lambda_{j}} \quad \text { for all } i \geq 0
$$

## 5. Examples: Diffuse strategies

In this section, we still let $f_{n}(x, y, z)=q_{n}(x)$ but $\alpha=\{\alpha(x): x \in S\}$ is any kernel on ( $S, \mathcal{B}$ ) satisfying conditions (a)-(b). For the sake of simplicity, we suppose that (b) holds with $A=S$ (this can actually be assumed without loss of generality). We denote by $\sigma \in \Sigma$ the strategy induced by $\sigma_{0}, \alpha$ and ( $q_{n}: n \geq 0$ ). Such a $\sigma$ can be written as

$$
\sigma_{n}(x)=\sigma_{0} \prod_{i=0}^{n-1} q_{i}+\alpha\left(x_{n}\right)\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1} \alpha\left(x_{i}\right)\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j}
$$

for all $n \geq 1$ and $x \in S^{n}$, where $q_{i}=q_{i}\left(x_{1}, \ldots, x_{i}\right)$.
To our knowledge, none of the strategies exhibited in this section have ever been proposed before. Two more remarks are in order.

First, $\sigma$ is diffuse whenever $\sigma_{0}$ and $\alpha$ are diffuse. Instead, the strategy (4), as well as many popular strategies, has a discrete part in correspondence with the observed data.

Second, let $\mathcal{H} \subset \mathcal{B}$ be a countable partition of $S$ and let $H(x)$ denote the unique $H \in \mathcal{H}$ which includes the point $x \in S$. For definiteness, suppose $\sigma_{0}(H)>0$ for all $H \in \mathcal{H}$. Then, a simple kernel satisfying conditions (a)-(b) is

$$
\begin{equation*}
\alpha(x)=\sum_{H \in \mathcal{H}} 1_{H}(x) \sigma_{0}(\cdot \mid H)=\sigma_{0}(\cdot \mid H(x)) \quad \text { for all } x \in S \tag{5}
\end{equation*}
$$

Let us turn to specific examples.

Example 11 (Examples 7 and 8 continued). For $x=\left(x_{1}, \ldots, x_{n}\right) \in S^{n}$, the strategies of Examples 7 and 8 turn into

$$
\sigma_{n}(x)=q^{n} \sigma_{0}+(1-q) \sum_{i=1}^{n} q^{n-i} \alpha\left(x_{i}\right) \quad \text { and } \quad \sigma_{n}(x)=\frac{c \sigma_{0}+\sum_{i=1}^{n} \alpha\left(x_{i}\right)}{n+c}
$$

respectively. For definiteness, we focus on the second one, which can be viewed as a version of the predictive distributions of Dirichlet sequences.

As a first example, fix a countable partition $\mathcal{H}$ of $S$ and define $\alpha$ according to (5). Then,

$$
\sigma_{n}(x)=\frac{c \sigma_{0}+\sum_{i=1}^{n} \sigma_{0}\left(\cdot \mid H\left(x_{i}\right)\right)}{n+c} .
$$

Note that $\sigma_{n}(x)$ is absolutely continuous with respect to $\sigma_{0}$ for all $n \geq 0$ and $x \in S^{n}$. Such a strategy $\sigma$ could be reasonable when $H\left(x_{i}\right)$ is the basic information coming from the observation $x_{i}$. Roughly speaking, due to the precision of the measuring tool, one is actually observing an element of $\mathcal{H}$ rather than a point of $S$.

For a more elaborate example, take $S=\mathbb{R}^{2}$ and denote by $\mathcal{R}$ the Borel $\sigma$-field on $\mathbb{R}$ (so that $\mathcal{B}=\mathcal{R}^{2}$ ). Fix a probability measure $r$ on $\mathcal{R}$ and define

$$
\sigma_{0}(A \times B)=\int_{A} N(u)(B) r(d u) \quad \text { for all } A, B \in \mathcal{R},
$$

where $N(u)=N(u, 1)$ is the Gaussian law on $\mathcal{R}$ with mean $u$ and variance 1 . Then, a kernel satisfying conditions (a)-(b) is

$$
\alpha(u, v)=\delta_{u} \times N(u) \quad \text { for all }(u, v) \in \mathbb{R}^{2} .
$$

Thus, if a point in the plane is selected through $\alpha(u, v)$, the abscissa agrees with $u$ a.s. while the ordinate is distributed according to $N(u)$. Using this $\alpha$, one obtains

$$
\sigma_{n}(x)(A \times B)=\frac{c \sigma_{0}(A \times B)+\sum_{i=1}^{n} 1_{A}\left(u_{i}\right) N\left(u_{i}\right)(B)}{n+c} \quad \text { for all } A, B \in \mathcal{R},
$$

where the $i$-th observation $x_{i}$ is written as $x_{i}=\left(u_{i}, v_{i}\right)$. This strategy $\sigma$ comes into play when the basic information which stems from $x_{i}=\left(u_{i}, v_{i}\right)$ is the abscissa $u_{i}$. Note that $\sigma_{n}(x)(A \times B)$ is small if $n$ is large but $u_{i} \notin A$ for all $i$. In a sense, the classical Dirichlet prediction scheme is preserved as regards the abscissas $u_{1}, \ldots, u_{n}$ of the observed data. Note also that $N(u)$ could be replaced by $Q(u)$, where $\{Q(u): u \in \mathbb{R}\}$ is any measurable collection of probabilities on $\mathcal{R}$.

