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HOW FAST DOES A MARKOV CHAIN FORGET THE INITIAL STATE? A DECISION THEORETIC APPROACH.
by

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## ABSTRACT.

Let $\left\{X_{n}\right\}_{n=0}^{\infty}$ be a Markov chain with finite state space. We study the experiment $\mathcal{G}_{n}$ of observing $X_{n}$ in order to obtain information about the initial state $X_{0}$. The sequence $\left\{\xi_{n}\right\}_{n=1}^{\infty}$ is proved to have a limit $\xi_{\infty}^{\infty}$, which is characterized in terms of properties of the given chain. A measure for the information contained in $G_{n}$ is proposed, and an asymptotic expression is derived. If states of the chain are lumped, then some information may get lost. Two methods of measuring the loss are studied, the deficiency and the insufficiency. They are shown to have equal asymptotic properties.

Key words: Markov chain, limit experiment, information, lumping, insufficiency.

## 1. Introduction.

Let $\left\{X_{n}\right\}_{n=0}^{\infty}$ be a Markov chain on the finite state space $X=\{1,2, \ldots, r\}$, with transition matrix $P$. We shall consider a statistical model where $X_{n}$ is observed in order to obtain information about the initial state $X_{0}$. Thus $X_{o}$ is the unknown parameter, and parameter set = sample space $=$ the state space of our chain.

Such models are frequently appearing in genetics. One observes some generation of a population and wants information about previous generations. An example is given in section 3.

The main question is, how can we measure the information ? Le Cam has in [7] introduced an information-distance between experiments, the deficiency, which generalizes Blackwell's ([1]) relation "more informative". The deficiency gives rise to a semi-metric in the set of experiments, the $\Delta$-distance.

Let $\mathcal{G}_{n}$ be the experiment of observing $X_{n}$. The information contained in $\xi_{n}$ of course decreases as $n$ grows. It is shown that the sequence $\left\{\xi_{n}\right\}$ has a limit $\xi_{\infty}$ (in the $\Delta$-metric). Section 4 is devoted to characterizing the limit experiment in terms of properties of the given Markov chain. It turns out that if the chain is aperiodic, then $\xi_{\infty}$ is simply the experiment of observing the recurrent class into which the chain ends up. We note that the limit $\xi_{\infty}$ exists, even though the sequence $\left\{P^{n}\right\}$ diverges.

In section 5 we show how the $\Delta$-distance provides us with a measure of the information contained in $X_{n}$. The quantity studied is $\Delta_{n}=\Delta\left(\xi_{n}, \mathcal{G}_{\infty}\right)$. It is shown that $\Delta_{n}$ asymptotically behaves like $n^{T-1}{ }^{n}$, where

$$
\rho=\max \{|\lambda|: \lambda \text { eigenvalue of } P,|\lambda|<1\}
$$

and $\tau$ is an integer, $\tau \geq 1$.
The computations needed are based on a formula of $P^{n}$ derived from a result on linear operators in finite dimensional spaces taken from Dunford-Schwarz [2] . The formula involves the spectrum of $P$ and is given in section 2.

Section 8 treats the subject of lumping states of a Markov chain. Kemeny and Snell [6] have defined the concept lumpability of a Markov chain. In this paper we study how information gets lost by lumping. The efficiency of lumping may be measured both by deficiency and insufficiency, a concept introduced by Le Cam in [8].

Some results on insufficiency are given in section 7 for finite experiments. We note that the given results are valid for quite more general experiments. The given proofs, however, turn out to be very "clean" and may illustrate the technique of proving the general results.

It is proved in section 8 that in the present case, the deficiency and the insufficiency behave asymptotically in the same manner. As is pointed out by Le Cam in [8], this may not always be the case.

Sections 1-6 corresponds to Lindqvist [10]. The theory is, however, rewritten and some proofs are simplified.

## 2. Preliminaries.

If $A=\left(a_{i j}\right)$ is a $m \times n$-matrix with complex entries, we shall define the norm $\|A\|$ of $A$ by

$$
\|A\|=\max _{i} \sum_{j}\left|a_{i j}\right|
$$

We note that convergence in this norm is equivalent to convergence in each entry.

