# $G$ method in action: Fast exact sampling from set of permutations of order $\boldsymbol{n}$ according to Mallows model through Cayley metric 

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#### Abstract

Using $G$ method, we give a fast exact (not approximate) Markovian method for sampling from $\mathbb{S}_{n}$, the set of permutations of order $n$, according to the Mallows model through Cayley metric (a model for ranked data). This method has something in common with the cyclic Gibbs sampler and something in common with the swapping method. The number of steps of our method is equal to the number of steps of swapping method, that is, $n-1$; moreover, both methods use the best probability distributions on sampling, the swapping method uses uniform probability distributions while our method uses almost uniform probability distributions (all the components of an almost uniform probability distribution are, here, identical, excepting at most one of them). But, besides sampling, we can do other things for the Mallows model through Cayley metric-we compute the normalizing constant and, by Uniqueness theorem, certain important probabilities.


## 1 The basic result we need

In this section, we present the basic result from Păun (2010) we need.
Set

$$
\operatorname{Par}(E)=\{\Delta \mid \Delta \text { is a partition of } E\},
$$

where $E$ is a nonempty set. We shall agree that the partitions do not contain the empty set.

Definition 1.1. Let $\Delta_{1}, \Delta_{2} \in \operatorname{Par}(E)$. We say that $\Delta_{1}$ is finer than $\Delta_{2}$ if $\forall V \in \Delta_{1}$, $\exists W \in \Delta_{2}$ such that $V \subseteq W$.

Write $\Delta_{1} \preceq \Delta_{2}$ when $\Delta_{1}$ is finer than $\Delta_{2}$.
In this article, a vector is a row vector and a stochastic matrix is a row stochastic matrix.

The entry $(i, j)$ of a matrix $Z$ will be denoted by $Z_{i j}$ or, if confusion can arise, $Z_{i \rightarrow j}$.

Key words and phrases. $G$ method, exact sampling, Gibbs sampler in a generalized sense, swapping method, Cayley metric, Mallows model, normalizing constant, important probabilities.

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Set

$$
\begin{aligned}
\langle m\rangle & =\{1,2, \ldots, m\} \quad(m \geq 1), \\
N_{m, n} & =\{P \mid P \text { is a nonnegative } m \times n \text { matrix }\}, \\
S_{m, n} & =\{P \mid P \text { is a stochastic } m \times n \text { matrix }\}, \\
N_{n} & =N_{n, n}, \\
S_{n} & =S_{n, n} .
\end{aligned}
$$

Let $P=\left(P_{i j}\right) \in N_{m, n}$. Let $\emptyset \neq U \subseteq\langle m\rangle$ and $\emptyset \neq V \subseteq\langle n\rangle$. Set the matrices
$P_{U}=\left(P_{i j}\right)_{i \in U, j \in\langle n\rangle}, \quad P^{V}=\left(P_{i j}\right)_{i \in\langle m\rangle, j \in V}, \quad$ and $\quad P_{U}^{V}=\left(P_{i j}\right)_{i \in U, j \in V}$.
Set

$$
\begin{aligned}
(\{i\})_{i \in\left\{s_{1}, s_{2}, \ldots, s_{t}\right\}} & =\left(\left\{s_{1}\right\},\left\{s_{2}\right\}, \ldots,\left\{s_{t}\right\}\right) ; \\
(\{i\})_{i \in\left\{s_{1}, s_{2}, \ldots, s_{t}\right\}} & \in \operatorname{Par}\left(\left\{s_{1}, s_{2}, \ldots, s_{t}\right\}\right) .
\end{aligned}
$$

Definition 1.2. Let $P \in N_{m, n}$. We say that $P$ is a generalized stochastic matrix if $\exists a \geq 0, \exists Q \in S_{m, n}$ such that $P=a Q$.

Definition 1.3 (Păun (2010)). Let $P \in N_{m, n}$. Let $\Delta \in \operatorname{Par}(\langle m\rangle)$ and $\Sigma \in \operatorname{Par}(\langle n\rangle)$. We say that $P$ is a $[\Delta]$-stable matrix on $\Sigma$ if $P_{K}^{L}$ is a generalized stochastic matrix, $\forall K \in \Delta, \forall L \in \Sigma$. In particular, a [ $\Delta$ ]-stable matrix on $(\{i\})_{i \in\langle n\rangle}$ is called [ $\Delta$ ]stable for short.

Definition 1.4 (Păun (2010)). Let $P \in N_{m, n}$. Let $\Delta \in \operatorname{Par}(\langle m\rangle)$ and $\Sigma \in \operatorname{Par}(\langle n\rangle)$. We say that $P$ is a $\Delta$-stable matrix on $\Sigma$ if $\Delta$ is the least fine partition for which $P$ is a [ $\Delta]$-stable matrix on $\Sigma$. In particular, a $\Delta$-stable matrix on $(\{i\})_{i \in\langle n\rangle}$ is called $\Delta$-stable while a $(\langle m\rangle)$-stable matrix on $\Sigma$ is called stable on $\Sigma$ for short. A stable matrix on $(\{i\})_{i \in\langle n\rangle}$ is called stable for short.

Let $\Delta_{1} \in \operatorname{Par}(\langle m\rangle)$ and $\Delta_{2} \in \operatorname{Par}(\langle n\rangle)$. Set (see Păun (2010) for $G_{\Delta_{1}, \Delta_{2}}$ and Păun (2011) for $\bar{G}_{\Delta_{1}, \Delta_{2}}$ )

$$
G_{\Delta_{1}, \Delta_{2}}=\left\{P \mid P \in S_{m, n} \text { and } P \text { is a }\left[\Delta_{1}\right] \text {-stable matrix on } \Delta_{2}\right\}
$$

and

$$
\bar{G}_{\Delta_{1}, \Delta_{2}}=\left\{P \mid P \in N_{m, n} \text { and } P \text { is a }\left[\Delta_{1}\right] \text {-stable matrix on } \Delta_{2}\right\} .
$$

When we study or even when we construct products of nonnegative matrices (in particular, products of stochastic matrices) using $G_{\Delta_{1}, \Delta_{2}}$ or $\bar{G}_{\Delta_{1}, \Delta_{2}}$ we shall refer this as the $G$ method.