Next, given a countable class $G$ of measurable maps $g: S \rightarrow S$, say that $\sigma_{0}$ is $G$-invariant if

$$
\sigma_{0}\left(g^{-1} B\right)=\sigma_{0}(B) \quad \text { for all } g \in G \text { and } B \in \mathcal{B} .
$$

In this case, the inferrer may wish his/her predictions were $G$-invariant as well.
Example 12 (Invariant strategies). Suppose $\sigma_{0}$ is $G$-invariant and

$$
\mathcal{G}=\left\{B \in \mathcal{B}: g^{-1} B=B \text { for all } g \in G\right\} .
$$

As noted in Section 2.2, since $S$ is nice, there is a regular conditional distribution $\alpha=\{\alpha(x): x \in S\}$ for $\sigma_{0}$ given $\mathcal{G}$. Because of Lemma $6, \alpha$ satisfies conditions (a)-(b). By standard arguments, since $G$ is countable and $\mathcal{B}$ countably generated, $\alpha$ can be taken in such a way that $\alpha(x)$ is $G$-invariant for all $x \in S$. In turn, this implies that $\sigma_{n}(x)$ is $G$-invariant for all $n \geq 0$ and $x \in S^{n}$.

As a simple example, let $S=\mathbb{R}$ and $\sigma_{0}$ symmetric. Take $G=\{g\}$ where $g(x)=-x$, and

$$
\alpha(x)=\frac{\delta_{x}+\delta_{-x}}{2} .
$$

Then,

$$
\begin{aligned}
\sigma_{n}(x)= & \sigma_{0} \prod_{i=0}^{n-1} q_{i}+\frac{1}{2}\left(\delta_{x_{n}}\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1} \delta_{x_{i}}\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j}\right) \\
& +\frac{1}{2}\left(\delta_{-x_{n}}\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1} \delta_{-x_{i}}\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j}\right)
\end{aligned}
$$

is a symmetric strategy which makes $\left(X_{n}\right)$ c.i.d.
As a further example, let $S=T^{d}$ and $\sigma_{0}$ exchangeable, where $T$ is a Borel subset of a Polish space. Take $G$ as the set of all permutations of $T^{d}$, and

$$
\alpha(x)=\frac{\sum_{g \in G} \delta_{g(x)}}{d!} .
$$

Then, $\sigma$ is an exchangeable strategy which makes $\left(X_{n}\right)$ c.i.d.
The strategy $\lambda$ obtained in the last example does not belong to $\Sigma$. However, $\lambda$ comes from essentially the same idea of $\Sigma$ and $\left(X_{n}\right)$ is c.i.d. under $P_{\lambda}$.

Example 13 (Another strategy dominated by $\sigma_{0}$ ). For each $n \geq 1$, take a countable partition $\mathcal{H}_{n}$ of $S$ and assume

$$
\mathcal{H}_{n} \subset \mathcal{B}, \quad \mathcal{H}_{n+1} \text { finer than } \mathcal{H}_{n} \text { and } \sigma_{0}(H)>0 \quad \text { for all } H \in \mathcal{H}_{n}
$$

To avoid trivialities, assume also that $q_{n}>0$ for all $n \geq 0$.
For every $n \geq 1$ and $\tau \in \mathcal{P}$, a kernel $\alpha_{n}=\left\{\alpha_{n}(x): x \in S\right\}$ which admits $\tau$ as a stationary distribution is

$$
\alpha_{n}(x)=\sum_{H \in \mathcal{H}_{n}} 1_{H}(x) \tau(\cdot \mid H)=\tau\left(\cdot \mid H_{n}(x)\right)
$$

(Here, $H_{n}(x)$ is the unique $H \in \mathcal{H}_{n}$ such that $x \in H$ and we tacitly assumed $\tau(H)>0$ for all $H \in \mathcal{H}_{n}$.) Let us define a strategy $\lambda$ as follows. Let $\lambda_{0}=\sigma_{0}$ and

$$
\lambda_{1}(x)=q_{0} \sigma_{0}+\left(1-q_{0}\right) \sigma_{0}\left(\cdot \mid H_{1}(x)\right) \quad \text { for all } x \in S .
$$

By Theorem 4, $\lambda_{0}$ is a stationary distribution for the kernel $\left\{\lambda_{1}(x): x \in S\right\}$. Next, for every $(x, y) \in S^{2}$, define

$$
\lambda_{2}(x, y)=q_{1}(x) \lambda_{1}(x)+\left(1-q_{1}(x)\right) \lambda_{1}(x)\left(\cdot \mid H_{2}(y)\right) .
$$

The kernel $\left\{\lambda_{2}(x, y): y \in S\right\}$ admits $\lambda_{1}(x)$ as a stationary distribution. Moreover, since $\mathcal{H}_{2}$ is finer than $\mathcal{H}_{1}$, one obtains

$$
\lambda_{1}(x)\left(B \mid H_{2}(y)\right)=\sigma_{0}\left(B \mid H_{2}(y)\right) \quad \text { for all } B \in \mathcal{B} .
$$