A $m \times n$-matrix $M=\left(m_{i j}\right)$ is called a Markov matrix if $m_{i j} \geq 0$ for all $i, j$ and $\sum_{j=1}^{n} m_{i j}=1$ for all $i$.

The set of Markov-matrices with dimension $m \times n$ is denoted $\mathbb{N} n_{m, n}$. The set $\mathbb{M}_{\mathrm{m}, \mathrm{n}}$ is compact in the metric space of $\mathrm{m} \times \mathrm{n}-$ matrices with metric induced from || • || .

We shall need a suitable representation of powers of matrices. A complex $r \times r$-matrix $T$ may be considered as a linear operator in $C^{\Upsilon}$ with scalar field $C$. Let $\lambda_{1}, \ldots, \lambda_{S}$ be the distinct eigenvalues of $T$. If $\lambda$ is an eigenvalue, we define the index $\nu(\lambda)$ of $\lambda$ to be the smallest non-negative integer $\nu$ such that $(\lambda I-T)^{\nu}{ }_{x}=0$ for every vector $x$ such that $(\lambda I-T)^{\nu+1} x=0$. We shall put $m_{i}=\nu\left(\lambda_{i}\right) ; i=1, \ldots, s$. It is seen that $m_{i} \geq 1$ for all i . The index of an eigenvalue may alternatively be defined as follows: Let $\Psi(t)$ be the minimal polynomial of $\mathbb{T}$, i.e. $\Psi(t)$ is a complex polynomial with leading coefficient 1 and of lowest degree such that $\Psi(\mathbb{T})=0$. It may be shown (see Gantmacher [4]) that we may write

$$
\Psi(\lambda)=\left(\lambda-\lambda_{1}\right)^{m_{1}}\left(\lambda-\lambda_{2}\right)^{m_{2}}, \ldots,\left(\lambda-\lambda_{S}\right)^{m_{S}}
$$

Hence the index of $\lambda_{i}$ is the multiplicity of $\lambda_{i}$ as a root of the minimal polynomial of $T$.

Theorem 8 of ch. VII in Dunford-Schwarz [2] now states that if $f$ is a polynomial with complex coefficients, then we have the expansion

$$
\begin{equation*}
f(T)=\sum_{i=1}^{s} \sum_{j=0}^{m_{i}^{-1}} \frac{\left(T-\lambda_{i} I\right)^{j}}{j!} f^{(j)}\left(\lambda_{i}\right) Z_{i} \tag{2.1}
\end{equation*}
$$

where $Z_{1}, \ldots, Z_{s}$ are complex $r \times r$-matrices which are independent of $f$ and satisfy

$$
\begin{array}{ll}
\text { (i) } Z_{i} & \neq 0 ; i=1, \ldots, s \\
\text { (ii) } Z_{i}^{2} & =Z_{i} ; i=1, \ldots, s \\
\text { (iii) } Z_{i} Z_{k}=0 ; \text { for } i \neq k \\
\text { (iv) } \sum_{i=1}^{S} Z_{i}=I .
\end{array}
$$

If $T$ is given, then the component matrices $Z_{1}, \ldots, Z_{s}$ may easily be determined, by inserting suitable $f(t)$ in (2.1). In fact, for any $i=1, \ldots, s$ there exists a polynomial $f_{i}(t)$ such that $Z_{i}=f_{i}(T)$. The idea of the construction is to define $f_{i}(t)$ so that $f_{i}^{(j)}\left(\lambda_{k}\right)=0$ for $k \neq i, j=0, \ldots, m_{i}-1$. Putting $f(t)=t^{n}$ in (2.1) we obtain an expansion of $T^{n}$, which is the fundamental formula used in this paper

$$
\begin{equation*}
T^{n}=\sum_{i=1}^{s} \sum_{j=0}^{m_{i}^{-1}} \frac{\left(T-\lambda_{i} I\right)^{j}}{j!} n(j)_{\lambda_{i}}^{n-j_{Z_{i}}} \tag{2.2}
\end{equation*}
$$

where $n^{(r)}=n(n-1) \ldots(n-r+1)$ if $r \geq 1$ and $n^{(0)}=1$.