Below, we give the basic result from Păun (2010) we need.

Theorem 1.5 (Păun (2010)). Let $P_{1} \in G_{\left(\left\langle m_{1}\right\rangle\right), \Delta_{2}} \subseteq S_{m_{1}, m_{2}}, P_{2} \in G_{\Delta_{2}, \Delta_{3}} \subseteq$


$$
P_{1} P_{2} \cdots P_{n}
$$

is a stable matrix (i.e., a matrix with identical rows, see Definition 1.4).
Proof. See Păun (2010). (Theorem 1.5 is part of Theorem 2.10 from Păun (2010); a generalization of Theorem 2.10 from Păun (2010) is Theorem 1.6 from Păun (2011).)

## 2 The Markovian method

In this section, we present the Mallows model and our fast Markovian method for sampling exactly (not approximately) from $\mathbb{S}_{n}$, the set of permutations of order $n$, according to the Mallows model through Cayley metric. In addition to sampling, for this special Mallows model, we compute the normalizing constant and, by Uniqueness theorem, certain important probabilities.

Consider the group $\left(\mathbb{S}_{n}, \circ\right)$, where $\circ$ is the usual composition of functions. $(u, v)$ is a transposition, $\forall u, v \in\langle n\rangle, u \neq v$. Set $(u, u)=\mathrm{Id}, \forall u \in\langle n\rangle$, where Id is the identity permutation.

Theorem 2.1. Let $n \geq 2$. Let $\sigma_{0} \in \mathbb{S}_{n}$. Let

$$
\begin{aligned}
\mathbb{M}_{n, l}= & \left\{\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l} \mid i_{1}, i_{2}, \ldots, i_{l} \in\langle n\rangle, 1 \leq i_{1} \leq n,\right. \\
& \left.2 \leq i_{2} \leq n, \ldots, l \leq i_{l} \leq n, \sigma_{l} \in \mathbb{S}_{n}, \sigma_{l}(v)=v, \forall v \in\langle l\rangle\right\} \quad \forall l \in\langle n-1\rangle .
\end{aligned}
$$

Then

$$
\mathbb{M}_{n, l}=\mathbb{S}_{n} \quad \forall l \in\langle n-1\rangle
$$

Proof. Let $l \in\langle n-1\rangle$. Since $\left(\mathbb{S}_{n}, o\right)$ is a group, we have $\mathbb{M}_{n, l} \subseteq \mathbb{S}_{n}$. Therefore, $\left|\mathbb{M}_{n, l}\right| \leq\left|\mathbb{S}_{n}\right|=n!(|\cdot|$ is the cardinal). To finish the proof, we show that $\left|\mathbb{M}_{n, l}\right|=n!$.

The number of permutations $\sigma_{l} \in \mathbb{S}_{n}$ with $\sigma_{l}(v)=v, \forall v \in\langle l\rangle$, is equal to $(n-l)!$. Since $1 \leq i_{1} \leq n, 2 \leq i_{2} \leq n, \ldots, l \leq i_{l} \leq n$, it follows that $\left|\mathbb{M}_{n, l}\right|$ is at most equal to

$$
n(n-1) \cdots(n-l+1)[(n-l)!]=n!.
$$

We show that

$$
\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l}=\sigma_{0} \circ\left(1, j_{1}\right) \circ\left(2, j_{2}\right) \circ \cdots \circ\left(l, j_{l}\right) \circ \tau_{l}
$$

if and only if $i_{k}=j_{k}, \forall k \in\langle l\rangle$, and $\sigma_{l}=\tau_{l}$, where $i_{1}, j_{1}, i_{2}, j_{2}, \ldots, i_{l}, j_{l} \in\langle n\rangle$, $1 \leq i_{1}, j_{1} \leq n, 2 \leq i_{2}, j_{2} \leq n, \ldots, l \leq i_{l}, j_{l} \leq n, \sigma_{l}, \tau_{l} \in \mathbb{S}_{n}, \sigma_{l}(v)=\tau_{l}(v)=v$, $\forall v \in\langle l\rangle$.
" $\Longleftarrow "$ Obvious.
" $\Longrightarrow " \sigma_{0}$ can be removed. We remove it, so, we suppose that

$$
\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l}=\left(1, j_{1}\right) \circ\left(2, j_{2}\right) \circ \cdots \circ\left(l, j_{l}\right) \circ \tau_{l} .
$$

It follows that

$$
\left[\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l}\right](1)=\left[\left(1, j_{1}\right) \circ\left(2, j_{2}\right) \circ \cdots \circ\left(l, j_{l}\right) \circ \tau_{l}\right](1)
$$

Therefore,

$$
i_{1}=j_{1}
$$

Since $i_{1}=j_{1}$, removing $\left(1, i_{1}\right)$ and $\left(1, j_{1}\right)$, we have

$$
\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l}=\left(2, j_{2}\right) \circ \cdots \circ\left(l, j_{l}\right) \circ \tau_{l} .
$$

It follows that

$$
\left[\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l}\right](2)=\left[\left(2, j_{2}\right) \circ \cdots \circ\left(l, j_{l}\right) \circ \tau_{l}\right](2)
$$

Therefore,

$$
i_{2}=j_{2}
$$

Proceeding in this way, we obtain

$$
i_{1}=j_{1}, \quad i_{2}=j_{2}, \quad \ldots, \quad i_{l}=j_{l}
$$

and, as a result of these equations,

$$
\sigma_{l}=\tau_{l}
$$

We conclude that

$$
\left|\mathbb{M}_{n, l}\right|=n!
$$

Theorem 2.1 says that we can work with $\mathbb{M}_{n, l}$ instead of $\mathbb{S}_{n}, \forall l \in\langle n-1\rangle$ (this fact will be used in Theorem 2.3).