Therefore, $\lambda_{2}(x, y)$ can be written as

$$
\lambda_{2}(x, y)=q_{0} q_{1}(x) \sigma_{0}+\left(1-q_{0}\right) q_{1}(x) \sigma_{0}\left(\cdot \mid H_{1}(x)\right)+\left(1-q_{1}(x)\right) \sigma_{0}\left(\cdot \mid H_{2}(y)\right)
$$

In general, for every $n \geq 1$ and $x=\left(x_{1}, \ldots, x_{n}\right) \in S^{n}$, define

$$
\lambda_{n}(x)=\sigma_{0} \prod_{i=0}^{n-1} q_{i}+\sigma_{0}\left(\cdot \mid H_{n}\left(x_{n}\right)\right)\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1} \sigma_{0}\left(\cdot \mid H_{i}\left(x_{i}\right)\right)\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j}
$$

where $q_{i}$ stands for $q_{i}\left(x_{1}, \ldots, x_{i}\right)$. Arguing as above, it is easily seen that, for a fixed $x \in S^{n}$, the kernel $\left\{\lambda_{n+1}(x, y): y \in S\right\}$ admits $\lambda_{n}(x)$ as a stationary distribution. Hence, Theorem 3 implies that $\left(X_{n}\right)$ is c.i.d. under $P_{\lambda}$.

The strategy $\lambda$ is reminiscent of (4). As a matter of fact, $\lambda$ agrees with (4) up to replacing $\sigma_{0}(\cdot \mid$ $\left.H_{i}\left(x_{i}\right)\right)$ with $\delta_{x_{i}}$. Furthermore, the partitions $\mathcal{H}_{n}$ can be chosen such that

$$
\{x\}=\bigcap_{n} H_{n}(x) \quad \text { for each } x \in S
$$

Unlike (4), however, $\lambda_{n}(x)$ is absolutely continuous with respect to $\sigma_{0}$ for all $n \geq 0$ and $x \in S^{n}$.

## 6. Examples: Other choices of $\boldsymbol{f}_{\boldsymbol{n}}$

In the previous examples, $f_{n}(x, y, z)=q_{n}(x)$ does not depend on $(y, z)$. This is not so in the present section. We let $\alpha(x)=\delta_{x}$ and we denote by $\sigma \in \Sigma$ the strategy induced by $\sigma_{0}, \alpha$ and ( $f_{n}: n \geq 0$ ). Once again, to our knowledge, the strategies obtained in this section have never been proposed before.

Example 14 (Separating sets). For each $n \geq 0$ and $x \in S^{n}$, take a set $A_{n}(x) \in \mathcal{B}$ and define

$$
f_{n}(x, y, z)=1_{A_{n}(x)}(y) 1_{A_{n}(x)}(z)+1_{A_{n}^{c}(x)}(y) 1_{A_{n}^{c}(x)}(z)
$$

where $A_{n}^{c}(x)$ is the complement of $A_{n}(x)$. Thus, $f_{n}(x, y, z)=0$ or $f_{n}(x, y, z)=1$ according to whether $y$ and $z$ can, or cannot, be separated by the set $A_{n}(x)$. A direct calculation shows that

$$
\sigma_{n+1}(x, y)=\sigma_{n}(x)\left(A_{n}(x)\right) \sigma_{n}(x)\left(\cdot \mid A_{n}(x)\right)+\sigma_{n}(x)\left(A_{n}^{c}(x)\right) \delta_{y} \quad \text { if } y \in A_{n}(x)
$$

where the first summand on the right is meant to be 0 in case $\sigma_{n}(x)\left(A_{n}(x)\right)=0$. Similarly,

$$
\sigma_{n+1}(x, y)=\sigma_{n}(x)\left(A_{n}^{c}(x)\right) \sigma_{n}(x)\left(\cdot \mid A_{n}^{c}(x)\right)+\sigma_{n}(x)\left(A_{n}(x)\right) \delta_{y} \quad \text { if } y \notin A_{n}(x) .
$$

According to the heuristic interpretation of Section 3, such a strategy $\sigma$ can be described as follows. At time $n+1$, after observing $(x, y) \in S^{n+1}$, the inferrer selects a new state $z$ according to $\sigma_{n}(x)$. Then, he/she remains in $y$ or goes to $z$ according to whether $y$ and $z$ are, or are not, separated by $A_{n}(x)$. This could be reasonable, for instance, if the inferrer has some reason to request

$$
\sigma_{n+1}(x, y)\left(A_{n}(x)\right)=1_{A_{n}(x)}(y) .
$$

Example 15 (Decreasing functions of the distance). In the spirit of Example 14, let

$$
f_{n}(x, y, z)=g_{n}[x, d(y, z)]
$$

where $d$ is the distance on $S$ and $g_{n}: S^{n} \times[0, \infty) \rightarrow[0,1]$ a measurable function such that

$$
g_{n}(x, t)<g_{n}(x, s)<g_{n}(x, 0)=1 \quad \text { for all } x \in S^{n} \text { and } 0<s<t .
$$

Then, $\sigma$ can be attached an interpretation similar to Example 14. Again, after observing $(x, y) \in S^{n+1}$, the inferrer selects a new state $z$ according to $\sigma_{n}(x)$. Then, he/she goes to $z$ with probability $f_{n}(x, y, z)$ or remains in $y$ with probability $1-f_{n}(x, y, z)$. Moreover, the chance of reaching $z$ starting from $y$ is a decreasing function of $d(y, z)$ and is 1 if and only if $y=z$.