A finite experiment will here be defined as an experiment with finite parameter set $\Theta=\{1,2, \ldots, r\}$ and finite sample space $x=\{1,2, \ldots, s\}$. Such an experiment is completely given by a matrix $P=\left(p_{\theta j}\right) \in M_{r, s}$; the $\theta$-th row vector defining the probability distribution over $X$ when $\theta$ is the "true" parameter. We shall use the symbol $\mathcal{E}_{\mathrm{P}}$ for the experiment defined by $P$.

In [7] Le Cam introduced an information distance between experiments: the deficiency. The deficiency is defined in terms of risk-functions. It follows, however, from [7] that for finite experiments an equivalent definition may be given as follows:

Let $\mathcal{G}_{P}$ and $\mathcal{E}_{Q}$ be experiments such that $P \in \ln _{r}, s$, $Q \in \mathbb{A}_{I, t}$ - Then the deficiency of $\mathcal{G}_{P}$ relative to $\mathcal{E}_{Q}$ is given by

$$
\delta\left(\xi_{P}, \varepsilon_{Q}\right)=\inf \left\{\|P M-Q\|: M \in M M_{s, t}\right\}
$$

The quantity $\delta\left(\varepsilon_{P}, \xi_{Q}\right)$ is interpreted as the maximal loss of information by observing $\xi_{P}$ instead of $\mathcal{\xi}_{Q}$.

It follows from the above definition that the deficiency equals 0 if and only if $\mathcal{G}_{P}$ is sufficient for $\xi_{Q}$ according to Blackwell's definition (see [1]). if $\delta\left(\xi_{P}, \xi_{Q}\right)=0$, then we say that $\varepsilon_{\mathrm{E}}$ is more informative than $\mathcal{E}_{Q}-$ We write this
 $M \in \mathbb{M}_{S, t}$ such that $P M=Q$. This result will be referred to as the randomization criterion.

Le Cam [7] also introduced a distance in the set of experiments It is defined by

$$
\Delta\left(\xi_{P}, \xi_{Q}\right)=\delta\left(\xi_{P}, \xi_{Q}\right) \vee \delta\left(\xi_{Q}, \xi_{P}\right)
$$

( $V$ denotes maximum) and is easily seen to have the properties of a semi-metric. If $\Delta\left(\xi_{P}, \xi_{Q}\right)=0$, then $\xi_{P}$ and $\xi_{Q}$ are said to be equivalent; in symbols $\mathcal{G}_{P} \sim \mathscr{G}_{Q}$.

Let $\left\{\mathcal{G}^{(n)}\right\}_{n=1}^{\infty}$ be a sequence of experiments and let $\mathcal{G}$ be an experiment, all with the same parameter set.

Then we say that $\mathscr{G}(n)$ converges to $\mathcal{E}$, in symbols $\xi(n) \rightarrow \mathcal{E}$ as $n \rightarrow \infty$ if $\Delta(\xi(n), \mathcal{\xi}) \rightarrow 0$ as $n \rightarrow \infty$. The limit $\xi$ is uniquely determined up to $\Delta$-equivalence. In the case of finite parameter set, to each sequence $\left\{\mathcal{G}^{(n)}\right\}$ of experiments there exists a sub-sequence $\left\{\mathscr{G}\left(n^{\prime}\right)\right\}$ and an experiment $\mathcal{E}$ such that $\mathcal{G}\left(n^{\prime}\right) \rightarrow \mathcal{G}$.

This follows from a remark on p. 228 in Torgersen [11], which states that the set of (standard) experiments is compact in the metric $\Delta$. As a consequence of this we have:

## Lemma 2.1.

Let $\left\{\xi^{(n)}\right\}$ be a sequence of experiments with common finite parameter set $\Theta$. Assume that $\xi^{(1)} \geq \xi^{(2)} \geq \ldots$. Then there is an experiment $\mathcal{G}$ such that

$$
\begin{aligned}
& \text { (i) } \Delta(\xi(n), \xi) \backslash 0 \text { as } n \rightarrow \infty . \\
& \text { (ii) } \xi(n) \geq \xi \quad \text { for each } n .
\end{aligned}
$$

Proof: There is a subsequence $\left\{\xi\left(n^{\prime}\right)\right\}$ and an $\xi$ such that

$$
\begin{equation*}
\Delta\left(\xi\left(n^{\prime}\right), \xi\right) \rightarrow 0 \tag{2.3}
\end{equation*}
$$