Let $C(\sigma, \tau)=$ minimum number of transpositions required to bring $\sigma$ to $\tau$, $\forall \sigma, \tau \in \mathbb{S}_{n} . C$ is a metric on $\mathbb{S}_{n}$, called the Cayley metric (see, e.g., Diaconis and Saloff-Coste (1998)).

Theorem 2.2. Let $n \geq 2$. Let $\sigma_{0} \in \mathbb{S}_{n}$. Consider on $\mathbb{S}_{n}$ the Cayley metric. Then

$$
\begin{aligned}
C & \left(\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l}, \sigma_{0}\right) \\
& =\left\{\begin{array}{l}
C\left(\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, k) \circ \sigma_{l}, \sigma_{0}\right) \\
\text { if } j=k=l \text { or } j, k>l, \\
C\left(\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, k) \circ \sigma_{l}, \sigma_{0}\right)-1 \\
\text { if } j=l, k>l, \\
C\left(\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, k) \circ \sigma_{l}, \sigma_{0}\right)+1 \\
\text { if } j>l, k=l,
\end{array}\right.
\end{aligned}
$$

$\forall l \in\langle n-1\rangle, \forall i_{1}, i_{2}, \ldots, i_{l-1}, j, k \in\langle n\rangle, 1 \leq i_{1} \leq n, 2 \leq i_{2} \leq n, \ldots, l-1 \leq i_{l-1} \leq$ $n, l \leq j, k \leq n, \forall \sigma_{l} \in \mathbb{S}_{n}, \sigma_{l}(v)=v, \forall v \in\langle l\rangle\left(\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right)\right.$, etc. vanish when $l=1$ ).

Proof. Case 1. $j=k=l$ or $j, k>l$.
Subcase 1.1. $j=k=l$. Obvious $((l, j)=(l, k)=\mathrm{Id})$.
Subcase 1.2. $j, k>l$. It is known that $(u, v) \circ(u, v)=\mathrm{Id}, \forall(u, v),(u, v)$ is a transposition, $\psi_{1} \circ \psi_{2}=\psi_{2} \circ \psi_{1}, \forall \psi_{1}, \psi_{2} \in \mathbb{S}_{n}, \psi_{1}$ and $\psi_{2}$ are disjoint cycles, any permutation can be factored uniquely (leaving the order of factors aside) into a product of pair-wise disjoint cycles ("factored", "factors", and "product" are improper words), and any cycle of length $s(s \geq 2)$ can be factored into a product of $s-1$ transpositions in $s$ different ways. Since $\sigma_{l}(v)=v, \forall v \in\langle l\rangle$, it follows that, for any cycle of $\sigma_{l}$, any factorization of the cycle into a product of $s-1$ transpositions, where $s$ is the length of cycle ("factorization" is an improper word), does not contain the transpositions $\left(1, i_{1}\right),\left(2, i_{2}\right), \ldots,\left(l-1, i_{l-1}\right),(l, j)$, and $(l, k)$. So,

$$
\begin{aligned}
& \left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l} \\
& \left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, k) \circ \sigma_{l}
\end{aligned}
$$

cannot be simplified. The equation we need prove is now obvious.
Case 2. $j=l, k>l$. Since $j=l$, we have $(l, j)=$ Id. Proceeding similar to Subcase 1.2,

$$
\begin{aligned}
& \left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ \sigma_{l}, \\
& \left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, k) \circ \sigma_{l}
\end{aligned}
$$

cannot be simplified. Therefore, the equation we need prove holds.
Case 3. $j>l, k=l$. Similar to Case 2.
Recall that $\mathbb{R}^{+}=\{x \mid x \in \mathbb{R}$ and $x>0\}$.
Let

$$
\pi_{\sigma}=\frac{\theta^{d\left(\sigma, \sigma_{0}\right)}}{Z} \quad \forall \sigma \in \mathbb{S}_{n},
$$

where $\theta \in \mathbb{R}^{+}$(cases of interest: $\left.0<\theta \leq 1 ; \theta>1\right), \sigma_{0} \in \mathbb{S}_{n}(n \geq 1), d$ is a metric on $\mathbb{S}_{n}$, and

$$
Z=\sum_{\sigma \in \mathbb{S}_{n}} \theta^{d\left(\sigma, \sigma_{0}\right)}
$$

The probability distribution $\pi=\left(\pi_{\sigma}\right)_{\sigma \in \mathbb{S}_{n}}$ (on $\mathbb{S}_{n}$ ) is called the Mallows model through metric $d$ (see Mallows (1957); see, e.g., also Critchlow (1985), Diaconis (2009), Diaconis and Saloff-Coste (1998), Fligner and Verducci (1993), and Marden (1995)). This is a model-an exponential model when $\theta \neq 1$-for ranked data (see the above references).

In this article, the transpose of a vector $x$ is denoted by $x^{\prime}$. Set $e=e(n)=$ $(1,1, \ldots, 1) \in \mathbb{R}^{n}, \forall n \geq 1$.