Example 16 (Ehrenfest-like models). Theorem 4 still works if the assumption $f \leq 1$ is weakened. More specifically, define $\beta$ according to Theorem 4 with $\alpha(x)=\delta_{x}$ and $f$ a measurable symmetric function such that $0 \leq f \leq c$, where $c$ is any constant. Then, $\beta$ is a reversible kernel provided $\beta(x)(B) \geq 0$ for all $x \in S$ and $B \in \mathcal{B}$. Note that the latter condition is trivially true if $c \leq 1$.

As an example, take $S=\{0,1\}$ and $f_{n}$ a non-negative function on $S^{n+2}$ such that $f_{n}(x, y, z)=$ $f_{n}(x, z, y)$. If

$$
\begin{equation*}
f_{n}(x, 0,1)-1 \leq f_{n}(x, 0,1) \sigma_{n}(x)(\{1\}) \leq 1 \quad \text { for all } n \geq 0 \text { and } x \in S^{n}, \tag{6}
\end{equation*}
$$

then $\sigma_{n}(x)(B) \geq 0$ for all $n, x$ and $B$. Hence, $\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$ whenever condition (6) holds. On the other hand, if $f_{n}(x, 0,1)>1$, then

$$
\sigma_{n+1}(x, y)(\{y\})=\sigma_{n}(x)(\{y\})+\left(1-f_{n}(x, 0,1)\right) \sigma_{n}(x)(\{1-y\})<\sigma_{n}(x)(\{y\}) .
$$

In other terms, observing $y$ at step $n+1$ makes the probability of $y$ at step $n+2$ strictly less than the probability of $y$ at step $n+1$. This may look counterintuitive but it makes sense in some problems.

Think of two water-containers $C_{0}$ and $C_{1}$. At each time $n \geq 1$, either $C_{0}$ or $C_{1}$ is selected and a part of its water is transferred into the other one. The total quantity of water, say $w$, remains constant in time. The data are the selected containers. To model this situation, it is quite natural to let $S=\{0,1\}$ and

$$
\lambda_{n}(x)(\{y\})=\frac{\text { quantity of water in } C_{y} \text { after observing } x}{w}
$$

for all $n \geq 0, x \in S^{n}$ and $y \in S$. Such a strategy $\lambda$ belongs to $\Sigma$ under some assumptions on the quantity of water moving from one container to the other. For instance suppose that, after observing $(x, y)$ for
some $x \in S^{n}$ and $y \in S$, the quantity of water transferred from $C_{y}$ into $C_{1-y}$ is

$$
\lambda_{n}(x)(\{1-y\})^{2} \lambda_{n}(x)(\{y\}) w .
$$

Then, $\lambda \in \Sigma$. In fact, $\lambda$ is induced by $\lambda_{0},\left\{\delta_{x}: x \in S\right\}$ and

$$
f_{n}(x, 0,1)=1+\lambda_{n}(x)(\{0\}) \lambda_{n}(x)(\{1\}) .
$$

## 7. Discreteness of the limit of $\sigma_{\boldsymbol{n}}$

This section investigates the limit $\mu$ of $\sigma_{n}$ as $n \rightarrow \infty$. It is split into two subsections. The first introduces a sequence of random variables the predictive distributions of which are given by (4), while the second includes the main results.

### 7.1. An explicit construction

Let $\sigma$ be the strategy (4). To better understand the meaning of $\sigma$, it may be useful to build a sequence $\left(Y_{n}\right)$ of random variables satisfying $Y_{1} \sim \sigma_{0}$ and

$$
\begin{align*}
P\left(Y_{n+1} \in \cdot \mid Y_{1}, \ldots, Y_{n}\right) & =\sigma_{n}\left(Y_{1}, \ldots, Y_{n}\right) \\
& =\sigma_{0} \prod_{i=0}^{n-1} q_{i}+\delta_{Y_{n}}\left(1-q_{n-1}\right)+\sum_{i=1}^{n-1} \delta_{Y_{i}}\left(1-q_{i-1}\right) \prod_{j=i}^{n-1} q_{j} \quad \text { a.s. for } n \geq 1 \tag{7}
\end{align*}
$$