Given $n$, choose $n^{\prime}>n$. Then since $\xi^{(n)} \geq \xi\left(n^{\prime}\right)$,

$$
\delta(\xi(n), \xi) \leq \delta\left(\xi(n), \xi\left(n^{\prime}\right)\right)+\delta\left(\xi\left(n^{\prime}\right), \xi\right)=\delta\left(\xi\left(n^{\prime}\right), \xi\right) .
$$

Letting $n^{\prime} \rightarrow \infty$ it is seen that $\delta(\xi(n), \xi)=0$ and hence $\xi^{(n)} \geq \xi$.
It follows that

$$
\begin{aligned}
& \Delta\left(\xi^{(n+1)}, \xi\right)=\delta(\xi, \xi(n+1) \\
& \quad \leq \delta(\xi, \xi(n))+\delta(\xi(n), \xi(n+1)) \\
& \quad=\delta(\xi, \xi(n))=\Delta(\xi(n), \xi) .
\end{aligned}
$$

Hence the sequence $\{\Delta(\xi(n), \xi)\}$ is monotonically decreasing and thus convergent. By (2.3) we must have $\Delta\left(\xi^{(n)}, \xi\right) \downharpoonright 0$.

If the $\xi^{(n)}$ are finite experiments, then under some conditions, the limit experiment will also be a finite experiment.

## Lemma 2.2.

Let $\left\{\xi^{(n)}\right\}$ be a sequence of finite experiments which may be defined by matrices $p^{(n)} \in M M_{r, s_{n}}$. Assume that the sequence $\left\{S_{n}\right\}$ is bounded. If $\left\{\xi^{(n)}\right\}$ converges, then the limit experiment $\mathcal{E}$ is a finite experiment (i.e. may be given by a finite Markov matrix).

Proof. Let $s=\max s_{n}$. Extend the matrices $P^{(n)}$ to $r \times s-$ matrices $Q^{(n)}$ by adding colums of zeros. Of course

$$
\xi_{Q}(n) \sim \xi_{P}(n)
$$

By compactness of $M_{r, s}$, there is a subsequence $Q^{\left(n^{\prime}\right)}$ which converges to a matrix $Q \in \mathbb{M}_{r, s}$. It remains to prove that $\xi_{Q}(n) \rightarrow \xi_{Q}$. Clearly $\Delta\left(\xi_{Q}(n), \xi_{Q}\right) \leq\|Q(n)-Q\|$. Hence $\xi_{Q}\left(n^{\prime}\right) \rightarrow \xi_{Q}$. By uniqueness of limits of experiments, this implies $\xi^{G} \sim \mathcal{E}_{Q}$.

Some other notations which will be used are:
If $B_{1}, \ldots, B_{n}$ are quadratic matrices, then we denote by diag( $B_{1}$, $\ldots, B_{n}$ ) the matrix

$$
\left(\begin{array}{cccc}
B_{1} & 0 & \ldots & 0 \\
0 & B_{2} & \ldots & 0 \\
\vdots & & \ddots & \\
0 & \ldots & \ldots & B_{n}
\end{array}\right)
$$

$M_{0}$ is the set of $r \times r$ Markov matrices $M=\left(m_{i j}\right)$ with identical row vectors, ie. there exists $m_{1}, \ldots, m_{r} \geq 0, \sum m_{j}=1$ such that $m_{i j}=m_{j}$ for all $i, j$.
N. o is the set of real rxr-matrices $N=\left(n_{i j}\right)$ with identical rows and with row sums equal to 0 , ie. $\sum_{j=1}^{r} n_{i j}=0$.
3. The model.

We shall in this paper consider a finite, discrete Markov chain $\left\{X_{n}\right\}_{n=0}^{\infty}$ with state space $x=\{1,2, \ldots, r\}$ and transition matrix $P=\left(p_{i j}\right)$. The elements of $P^{n}$ will be denoted $p_{i j}(n)$, so that for $i, j=1,2, \ldots, r$

$$
\begin{equation*}
p_{i j}^{(n)}=\operatorname{Pr}\left(X_{n}=j \mid X_{0}=i\right) \tag{3.1}
\end{equation*}
$$

We let $\theta=X_{o}$ be the parameter and let our experiment consist in observing $X_{n}$. Thus the parameter set $\Theta=x$. Since by (3.1) the $\theta$-th row vector of $P^{n}$ defines the probability distribution of $X_{n}$ given that $X_{0}=\theta$, it follows that the experiment of observing $X_{n}$ is given by $P^{n}$. We shall denote if by $G_{n}$.