Below, we give the main result of this work.
Theorem 2.3. Let $n \geq 2$. Let $\pi=\left(\pi_{\sigma}\right)_{\sigma \in \mathbb{S}_{n}}$ be the Mallows model through Cayley metric. Consider a Markov chain with state space $\mathbb{S}_{n}$, initial probability distribution $p_{0}$, and transition matrix $P=P_{1} P_{2} \cdots P_{n-1}$, where $P_{l}, l \in\langle n-1\rangle$, are stochastic matrices on $\mathbb{S}_{n}$,

$$
\begin{aligned}
& \left(P_{l}\right)_{\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l} \rightarrow \xi} \\
& =\left\{\begin{array}{l}
\frac{\pi_{\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l}}}{\sum_{l \leq k \leq n} \pi_{\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, k) \circ \sigma_{l}}} \\
\quad \text { if } \xi=\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l} \\
\quad \text { for some } j, l \leq j \leq n, \\
\quad \text { if } \xi \neq \sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l}, \\
\quad \forall j, l \leq j \leq n,
\end{array}\right.
\end{aligned}
$$

$\forall l \in\langle n-1\rangle-\left(\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right)\right.$ vanishes when $\left.l=1\right), \forall i_{1}, i_{2}, \ldots$, $i_{l} \in\langle n\rangle, 1 \leq i_{1} \leq n, 2 \leq i_{2} \leq n, \ldots, l \leq i_{l} \leq n, \forall \sigma_{l} \in \mathbb{S}_{n}, \sigma_{l}(v)=v, \forall v \in\langle l\rangle$, $\forall \xi \in \mathbb{S}_{n}$, where $\sigma_{0}$ is the parameter from $\mathbb{S}_{n}$ of Mallows model through Cayley metric. Then

$$
P=e^{\prime} \pi
$$

(therefore, the chain attains its stationarity at time 1, its stationary (limit) probability distribution being, obviously, $\pi$ ).

Proof. Set
$K_{\left(i_{1}, i_{2}, \ldots, i_{l}\right)}=\left\{\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l} \mid \sigma_{l} \in \mathbb{S}_{n}, \sigma_{l}(v)=v, \forall v \in\langle l\rangle\right\}$, $\forall l \in\langle n-1\rangle, \forall i_{1}, i_{2}, \ldots, i_{l} \in\langle n\rangle, 1 \leq i_{1} \leq n, 2 \leq i_{2} \leq n, \ldots, l \leq i_{l} \leq n$. We have

$$
\bigcup_{\substack{i_{1}, i_{2}, \ldots, i_{l} \in\langle n\rangle \\ 1 \leq i_{1} \leq n \\ 2 \leq i_{2} \leq n}} K_{\left(i_{1}, i_{2}, \ldots, i_{l}\right)}=\mathbb{M}_{n, l}=\mathbb{S}_{n} \quad \forall l \in\langle n-1\rangle
$$

(see Theorem 2.1). We show that

$$
K_{\left(i_{1}, i_{2}, \ldots, i_{l}\right)} \cap K_{\left(j_{1}, j_{2}, \ldots, j_{l}\right)}=\emptyset
$$

if $\exists u \in\langle l\rangle$ such that $i_{u} \neq j_{u}$, where $l \in\langle n-1\rangle, i_{1}, j_{1}, i_{2}, j_{2}, \ldots, i_{l}, j_{l} \in\langle n\rangle, 1 \leq i_{1}$, $j_{1} \leq n, 2 \leq i_{2}, j_{2} \leq n, \ldots, l \leq i_{l}, j_{l} \leq n$. To see this, we suppose that $\exists u \in\langle l\rangle$ with $i_{u} \neq j_{u}$ such that

$$
K_{\left(i_{1}, i_{2}, \ldots, i_{l}\right)} \cap K_{\left(j_{1}, j_{2}, \ldots, j_{l}\right)} \neq \emptyset
$$

Let $\omega \in K_{\left(i_{1}, i_{2}, \ldots, i_{l}\right)} \cap K_{\left(j_{1}, j_{2}, \ldots, j_{l}\right)}$. We have

$$
\begin{aligned}
\omega & =\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l} \\
& =\sigma_{0} \circ\left(1, j_{1}\right) \circ\left(2, j_{2}\right) \circ \cdots \circ\left(l, j_{l}\right) \circ \tau_{l}
\end{aligned}
$$

where $\sigma_{l}, \tau_{l} \in \mathbb{S}_{n}, \sigma_{l}(v)=\tau_{l}(v)=v, \forall v \in\langle l\rangle$. Proceeding as in the proof of Theorem 2.1 ( $\sigma_{0}$ is removed, $\ldots$ ), we obtain

$$
i_{1}=j_{1}, \quad i_{2}=j_{2}, \quad \ldots, \quad i_{l}=j_{l}, \quad \sigma_{l}=\tau_{l}
$$

Therefore, we obtained a contradiction.
The above results lead to the fact that

is a partition of $\mathbb{M}_{n, l}\left(\mathbb{M}_{n, l}=\mathbb{S}_{n}\right), \forall l \in\langle n-1\rangle$. Set the partitions (this can now be done)

$$
\begin{array}{cl}
\Delta_{1}=\left(\mathbb{S}_{n}\right), \\
\Delta_{l+1}=\left(K_{\left(i_{1}, i_{2}, \ldots, i_{l}\right)}\right)_{i_{1}, i_{2}, \ldots, i_{l} \in\langle n\rangle} 1 \leq i_{1} \leq n \\
2 \leq i_{2} \leq n \\
& \vdots \\
& l \leq i_{l} \leq n
\end{array}
$$

$\forall l \in\langle n-1\rangle$. Obviously, we have $\Delta_{n}=(\{\sigma\})_{\sigma \in \mathbb{S}_{n}}$.
By hypothesis and Theorem 2.2, we have

$$
\begin{aligned}
& \left(P_{l}\right)_{\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l, i_{l}\right) \circ \sigma_{l} \rightarrow \xi} \\
& =\left\{\begin{array}{l}
\frac{\theta^{C\left(\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l}, \sigma_{0}\right)}}{\sum_{l \leq k \leq n} \theta^{C\left(\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, k) \circ \sigma_{l}, \sigma_{0}\right)}} \\
\quad \text { if } \xi=\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l} \\
\quad \text { for some } j, l \leq j \leq n, \\
\quad \text { if } \xi \neq \sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l}, \\
\quad \forall j, l \leq j \leq n,
\end{array}\right.
\end{aligned}
$$

$$
=\left\{\begin{array}{l}
\frac{1}{1+(n-l) \theta} \\
\text { if } \xi=\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, l) \circ \sigma_{l}, \\
\frac{\theta}{1+(n-l) \theta} \\
\quad \begin{array}{l}
\text { if } \xi=\sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l} \\
\quad \text { for some } j, l<j \leq n, \\
\text { if } \xi \neq \sigma_{0} \circ\left(1, i_{1}\right) \circ\left(2, i_{2}\right) \circ \cdots \circ\left(l-1, i_{l-1}\right) \circ(l, j) \circ \sigma_{l}, \\
\forall j, l \leq j \leq n,
\end{array}
\end{array}\right.
$$