where $q_{i}=q_{i}\left(Y_{1}, \ldots, Y_{i}\right)$. One such $\left(Y_{n}\right)$ is provided by [4].
Let ( $T_{n}: n \geq 1$ ) and ( $U_{i, j}: j \geq 1,0 \leq i<j$ ) be random variables such that:
(j) $\left(T_{n}\right)$ is an i.i.d. sequence of $S$-valued random variables with $T_{1} \sim \sigma_{0}$;
(ji) $\left(U_{i, j}\right)$ is an i.i.d. array of $[0,1]$-valued random variables with $U_{0,1}$ uniformly distributed on $[0,1]$;
( jjj$)\left(T_{n}\right)$ is independent of $\left(U_{i, j}\right)$.
Using ( $T_{n}$ ) and ( $U_{i, j}$ ) as building blocks, the sequence $\left(Y_{n}\right)$ is obtained as follows.
Let $Y_{1}=T_{1}$. Then, define $Y_{2}=T_{2}$ or $Y_{2}=Y_{1}$ according to whether $U_{0,1} \leq q_{0}$ or $U_{0,1}>q_{0}$. At step $n+1$, after $Y_{1}, \ldots, Y_{n}$ have been defined, let

$$
\begin{array}{ll}
Y_{n+1}=T_{n+1} & \text { if } U_{i, n} \leq q_{i}\left(Y_{1}, \ldots, Y_{i}\right) \text { for all } 0 \leq i<n, \\
Y_{n+1}=Y_{i+1} & \text { if } U_{i, n}>q_{i}\left(Y_{1}, \ldots, Y_{i}\right) \text { and } U_{j, n} \leq q_{j}\left(Y_{1}, \ldots, Y_{j}\right) \\
& \text { for some } 0 \leq i<n \text { and all } j>i .
\end{array}
$$

It is not hard to verify that $Y_{1} \sim \sigma_{0}$ and condition (7) holds; see [4], Lemma 3.

### 7.2. Asymptotics

Let $s=\left(s_{1}, \ldots, s_{n}, \ldots\right)$ denote a point of $S^{\infty}$. For any strategy $\sigma$ which makes $\left(X_{n}\right)$ c.i.d., there is a random probability measure $\mu$ on $(S, \mathcal{B})$ such that, for every fixed $B \in \mathcal{B}$,

$$
\sigma_{n}\left(s_{1}, \ldots, s_{n}\right)(B) \longrightarrow \mu(s)(B) \quad \text { for } P_{\sigma} \text {-almost all } s \in S^{\infty}
$$

As noted in Section 2.1, the role played by $\mu$ is not as crucial as in the exchangeable case, as $P_{\sigma}$ is not completely determined by $\mu$.

Example 17. Take ( $X_{n}$ ) c.i.d. but not exchangeable under $P_{\sigma}$ and define

$$
Q(A)=E_{P_{\sigma}}\left\{\mu^{\infty}(A)\right\} \quad \text { for all } A \in \mathcal{B}^{\infty} .
$$

By definition, $\left(X_{n}\right)$ is exchangeable under $Q$. Also, since $S$ is nice, $Q=P_{\sigma^{*}}$ for some strategy $\sigma^{*}$. Thus, $P_{\sigma} \neq P_{\sigma^{*}}$. However, for every fixed $B \in \mathcal{B}$,

$$
\sigma_{n}^{*}\left(s_{1}, \ldots, s_{n}\right)(B) \longrightarrow \mu(s)(B) \quad \text { for } P_{\sigma^{*}-\text {-almost }} \text { all } s \in S^{\infty} .
$$

Despite Example 17, $\mu$ is an important random parameter for $\left(X_{n}\right)$ and a (natural) question is: What kind of random probability measures $\mu$ can be obtained if $\sigma \in \Sigma$ ? We address this question when $\sigma$ is given by (4). To this end, we first prove a general result.

In the next statement, we write "a.s." to mean " $P_{\sigma}$-a.s." and we denote by $X_{1}^{*}, X_{2}^{*}, \ldots$ the (finite or infinite) sequence of distinct observations corresponding to $X_{1}, X_{2}, \ldots$ More specifically, if $N$ is the cardinality of the (random) set $\left\{X_{1}, X_{2}, \ldots\right\}$, we let

$$
\begin{aligned}
X_{n}^{*}=X_{\tau_{n}} & \text { for all integers } n \text { such that } 1 \leq n \leq N, \\
& \text { where } \tau_{1}=1 \text { and } \tau_{n}=\inf \left\{j: X_{j} \notin\left\{X_{1}^{*}, \ldots, X_{n-1}^{*}\right\}\right\} .
\end{aligned}
$$

Theorem 18. Suppose $\left(X_{n}\right)$ is c.i.d. under $P_{\sigma}$, where $\sigma$ is any strategy. Then,

$$
\begin{equation*}
\mu \stackrel{\text { a.s. }}{=} \sum_{k} W_{k} \delta_{X_{k}^{*}}, \tag{8}
\end{equation*}
$$

for some random variables $W_{k} \geq 0$ such that $\sum_{k} W_{k}=1$, if and only if

$$
\begin{equation*}
\lim _{n} P_{\sigma}\left(X_{n} \neq X_{i} \text { for each } i<n\right)=0 . \tag{9}
\end{equation*}
$$

In addition,

$$
\begin{equation*}
W_{k} \stackrel{\text { a.s. }}{=} \lim _{n} \frac{1}{n} \sum_{i=1}^{n} 1_{\left\{X_{i}=X_{k}^{*}\right\}} . \tag{10}
\end{equation*}
$$