It seems reasonable that the information obtained about $X_{0}$ by observing $X_{n}$ decreases as $n$ grows. In fact, if $0<m \leq n$ then $P^{n}=P^{m} P^{n-m}$ and the randomization criterion shows that $\xi_{m} \geq \xi_{n}$. Hence $\xi_{1} \geq \xi_{2} \geq \ldots$ and by lemma 2.1 there exists an experiment $\xi_{\infty}$ such that $\Delta\left(\xi_{n}, \xi_{\infty}\right){\underset{v}{l}}_{\dot{\prime}}$ and $\xi_{n} \geq \xi_{\infty}$. We remark that the limit experiment $\mathcal{E} \infty$ exists for any $P$. Hence the experiment given by $P^{n}$ will converge even if $P^{n}$ does not converge. In section 4 we shall characterize and give an interpretation of the experiment $\xi_{\infty}$ in terms of the properties of the given Markov chain.

At first sight one may think that the information about $X_{0}$ may be increased by observing the chain on times $n_{1}<n_{2}<\ldots<n_{k}$ instead of merely observing at $n_{1}$. However, by the Markov property, the conditional distribution of $X_{n_{1}}, \ldots, X_{n_{k}}$ given $X_{n_{1}}$ is independent of $\theta$, so $X_{n_{1}}$ is sufficient for the vector
$\left(X_{n_{1}}, \ldots, X_{n_{k}}\right)$. A question that may arise concerning the given model is : Let $P$ and $Q$ be transition matrices of the same dimension $r$. Let $\left\{\xi_{n}\right\}$ and $\left\{\mathcal{F}_{n}\right\}$ be the corresponding sequences of experiments according to the model given here. Assume that $P$ and $Q$ represent equivalent experiments, ie. $\xi_{1} \sim \xi_{1}$. Will this imply that $\xi_{n} \sim \mathcal{F}_{n}$ for all $n$ ? That this is not true is seen from the following simple example :

$$
\begin{array}{ll}
P=\left(\begin{array}{ll}
\frac{1}{2} & \frac{1}{2} \\
1 & 0
\end{array}\right) & ,
\end{array} \quad Q=\left(\begin{array}{cc}
\frac{1}{2} & \frac{1}{2}  \tag{3.2}\\
0 & 1
\end{array}\right) .
$$

clearly $\xi_{1} \sim \tilde{f}_{1}$. That $\xi_{2} \nsim \mathcal{F}_{2}$ is seen from the $\ddagger$-criterion (see Torgersen [11] , for computation see e.g. Lindqvist and Torgersen [9]).

Defining $\#\left(x_{1}, x_{2}\right)=\left|x_{1}-2 x_{2}\right|$ yields

$$
\Psi\left(\xi_{2}\right)=1 \neq \frac{3}{2}=\psi\left(\tilde{f}_{2}\right) .
$$

We remark that $P$ and $Q$, given in (3.2) defining equivalent experiments, have quite different properties as transition matrices for Markov chains. P defines an aperiodic irreducible chain, whence $Q$ defines a chain where state 2 is absorbing.

## Example 3.1.

This example is taken from Feller [3]. We consider a particular pair of genes, say $a$ and A. Each individual belongs to one of the three genotypes (a,a) , (a,A) or (A,A). Our population consists of $N$ individuals. However, we shall introduce the genes themselves as elements of the population, so that we deal with a population of 2 N elements which are either a or A . Here $N$ is a fixed number.

Under the assumption of random mating we may assume that the 2 N genes of any generation are formed in 2 N independent trials : if the parent population consists of $j$ a-genes and $2 N-j$ A-genes, then each trial results in $a$ or $A$ with probabilities

$$
p_{j}=j / 2 N \quad, \quad q_{j}=1-p_{j}
$$

respectively.
Hence we have a Markov chain with state space $X=\{0,1, \ldots, 2 N\}$ and transition probabilites given by the binomial distribution :

$$
p_{j k}=\binom{2 N}{k} p_{j}^{k} q_{j}^{2 N-k} .
$$

The states 0 and $2 N$ are absorbing, i.e. $\{0\}$ and $\{2 N\}$ are recurrent classes, whence $1,2, \ldots, 2 N-1$ are transient states.