$\forall l \in\langle n-1\rangle, \forall i_{1}, i_{2}, \ldots, i_{l} \in\langle n\rangle, 1 \leq i_{1} \leq n, 2 \leq i_{2} \leq n, \ldots, l \leq i_{l} \leq n, \forall \sigma_{l} \in \mathbb{S}_{n}$, $\sigma_{l}(v)=v, \forall v \in\langle l\rangle, \forall \xi \in \mathbb{S}_{n}$. It follows that

$$
P_{l} \in G_{\Delta_{l}, \Delta_{l+1}} \quad \forall l \in\langle n-1\rangle .
$$

Since $P=P_{1} P_{2} \cdots P_{n-1}$, by Theorem 1.5, $P$ is a stable matrix. Consequently, $\exists \psi, \psi$ is a probability distribution on $\mathbb{S}_{n}$, such that

$$
P=e^{\prime} \psi
$$

It is easy to see that

$$
\pi_{\sigma}\left(P_{l}\right)_{\sigma \tau}=\pi_{\tau}\left(P_{l}\right)_{\tau \sigma} \quad \forall l \in\langle n-1\rangle, \forall \sigma, \tau \in \mathbb{S}_{n}
$$

$\left(\mathbb{S}_{n}=\mathbb{M}_{n, l}, \forall l \in\langle n-1\rangle\right)$. This thing implies

$$
\pi P_{l}=\pi \quad \forall l \in\langle n-1\rangle,
$$

and, further,

$$
\pi P=\pi
$$

Finally, we have

$$
\pi=\pi P=\pi e^{\prime} \psi=\psi
$$

so,

$$
P=e^{\prime} \pi
$$

We comment on Theorem 2.3 and its proof.
First, we can work with the chain with transition matrix $P$ or, equivalently, with the chain with transition matrices $P_{1}, P_{2}, \ldots, P_{n-1}, P_{1}, P_{2}, \ldots, P_{n-1}, \ldots$ (the former chain is homogeneous while the latter one is nonhomogeneous when $n \geq 3$ ). We chose the first case. (For finite Markov chain theory, see, e.g., Iosifescu (2007).) Any 1 -step of the chain with transition matrix $P$ is performed via $P_{1}, P_{2}, \ldots, P_{n-1}$, that is, doing $n-1$ transitions: one using $P_{1}$, one using $P_{2}, \ldots$, one using $P_{n-1}$. By Theorem 2.3, the chain with transition matrix $P$ attains its stationarity at time 1 (to attain the stationarity, the chain with transition matrix $P$ makes one step while the other chain makes $n-1$ steps (due to $P_{1}, P_{2}, \ldots, P_{n-1}$ )).

We constructed the chain with transition matrix $P$ being guided by $G$ method, Theorem 1.5, our hybrid Metropolis-Hastings chain(s) from Păun (2011), and, especially, certain results and suggestions from Păun (2017). The chain with transition matrix $P$ belongs to our collection of hybrid Metropolis-Hastings chains from Păun (2011) (this follows from $K_{\left(i_{1}, i_{2}, \ldots, i_{l+1}\right)} \subset K_{\left(i_{1}, i_{2}, \ldots, i_{l}\right)}, \forall l \in\langle n-2\rangle$, $\forall i_{1}, i_{2}, \ldots, i_{l+1}, 1 \leq i_{1} \leq n, 2 \leq i_{2} \leq n, \ldots, l+1 \leq i_{l+1} \leq n$, etc.; for our collection, see also Păun (2017)). The chain with transition matrix $P$ is a cyclic Gibbs sampler in a generalized sense because the state space is, here, $\mathbb{S}_{n}$, ratios used to define the transition probabilities of matrices $P_{l}, l \in\langle n-1\rangle$, are similar to those of (usual) cyclic Gibbs sampler with-keeping the finite framework-finite state space (this chain also belongs to our collection of hybrid Metropolis-Hastings chains from Păun (2011), see Păun (2017)), and matrices $P_{1}, P_{2}, \ldots, P_{n-1}$ are used cyclically. (For the Gibbs sampler, see, e.g., Madras (2002).)

Second, to define transition probabilities of $P_{l}, l \in\langle n-1\rangle$ fixed, we used states from $\mathbb{M}_{n, l}$. So, using $P_{l}$, the chain passes from a state, say, $\gamma$ of $\mathbb{M}_{n, l}$ to a state, say, $\delta$ of $\mathbb{M}_{n, l}$ also. For $P_{l+1}$, when $l+1 \leq n-1$, we need states from $\mathbb{M}_{n, l+1}$, so, when we run the chain, we must rewrite $\delta$ using the generators of $\mathbb{M}_{n, l+1}$.