Proof. To make the notation easier, write $P=P_{\sigma}, E=E_{P_{\sigma}}$ and $I_{n-1}=\left(X_{1}, \ldots, X_{n-1}\right)$.
We first note a simple fact. Let

$$
\begin{aligned}
\gamma_{1} & =\delta_{I_{n-1}} \times \delta_{X_{n}}, \quad \gamma_{2}=\delta_{I_{n-1}} \times \mu, \quad \text { and } \\
H & =\left\{\left(s_{1}, \ldots, s_{n}\right) \in S^{n}: s_{n}=s_{i} \text { for some } i<n\right\} .
\end{aligned}
$$

Then, $\gamma_{1}$ and $\gamma_{2}$ are random probability measures on $\left(S^{n}, \mathcal{B}^{n}\right)$ such that

$$
\gamma_{1}(H)=\delta_{X_{n}}\left(\left\{X_{1}, \ldots, X_{n-1}\right\}\right) \quad \text { and } \quad \gamma_{2}(H)=\mu\left(\left\{X_{1}, \ldots, X_{n-1}\right\}\right) .
$$

Next, define two (non-random) probability measures on $\left(S^{n}, \mathcal{B}^{n}\right)$ as

$$
\gamma_{1}^{*}(C)=E\left\{\gamma_{1}(C)\right\} \quad \text { and } \quad \gamma_{2}^{*}(C)=E\left\{\gamma_{2}(C)\right\} \quad \text { for all } C \in \mathcal{B}^{n} .
$$

Since $\left(X_{n}\right)$ is c.i.d. under $P$, then $P\left(X_{n} \in B \mid I_{n-1}\right)=E\left(\mu(B) \mid I_{n-1}\right)$ a.s. for each $B \in \mathcal{B}$; see Section 2.1. Therefore,

$$
\begin{aligned}
\gamma_{1}^{*}(A \times B) & =P\left(I_{n-1} \in A, X_{n} \in B\right) \\
& =E\left\{1_{A}\left(I_{n-1}\right) P\left(X_{n} \in B \mid I_{n-1}\right)\right\} \\
& =E\left\{1_{A}\left(I_{n-1}\right) E\left(\mu(B) \mid I_{n-1}\right)\right\} \\
& =E\left\{1_{A}\left(I_{n-1}\right) \mu(B)\right\}=\gamma_{2}^{*}(A \times B)
\end{aligned}
$$

for all $A \in \mathcal{B}^{n-1}$ and $B \in \mathcal{B}$. Hence, $\gamma_{1}^{*}=\gamma_{2}^{*}$ on $\mathcal{B}^{n}$, which in turn implies

$$
\begin{aligned}
P\left(X_{n}=X_{i} \text { for some } i<n\right) & =E\left(\delta_{X_{n}}\left(\left\{X_{1}, \ldots, X_{n-1}\right\}\right)\right) \\
& =\gamma_{1}^{*}(H)=\gamma_{2}^{*}(H)=E\left(\mu\left(\left\{X_{1}, \ldots, X_{n-1}\right\}\right)\right)
\end{aligned}
$$

It follows that

$$
E\left(\mu\left(\left\{X_{1}^{*}, X_{2}^{*}, \ldots\right\}\right)\right)=\lim _{n} E\left(\mu\left(\left\{X_{1}, \ldots, X_{n-1}\right\}\right)\right)=\lim _{n} P\left(X_{n}=X_{i} \text { for some } i<n\right)
$$

This proves the equivalence between (8) and (9). In fact,

$$
\text { condition }(8) \Leftrightarrow \mu\left(\left\{X_{1}^{*}, X_{2}^{*}, \ldots\right\}\right) \stackrel{\text { a.s. }}{=} 1 \Leftrightarrow E\left(\mu\left(\left\{X_{1}^{*}, X_{2}^{*}, \ldots\right\}\right)\right)=1
$$

We finally turn to (10). As noted in Section 2.1, $\mu$ also satisfies

$$
\mu_{n}(B) \xrightarrow{\text { a.s. }} \mu(B) \quad \text { for every fixed } B \in \mathcal{B},
$$

where $\mu_{n}=\frac{1}{n} \sum_{i=1}^{n} \delta_{X_{i}}$ is the empirical measure. Hence,

$$
P\left(\mu_{n} \xrightarrow{\text { weakly }} \mu\right)=1 .
$$

If condition (9) holds, then

$$
\mu\left(\left\{X_{1}^{*}, X_{2}^{*}, \ldots\right\}\right) \stackrel{\text { a.s. }}{=} 1 \quad \text { and } \quad \mu_{n}\left(\left\{X_{1}^{*}, X_{2}^{*}, \ldots\right\}\right)=1 \quad \text { for each } n
$$

where the first equation has been proved above and the second is trivial. Hence, under (9), $\mu_{n}$ converges to $\mu$ in total variation norm with probability 1 , that is,

$$
\sup _{B \in \mathcal{B}}\left|\mu_{n}(B)-\mu(B)\right| \xrightarrow{\text { a.s. }} 0 .
$$

In particular,

$$
W_{k}=\mu\left(\left\{X_{k}^{*}\right\}\right)=\lim _{n} \mu_{n}\left(\left\{X_{k}^{*}\right\}\right)=\lim _{n} \frac{1}{n} \sum_{i=1}^{n} 1_{\left\{X_{i}=X_{k}^{*}\right\}} \quad \text { a.s. }
$$

Theorem 18 extends a result concerning exchangeability to the c.i.d. case. In fact, the equivalence between (8) and (9) is already known if $\left(X_{n}\right)$ is exchangeable under $P_{\sigma}$; see e.g. [23].