In genetics one is often faced with the following problem :
A process of the above type is observed after $n$ generations (we assume that the number of $a-$ and A-genes may be counted by some method), and one wants to draw conclusions about the initial population. This is a special case of the model treated in this paper. By the previous results, the information about the initial generation decreases as $n$ grows.

## 4. The limit experiment $G_{\infty}$.

Let the sequence $\left\{\xi_{n}\right\}$ be given as in section 3. As is noted there, for any transition matrix $P$, the corresponding sequence $\left\{\xi_{n}\right\}$ converges to some experiment $\xi_{\infty}$. By Lemma 2.2, $\xi_{\infty}$ is a finite experiment. It is the purpose of this section to determine the Markov matrix defining $\xi_{\infty}$. The results might have been developed using the expansion (2.2) together with some facts on eigenvalues of transition matrices. However, we shall make a probabilistic approach, using merely elementary properties of Markov chains. The following lemma will be useful:

## Lemma 4.1.

Assume that $Q$ is a limit at some subsequence $\left\{P^{n}\right\}$ of $\left\{P^{n}\right\}$. Then $\xi_{\infty}$ may be represented by $Q$. In particular, if $P^{n}$ converges, then $\mathscr{G}_{\infty}$ may be represented by the limiting matrix.

Proof: Analogous to the proof of Iemma 2.2

The totality of states of a Markov chain may be partitioned into equivalence classes, where the states in an equivalence class are those which communicate with each other. A class of states is called recurrent if the probability of leaving the class is 0 ; it is otherwise called transient. A Markov chain is called irreducible if there is only one equivalence class (which is then recurrent if the state space is finite). The period of a state is a constant in each equivalence class, so it makes sense to deal with the concept period of a class. A class is called aperiodic
if each state has period 1.
Given a Markov chain with state space $X=\{1,2, \ldots, r\}$ we may (eventually by rearranging the states) write $P$ on the "canonical" form

$$
P=\left(\begin{array}{ccccc}
P_{1} & 0 & \ldots & 0 & 0 \\
0 & P_{2} & \ldots & 0 & 0 \\
0 & \ddots & 0 \\
0 & \ldots & \ldots & \cdot & P_{V} \\
0 \\
\hline R & & & & Q
\end{array}\right)
$$

where $P_{1}, \ldots, P_{v}$ represent irreducible Markov chains, corresponding to the $v$ recurrent classes of the chain.

If the chain has no transient states, then $R$ and $Q$ may be removed from (4.1). Otherwise, $R \neq 0$ and $Q$ is a quadratic non-negative matrix (which is not a Markov matrix) .

It is seen that

$$
P^{n}=\left(\begin{array}{ccccc}
P_{1}^{n} & 0 & \cdots & \cdots & 0  \tag{4.2}\\
0 & P_{2}^{n} & \cdots & \cdots & 0 \\
\vdots & & \ddots & 0 & \vdots \\
\vdots & & \ddots & \cdot & \vdots \\
0 & \cdots & \cdots & \cdots & P_{v}^{n} \\
\hline R_{n} & & & & 0 \\
\hline
\end{array}\right)
$$

We shall at the outset make the assumption that the $v$ recurrent classes are all aperiodic. Then by the usual limit theorems of irreducible aperiodic chains, for each $k, P_{k}^{n}$ converges to a Markov matrix. $A_{k}$ such that all rows of $A_{k}$ are equal.

Since there is probability 0 of remaining in some transient state, we conclude that $Q^{n} \rightarrow 0$.

It remains to study the properties of $R_{n}$ as $n$ tends to infinity. For simplicity, let the transient states be fabled $1,2, \ldots, \beta$ and let $q_{j k}$ be the probability of being absorbed in the kith recurrent class, given that one starts in state $j$.