Third, there exists a case, a happy case, for which rewriting the states from Second is not necessary, namely, when $\sigma_{l}=$ Id. So, to avoid rewriting the states, we consider the chain with initial probability distribution $p_{0}$ with $\left(p_{0}\right)_{\sigma_{0}}=1$ (warning! here we have $\sigma_{0}$ and above we have $\sigma_{l}$ ). Since $P=e^{\prime} \pi$, we have

$$
p_{0} P^{m}=\pi \quad \forall m \geq 1, \forall p_{0}, p_{0}=\text { initial probability distribution. }
$$

So, for the initial probability distribution $p_{0}$ with $\left(p_{0}\right)_{\sigma_{0}}=1$, the above equations hold as well. From $\sigma_{0}=\sigma_{0} \circ(1,1) \circ \operatorname{Id} \in \mathbb{M}_{n, 1}\left(\sigma_{1}=\mathrm{Id}\right)$, the chain passes in one of the states

$$
\begin{aligned}
\sigma_{0} & =\sigma_{0} \circ(1,1)=\sigma_{0} \circ(1,1) \circ \mathrm{Id} \in \mathbb{M}_{n, 1}, \\
\sigma_{0} \circ(1,2) & =\sigma_{0} \circ(1,2) \circ \mathrm{Id} \in \mathbb{M}_{n, 1}, \\
& \vdots \\
\sigma_{0} \circ(1, n) & =\sigma_{0} \circ(1, n) \circ \mathrm{Id} \in \mathbb{M}_{n, 1} .
\end{aligned}
$$

Suppose that it passed in the state $\sigma_{0} \circ(1,2)$. From $\sigma_{0} \circ(1,2)=\sigma_{0} \circ(1,2) \circ(2,2) \circ$ $\operatorname{Id} \in \mathbb{M}_{n, 2}\left(\sigma_{2}=\mathrm{Id}\right)$, the chain passes in one of the states

$$
\begin{aligned}
\sigma_{0} \circ(1,2) & =\sigma_{0} \circ(1,2) \circ(2,2)=\sigma_{0} \circ(1,2) \circ(2,2) \circ \operatorname{Id} \in \mathbb{M}_{n, 2}, \\
\sigma_{0} \circ(1,2) \circ(2,3) & =\sigma_{0} \circ(1,2) \circ(2,3) \circ \mathrm{Id} \in \mathbb{M}_{n, 2},
\end{aligned}
$$

$$
\sigma_{0} \circ(1,2) \circ(2, n)=\sigma_{0} \circ(1,2) \circ(2, n) \circ \mathrm{Id} \in \mathbb{M}_{n, 2}
$$

Suppose that it passed in the state $\sigma_{0} \circ(1,2) \circ(2, n-1)$, etc. Therefore, the states are generated proceeding similar to the swapping method, the difference being that, here, we use the probability distributions

$$
\left(\frac{1}{1+(n-l) \theta}, \frac{\theta}{1+(n-l) \theta}, \frac{\theta}{1+(n-l) \theta}, \ldots, \frac{\theta}{1+(n-l) \theta}\right), \quad l \in\langle n-1\rangle
$$

instead of uniform probability distributions. (For the swapping method, see, e.g., Devroye (1986), pp. 645-646.) The above probability distributions, the former being almost uniform probability distributions-we call them almost uniform probability distributions because each of these probability distributions has identical components, excepting at most one of them (all the components are identical when $\theta=1$ )—and the latter, those of swapping method, being uniform probability distributions, are, concerning the implementation, the best ones. To see that this is also true for almost uniform probability distributions, we split each almost uniform probability distribution into two blocks,

$$
\left(\frac{1}{1+(n-l) \theta}\right), \quad\left(\frac{\theta}{1+(n-l) \theta}, \frac{\theta}{1+(n-l) \theta}, \ldots, \frac{\theta}{1+(n-l) \theta}\right)
$$

If

$$
X>\frac{1}{1+(n-l) \theta}, \quad X \sim U(0,1)
$$

further, we work with the latter block, which, by normalization, leads to the uniform probability distribution

$$
\left(\frac{1}{n-l}, \frac{1}{n-l}, \ldots, \frac{1}{n-l}\right)
$$

Therefore, our exact sampling Markovian method, having $n-1$ steps, is simple and good.

Fourth, using the equation $P=e^{\prime} \pi$, we can compute the normalizing constant $Z$. Set $\Delta \succ \Delta^{\prime}$ if $\Delta^{\prime} \preceq \Delta$ and $\Delta \neq \Delta^{\prime}$, where $\Delta, \Delta^{\prime} \in \operatorname{Par}(E), E$ is a nonempty set. The partitions $\Delta_{1}=\left(\mathbb{S}_{n}\right), \Delta_{2}, \ldots, \Delta_{n-1}, \Delta_{n}=(\{\sigma\})_{\sigma \in \mathbb{S}_{n}}$ from the proof of Theorem 2.3 have the property: $\Delta_{1} \succ \Delta_{2} \succ \cdots \succ \Delta_{n}$ (recall that $K_{\left(i_{1}, i_{2}, \ldots, i_{l+1}\right)} \subset K_{\left(i_{1}, i_{2}, \ldots, i_{l}\right)}, \forall l \in\langle n-2\rangle, \forall i_{1}, i_{2}, \ldots, i_{l+1}, 1 \leq i_{1} \leq n, 2 \leq i_{2} \leq$ $n, \ldots, l+1 \leq i_{l+1} \leq n$ ). $P_{l}$ is a block diagonal matrix (eventually by permutation of rows and columns), $\forall l \in\langle n-1\rangle-\{1\}$, and $\Delta_{l}$-stable matrix on $\Delta_{l}, \forall l \in\langle n-1\rangle$ (see First again - the fact that the chain with transition matrix $P$ belongs to our collection of hybrid Metropolis-Hastings chains from Păun (2011)). Moreover, $P_{l}$ is a $\Delta_{l}$-stable matrix on $\Delta_{l+1}, \forall l \in\langle n-1\rangle$. Due to these facts, using $P=e^{\prime} \pi$, it is easy to compute $\pi_{\sigma_{0}}$ (hint: use, directly, $\mathbb{S}_{n} \supset K_{(1)} \supset K_{(1,2)} \supset \cdots \supset K_{(1,2, \ldots, n-1)}=\left\{\sigma_{0}\right\}$ or, indirectly, the operator $(\cdot)^{-+}$from Păun (2010)); we have

$$
\pi_{\sigma_{0}}=\frac{1}{1+(n-1) \theta} \cdot \frac{1}{1+(n-2) \theta} \cdots \cdots \frac{1}{1+\theta} .
$$

Since, on the other hand,

$$
\pi_{\sigma_{0}}=\frac{\theta^{0}}{Z}=\frac{1}{Z}
$$

we obtain

$$
Z=(1+\theta)(1+2 \theta) \cdots(1+(n-1) \theta)
$$

This result is known (see, e.g., Diaconis (2009)), but our computation method is new, simple, and interesting.