Finally, we focus on the special case where $\sigma$ is assessed according to (4). Then, Theorem 18 provides conditions for $\mu$ to be a.s. discrete.

Theorem 19. Suppose the strategy $\sigma$ is given by (4) and

$$
\prod_{i=0}^{n-1} q_{i}\left(X_{1}, \ldots, X_{i}\right) \xrightarrow{P_{\sigma}} 0 .
$$

Then, $\mu$ admits representation (8) and the weights $W_{k}$ are given by (10).
Proof. Just note that

$$
\begin{aligned}
P_{\sigma}\left(X_{n+1} \notin\left\{X_{1}, \ldots, X_{n}\right\} \mid\left(X_{1}, \ldots, X_{n}\right)=x\right) & =\sigma_{n}(x)\left(\left\{x_{1}, \ldots, x_{n}\right\}^{c}\right) \\
& =\sigma_{0}\left(\left\{x_{1}, \ldots, x_{n}\right\}^{c}\right) \prod_{i=0}^{n-1} q_{i}
\end{aligned}
$$

where $n \geq 1, x=\left(x_{1}, \ldots, x_{n}\right) \in S^{n}$ and $q_{i}=q_{i}\left(x_{1}, \ldots, x_{i}\right)$. Hence,

$$
\begin{aligned}
P_{\sigma}\left(X_{n+1} \neq X_{i} \text { for each } i \leq n\right) & =E_{P_{\sigma}}\left\{\sigma_{0}\left(\left\{X_{1}, \ldots, X_{n}\right\}^{c}\right) \prod_{i=0}^{n-1} q_{i}\left(X_{1}, \ldots, X_{i}\right)\right\} \\
& \leq E_{P_{\sigma}}\left\{\prod_{i=0}^{n-1} q_{i}\left(X_{1}, \ldots, X_{i}\right)\right\} \longrightarrow 0 .
\end{aligned}
$$

An application of Theorem 18 concludes the proof.
Various popular random probability measures $v$ admit the representation

$$
\begin{equation*}
v \stackrel{\text { a.s. }}{=} \sum_{k} D_{k} \delta_{Z_{k}}, \tag{11}
\end{equation*}
$$

where $\left(Z_{k}\right)$ is an i.i.d. sequence of random variables and the weights $\left(D_{k}\right)$ are independent of $\left(Z_{k}\right)$. Our last result is that $\mu$ often admits representation (11) provided $\sigma$ is given by (4) and the $q_{i}$ are constant.

Theorem 20. Suppose the strategy $\sigma$ is given by (4) and $\sigma_{0}$ is diffuse. Suppose also that $q_{i}$ is constant for every $i \geq 0$, and

$$
\prod_{i=0}^{n-1} q_{i} \rightarrow 0 \quad \text { and } \quad \sum_{n=1}^{\infty} \prod_{i=0}^{n-1} q_{i}=\infty
$$

Then, $\mu$ admits representation (8) and the weights $W_{k}$ are given by (10). Moreover, the sequence $\left(X_{k}^{*}\right)$ is i.i.d., $X_{1}^{*} \sim \sigma_{0}$, and $\left(X_{k}^{*}\right)$ is independent of $\left(W_{k}\right)$.

Proof. Take ( $T_{n}$ ) and ( $U_{i, j}$ ) satisfying conditions (j)-(jjj) and define $\left(Y_{n}\right)$ as in Section 7.1. Since the predictive distributions of $\left(Y_{n}\right)$ are given by (4), we can replace $\left(X_{n}\right)$ with $\left(Y_{n}\right)$. In addition, since

$$
\sum_{n} P\left(Y_{n+1} \notin\left\{Y_{1}, \ldots, Y_{n}\right\} \mid Y_{1}, \ldots, Y_{n}\right) \stackrel{\text { a.s. }}{=} \sum_{n} \prod_{i=0}^{n-1} q_{i}=\infty,
$$

the Borel-Cantelli lemma yields

$$
P\left(Y_{n+1} \notin\left\{Y_{1}, \ldots, Y_{n}\right\} \text { for infinitely many } n\right)=1
$$

Hence, one can define

$$
Y_{n}^{*}=Y_{\rho_{n}} \quad \text { for all } n \geq 1,
$$

where $\rho_{1}=1$ and $\rho_{n}=\inf \left\{j: Y_{j} \notin\left\{Y_{1}^{*}, \ldots, Y_{n-1}^{*}\right\}\right\}$.
Let $v$ be a random probability measure on $(S, \mathcal{B})$ such that

$$
P\left(Y_{n+1} \in B \mid Y_{1}, \ldots, Y_{n}\right) \xrightarrow{\text { a.s. }} v(B) \quad \text { for each fixed } B \in \mathcal{B} .
$$