For each $j, \sum_{k=1}^{V} q_{j k}=1$.
Let the row vector of $A_{k}$ be $\bar{a}_{k}$. By the limit theorem for transition probabilities (see Karin [5] hm. 3.1) it follows that $R_{n}$ converges and that the row vector in $\lim R_{n}$ corresponding to transient state $j$ may be written

$$
\left(\bar{b}_{j 1}, \bar{b}_{j 2}, \ldots, \bar{b}_{j v}\right)
$$

where

$$
\bar{b}_{j k}=q_{j k} \bar{a}_{k}
$$

We have now shown that under the given conditions of aperiodicity, $P^{n}$ converges to a matrix $B$ of the form

$$
B=\left(\begin{array}{ccccc}
A_{1} & 0 & \ldots & 0 & 0  \tag{4.3}\\
0 & A_{2} & \ldots & 0 & 0 \\
\vdots & & & & \vdots \\
\dot{0} & & & A_{v} & \dot{0} \\
\hline \bar{b}_{11} & \bar{b}_{12} & \cdots & \bar{b}_{1 v} & 0 \\
\vdots & \vdots & & \vdots & \vdots \\
\bar{b}_{\beta 1} & \bar{b}_{\beta 2} & \cdots & \bar{b}_{\beta v} & 0
\end{array}\right)
$$

By lemma 4.1 $\xi_{n} \rightarrow \xi_{\mathrm{B}}$. The limit experiment $\xi_{\mathrm{B}}$ may be reduced to a minimal sufficient form:

## Theorem 4.2.

Let $P$ be given as in (4.1) and assume that $P_{1}, \ldots, P_{v}$ all represent aperiodic classes of states. Let $q_{j k}$ be given as before. Then the limit experiment $G_{\infty}$ is given by the rxv-matfix

$$
A=\left(\begin{array}{cccc}
I_{1} & 0 & \ldots . & 0 \\
0 & I_{2} \ldots \ldots & 0 \\
\vdots & \ddots & \vdots & \vdots \\
0 & \ldots & I_{v} \\
q_{11} & q_{12} \ldots . & q_{1 v} \\
\vdots & & \vdots \\
q_{\beta 1} & q_{\beta 2} \ldots . & q_{\beta v}
\end{array}\right)
$$

where for $k=1, \ldots, v, I_{k}$ is a column vector of the same dimension as $P_{k}$, with all entries equal to 1 .

Proof: We have to prove that $\xi_{A} \sim \xi_{B}$. The relation
$\xi_{\mathrm{B}} \geq \xi_{\mathrm{A}}$ is established by observing that $\mathrm{BA}=\mathrm{A}$.
Define now a matrix $M \in M_{v, r}$ by

$$
M=\left(\begin{array}{cccc}
\bar{a}_{1} & 0 & \ldots & 0 \\
0 & \bar{a}_{2} \ldots & 0 \\
\vdots & \ddots & \ddots & \bar{a}_{v}
\end{array}\right)
$$

A simple computation shows that $A M=B$, which implies

## Corollary 4.3.

If the Markov chain is irreducible and aperiodic, then $\left\{\xi_{n}\right\}$ converges to the minimal informative experiment.

Remark: The minimal informative experiment is characterized by the fact that the probability distribution does not depend on the parameter $\theta$. Any Markov matrix with equal rows will represent the minimal informative experiment.

The experiment $\mathscr{G}_{\mathrm{A}}$ given in theorem 4.2 may be given a reasonable interpretation as follows: Let $P$ be given as in the theorem. We note that absorbtion into one of the $v$ recurrent classes occurs with probability 1 . Thus we may define an experiment consisting in observing (the label of) the absorbing class. It is not difficult to see that the probability distribution of this experiment is exactly the one given by the matrix $A$.

We shall now discuss periodic chains. At first we assume that $P$ is the transition matrix of an irreducible, aperiodic chain with period $d>1$. $P$ may then (eventually by permuting the states) be represented on the form

$$
\left.P=\left(\begin{array}{cccc}
0 & P_{1} & 0 & \ldots \\
0 \\
0 & 0 & P_{2} \ldots & 0 \\
\vdots & & & \\
0 & 0 & 0 & \ldots \\
P_{d-1} \\
P_{d} & 0 & 0 & \ldots
\end{array}\right) 0.1\right)
$$

where the matrix in block no. (i,j) has dimension
$\alpha_{i} \times \alpha_{j} ; i, j=1, \ldots, d ; \alpha_{i}>0$ for all $i$ and $\sum_{i=1}^{d} \alpha_{i}=r$. By induction (or by a probabilistic line of argument) we show that