Fifth, using Uniqueness theorem from Păun (2017) (the presentation of this result is too long, so, we omit to give it here), we can compute certain important probabilities of the Mallows model through Cayley metric. Indeed, by Uniqueness theorem we have

$$
P\left(K_{\left(i_{1}\right)}\right)=\sum_{\sigma \in K_{\left(i_{1}\right)}} \pi_{\sigma}= \begin{cases}\frac{1}{1+(n-1) \theta}, & \text { if } i_{1}=1 \\ \frac{\theta}{1+(n-1) \theta}, & \text { if } 1<i_{1} \leq n\end{cases}
$$

Note that

$$
K_{\left(i_{1}\right)}=\left\{\sigma \mid \sigma \in \mathbb{S}_{n}, \sigma(1)=\sigma_{0}\left(i_{1}\right)\right\} \quad \forall i_{1}, 1 \leq i_{1} \leq n
$$

( $K_{\left(i_{1}\right)}$ is the set of permutations from $\mathbb{S}_{n}$, each permutation having the first component equal to $\sigma_{0}\left(i_{1}\right)$, the $i_{1}$ th component of $\left.\sigma_{0}\right)$. Further, by Uniqueness theorem we have

$$
\frac{P\left(K_{\left(i_{1}, i_{2}\right)}\right)}{P\left(K_{\left(i_{1}\right)}\right)}=\frac{\sum_{\sigma \in K_{\left(i_{1}, i_{2}\right)}} \pi_{\sigma}}{\sum_{\sigma \in K_{\left(i_{1}\right)}} \pi_{\sigma}}= \begin{cases}\frac{1}{1+(n-2) \theta}, & \text { if } i_{2}=2 \\ \frac{\theta}{1+(n-2) \theta}, & \text { if } 2<i_{2} \leq n\end{cases}
$$

$\forall i_{1}, 1 \leq i_{1} \leq n$, so,

$$
P\left(K_{\left(i_{1}, i_{2}\right)}\right)= \begin{cases}\frac{1}{[1+(n-1) \theta][1+(n-2) \theta]}, & \text { if } i_{1}=1, i_{2}=2 \\ \frac{\theta}{[1+(n-1) \theta][1+(n-2) \theta]}, & \text { if } i_{1}=1,2<i_{2} \leq n \\ \frac{\theta^{2}}{[1+(n-1) \theta][1+(n-2) \theta]}, & \text { if } 1<i_{1} \leq n, 2<i_{2} \leq n\end{cases}
$$

Note that

$$
K_{\left(i_{1}, i_{2}\right)}= \begin{cases}\left\{\sigma \mid \sigma \in \mathbb{S}_{n}, \sigma(1)=\sigma_{0}\left(i_{1}\right), \sigma(2)=\sigma_{0}(1)\right\}, & \text { if } i_{2}=i_{1} \\ \left\{\sigma \mid \sigma \in \mathbb{S}_{n}, \sigma(1)=\sigma_{0}\left(i_{1}\right), \sigma(2)=\sigma_{0}\left(i_{2}\right)\right\}, & \text { if } i_{2} \neq i_{1}\end{cases}
$$

$\forall i_{1}, 1 \leq i_{1} \leq n, \forall i_{2}, 2 \leq i_{2} \leq n$. To compute $P\left(K_{\left(i_{1}, i_{2}, i_{3}\right)}\right)$, etc., we use (see Uniqueness theorem)

$$
\frac{P\left(K_{\left(i_{1}, i_{2}, \ldots, i_{u}\right)}\right)}{P\left(K_{\left(i_{1}, i_{2}, \ldots, i_{u-1}\right)}\right)}=\frac{\sum_{\sigma \in K_{\left(i_{1}, i_{2}, \ldots, i_{u}\right)}} \pi_{\sigma}}{\sum_{\sigma \in K_{\left(i_{1}, i_{2}, \ldots, i_{u-1}\right)} \pi_{\sigma}} \pi_{\sigma},} \begin{array}{ll}
\frac{1}{1+(n-u) \theta} i_{u}=u \\
\frac{\theta}{1+(n-u) \theta}, & \text { if } u<i_{u} \leq n
\end{array}
$$

$\forall i_{1}, 1 \leq i_{1} \leq n, \forall i_{2}, 2 \leq i_{2} \leq n, \ldots, \forall i_{u-1}, u-1 \leq i_{u-1} \leq n(3 \leq u \leq n-1)$.
Finally, to illustrate Theorem 2.3, its proof, and the above comments, we give an example.