Since $\prod_{i=0}^{n-1} q_{i} \rightarrow 0$, Theorem 19 implies

$$
v \stackrel{\text { a.s. }}{=} \sum_{k} D_{k} \delta_{Y_{k}^{*}} \quad \text { where } D_{k} \stackrel{\text { a.s. }}{=} \lim _{n} \frac{1}{n} \sum_{i=1}^{n} 1_{\left\{Y_{i}=Y_{k}^{*}\right\}} .
$$

We now prove that $\left(Y_{k}^{*}\right)$ is i.i.d., $Y_{1}^{*} \sim \sigma_{0}$, and $\left(Y_{k}^{*}\right)$ is independent of $\left(D_{k}\right)$.
Let $\mathcal{U}$ be the $\sigma$-field generated by $U_{i, j}$ for all $i$ and $j$ and

$$
A=\left\{T_{i} \neq T_{j} \text { for all } i \neq j\right\}
$$

On the set $A$, one obtains $Y_{n} \notin\left\{Y_{1}, \ldots, Y_{n-1}\right\}$ if and only if $Y_{n}=T_{n}$. Furthermore, $P(A)=1$ for $\left(T_{n}\right)$ is i.i.d. and $\sigma_{0}$ diffuse. Thus, up to a negligible set, $\rho_{k}$ is $\mathcal{U}$-measurable for each $k$. Similarly, up to a negligible set, $D_{k}$ is $\mathcal{U}$-measurable for each $k$. Since $\left(T_{k}\right)$ is independent of $\mathcal{U}$, it follows that $\left(T_{k}\right)$ is independent of $\left(D_{k}, \rho_{k}\right)$. Therefore, for each event $H$ in the $\sigma$-field generated by $\left(D_{k}\right)$, one obtains

$$
\begin{aligned}
P & \left(H \cap\left\{Y_{1}^{*} \in B_{1}, \ldots, Y_{k}^{*} \in B_{k}\right\}\right) \\
& =\sum_{m_{1}, \ldots, m_{k}} P\left(H \cap\left\{\rho_{1}=m_{1}, \ldots, \rho_{k}=m_{k}, T_{m_{1}} \in B_{1}, \ldots, T_{m_{k}} \in B_{k}\right\}\right) \\
& =\sum_{m_{1}, \ldots, m_{k}} P\left(T_{m_{1}} \in B_{1}, \ldots, T_{m_{k}} \in B_{k}\right) P\left(H \cap\left\{\rho_{1}=m_{1}, \ldots, \rho_{k}=m_{k}\right\}\right) \\
& =\prod_{i=1}^{k} \sigma_{0}\left(B_{i}\right) \sum_{m_{1}, \ldots, m_{k}} P\left(H \cap\left\{\rho_{1}=m_{1}, \ldots, \rho_{k}=m_{k}\right\}\right)=P(H) \prod_{i=1}^{k} \sigma_{0}\left(B_{i}\right) .
\end{aligned}
$$

This concludes the proof.
If $\sigma_{0}$ is diffuse, Theorem 20 applies to Dirichlet sequences of Example 8. In this special case, however, and more generally in case of exchangeable species sampling models, the conclusions of Theorem 20 are due to Pitman [27], Proposition 11.

## 8. Concluding remarks

In this paper, only a very few indications on how to select $f_{n}$ and $\alpha$ have been given. Are there some general criterions which help the inferrer in the choice of $f_{n}$ and $\alpha$ ? Which problems are appropriate
for $f_{n}(x, y, z)=q_{n}(x)$ and which require more elaborate choices of $f_{n}$ ? Or else, when is it reasonable to take $\alpha(y)=\delta_{y}$ ?

These questions are certainly crucial from a practical point of view, but to answer them is beyond the scope of this paper. However, we make three brief remarks.

First, in the subjective view of probability, any general criterion to select $\sigma$ is possibly useful but never mandatory. At least in principle, the choice of $\sigma$ depends only on the knowledge/feelings of the inferrer about the dependence structure of $\left(X_{n}\right)$.

Second, the kernel $\alpha$ controls how much the next outcome is affected by the observed data. This is quite evident in equation (1), where $f_{n}(x, y, z)=q_{n}(x)$ and the updating rule reduces to $\sigma_{n+1}(x, y)=$ $q_{n}(x) \sigma_{n}(x)+\left\{1-q_{n}(x)\right\} \alpha(y)$. In this case, the prediction of $X_{n+2}$ is affected by the last observation $X_{n+1}=y$ only through $\alpha(y)$. Thus, in a sense, the choice of $\alpha$ has to do with how much the last observation should be "reinforced". In this respect, $\alpha(y)=\delta_{y}$ is the more extreme choice of $\alpha$.

Third, as suggested by an anonymous referee, the inferrer could apply a sort of empirical Bayes procedure and select $\sigma \in \Sigma$ based on the available data. As a simple example, the strategy of Example 7 depends on a single parameter $q \in[0,1]$. Such a $q$ could be estimated by the data. For instance, a prior for $q$ could be assessed and $q$ could be estimated by the mean of the corresponding posterior.

## Acknowledgements

This paper has been improved by the helpful remarks of three anonymous referees.

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Received March 2020 and revised June 2020