Example 2.4. Consider the Mallows model on $\mathbb{S}_{3}$ through Cayley metric with $\sigma_{0}=(213)$. By Theorem 2.3, we have (the rows and columns of $P_{1}$ and $P_{2}$ are labeled in lexicographic order)
and

|  | (123) | (132) | (213) | (231) | (312) | (321) |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| (123) | $\left(\frac{1}{1+\theta}\right.$ | $\frac{\theta}{1+\theta}$ |  |  |  |  |
| (132) | $\frac{1}{1+\theta}$ | $\frac{\theta}{1+\theta}$ |  |  |  |  |
| $P_{2}={ }^{(213)}$ |  |  | $\frac{1}{1+\theta}$ | $\frac{\theta}{1+\theta}$ |  |  |
| (231) |  |  | $\frac{1}{1+\theta}$ | $\frac{\theta}{1+\theta}$ |  |  |
| (312) |  |  |  |  | $\frac{1}{1+\theta}$ | $\frac{\theta}{1+\theta}$ |
| (321) |  |  |  |  | $\frac{1}{1+\theta}$ | $\frac{\theta}{1+\theta}$ |

Since, for $P_{1}$, we have

$$
\begin{aligned}
& (123)=(213) \circ(1,2) \circ \mathrm{Id} \in \mathbb{M}_{3,1}, \\
& (132)=(213) \circ(1,2) \circ(2,3) \in \mathbb{M}_{3,1}, \\
& (213)=(213) \circ(1,1) \circ \mathrm{Id} \in \mathbb{M}_{3,1}, \\
& (231)=(213) \circ(1,1) \circ(2,3) \in \mathbb{M}_{3,1}, \\
& (312)=(213) \circ(1,3) \circ \mathrm{Id} \in \mathbb{M}_{3,1}, \\
& (321)=(213) \circ(1,3) \circ(2,3) \in \mathbb{M}_{3,1},
\end{aligned}
$$

it follows that

$$
\begin{aligned}
& K_{(1)}=\{(213),(231)\}, \quad K_{(2)}=\{(123),(132)\}, \\
& K_{(3)}=\{(312),(321)\} .
\end{aligned}
$$

Since, for $P_{2}$, we have

$$
\begin{aligned}
& (123)=(213) \circ(1,2) \circ(2,2) \circ \operatorname{Id} \in \mathbb{M}_{3,2}, \\
& (132)=(213) \circ(1,2) \circ(2,3) \circ \operatorname{Id} \in \mathbb{M}_{3,2}, \\
& (213)=(213) \circ(1,1) \circ(2,2) \circ \operatorname{Id} \in \mathbb{M}_{3,2}, \\
& (231)=(213) \circ(1,1) \circ(2,3) \circ \operatorname{Id} \in \mathbb{M}_{3,2}, \\
& (312)=(213) \circ(1,3) \circ(2,2) \circ \operatorname{Id} \in \mathbb{M}_{3,2}, \\
& (321)=(213) \circ(1,3) \circ(2,3) \circ \operatorname{Id} \in \mathbb{M}_{3,2},
\end{aligned}
$$

it follows that

$$
\begin{array}{ll}
K_{(1,2)}=\{(213)\}, & K_{(1,3)}=\{(231)\}, \\
K_{(2,2)}=\{(123)\}, & K_{(2,3)}=\{(132)\}, \\
K_{(3,2)}=\{(312)\}, & K_{(3,3)}=\{(321)\} .
\end{array}
$$

Further, we have

$$
\begin{aligned}
& \Delta_{1}=\left(\mathbb{S}_{3}\right) \\
& \Delta_{2}=\left(K_{(1)}, K_{(2)}, K_{(3)}\right) \\
& \Delta_{3}=\left(K_{(1,2)}, K_{(1,3)}, K_{(2,2)}, K_{(2,3)}, K_{(3,2)}, K_{(3,3)}\right)
\end{aligned}
$$

Obviously, $\Delta_{1}=\left(\mathbb{S}_{3}\right) \succ \Delta_{2} \succ \Delta_{3}=(\{\sigma\})_{\sigma \in \mathbb{S}_{3}}$. It is easy to see that $P_{1} \in G_{\Delta_{1}, \Delta_{2}}$, $P_{2} \in G_{\Delta_{2}, \Delta_{3}}$, and $\pi_{\sigma}\left(P_{l}\right)_{\sigma \tau}=\pi_{\tau}\left(P_{l}\right)_{\tau \sigma}, \forall l \in\langle 2\rangle, \forall \sigma, \tau \in \mathbb{S}_{3}$. By Theorem 2.3 or direct computation, $P=e^{\prime} \pi$. Since $\pi_{\sigma_{0}}=\frac{1}{Z}$, it is easy to see, using $P=e^{\prime} \pi$,
that $Z=(1+\theta)(1+2 \theta)$. Obviously, $P_{2}$ is a block diagonal matrix and $\Delta_{2}$-stable matrix on $\Delta_{2}$. Moreover, $P_{2}$ is a $\Delta_{2}$-stable matrix, see Definition 1.4. $P_{1}$ is a stable matrix both on $\Delta_{1}$ and on $\Delta_{2}$. By Uniqueness theorem from Păun (2017) or direct computation, we have

$$
\begin{aligned}
P\left(K_{(1)}\right) & =\frac{1}{1+2 \theta}, \quad P\left(K_{(2)}\right)=P\left(K_{(3)}\right)=\frac{\theta}{1+2 \theta} \\
P\left(K_{(1,2)}\right) & =\frac{1}{(1+\theta)(1+2 \theta)}, \\
P\left(K_{(1,3)}\right) & =P\left(K_{(2,2)}\right)=P\left(K_{(3,2)}\right)=\frac{\theta}{(1+\theta)(1+2 \theta)}, \\
P\left(K_{(2,3)}\right) & =P\left(K_{(3,3)}\right)=\frac{\theta^{2}}{(1+\theta)(1+2 \theta)} .
\end{aligned}
$$

If the initial state of chain is $\sigma_{0}, \sigma_{0}=(213)$, then from this state the chain passes in one of the states $(213) \circ(1,1),(213) \circ(1,2),(213) \circ(1,3)$. Suppose that it passed in state $(213) \circ(1,3) .(213) \circ(1,3)=(312)$. From (312) the chain passes in one of the states $(312) \circ(2,2),(312) \circ(2,3)$. Suppose that it passed in state $(312) \circ(2,3) \cdot(312) \circ(2,3)=(321) .(321)$ is the state selected from $\mathbb{S}_{3}$ with our method, having, here, two ( $3-1=2$ ) steps.

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