$$
\begin{equation*}
P^{d}=\operatorname{diag}\left(Q_{1}, \ldots, Q_{d}\right) \tag{4.4}
\end{equation*}
$$

where $\boldsymbol{Q}_{\boldsymbol{i}}$ has dimension $\boldsymbol{\alpha}_{\boldsymbol{i}} \times \boldsymbol{\alpha}_{\boldsymbol{i}}$
and

$$
\begin{aligned}
& Q_{1}=P_{1} P_{2} \ldots P_{d} \\
& Q_{2}=P_{2} P_{3} \ldots P_{d} P_{1} \\
& \vdots \\
& Q_{d}=P_{d} P_{1} P_{2} \cdot P_{d-1}
\end{aligned}
$$

It is easily verified that $Q_{1}, \ldots, Q_{d}$ correspond to irreducible, aperiodic chains, so (4.4) shows that $P^{d}$ is of the form (4.1) with $R=0, Q=0$. Hence $\left\{P^{n d}\right\}_{n=1}^{\infty}$ converges to a matrix of the form (4.3) with no transient states and $v=d$. Furthermore, by Iemma 4.1 and theorem $4.2\left\{\xi_{n}\right\}$ converges to the experiment given by the matrix

$$
\left(\begin{array}{cccc}
I_{1} & 0 & \cdots & 0 \\
0 & I_{2} & \cdots & 0 \\
\vdots & & \ddots & \vdots \\
0 & \cdots & \cdots & I_{d}
\end{array}\right)
$$

where $I_{i}$ is a column vector of all $1^{\prime}$ s with dimension $a_{i} ; i=1, \ldots, d$.

We remark at once that the limit experiment in the case of an irreducible chain with period $d$ has the same form as the limit
experiment of a chain with $d$ recurrent classes (the recurrent classes correspond to the cyclic classes in the periodic case). This may be explained as follows:

In the former case we may by observing $X_{n}$, since $n$ is known , draw the conclusion as to which cyclic class $X_{o}$ belongs. This corresponds to, in the latter case, knowing that $X_{0}$ belongs to the same recurrent class as $X_{n}$.

We now return to the general case where $P$ is given as in (4.1) and $P_{i}$ represents a periodic chain with period $d_{i} \geq 1 ; i=1, \ldots, v$. By (4.2) and the remarks about irreducible periodic chains we may write (possibly by permuting states in each recurrent class)

$$
P^{d}=\left(\begin{array}{ccccc}
P_{1}^{\prime} & 0 & \cdots & 0 & 0  \tag{4.5}\\
0 & P_{2}^{\prime} & & \vdots & \vdots \\
\vdots & & & & \vdots \\
0 & \cdots & \ldots & P_{\gamma}^{\prime} & 0 \\
\hline R_{d} & & & & Q^{\alpha}
\end{array}\right)
$$

with $d$ equal to the least common multiple (l.c.m.) of $d_{1}, \ldots, d_{V}$ (i.e. $d$ is the least positive integer so that $d_{1}, \ldots d_{v}$ all divide d)
and $\quad \gamma=\sum_{i=1}^{V} d_{i}$.
Now $P_{1}^{\prime}, \ldots, P_{\gamma}^{\prime}$ all represent irreducible, aperiodic chains, so it is seen that $\left\{P^{n d}\right\}_{n=1}^{\infty}$ converges to a matrix of the form $B$ given in (4.3).

Thus by lemma 4.1 we have (this is the main result of this section):

## Theorem 4.4.

Consider a Markov chain with state space $X=\{1,2, \ldots, r\}$. We assume that the transition matrix $P$ is given by (4.1), where $P_{i}$ has period $d_{i} \geq 1$ and $P_{i}$ is written on cyclic form

$$
P_{i}=\left(\begin{array}{ccccc}
0 & P_{i 1} & 0 & \ldots & 0 \\
0 & 0 & P_{i 2} & \ldots & 0 \\
\vdots & & & \\
0 & 0 & 0 & \ldots & P_{i, d_{i}-1} \\
P_{i, d_{i}} & 0 & 0 & \ldots & 0
\end{array}\right)
$$

where the matrix in block no. ( $k, I$ ) has dimension $\alpha_{i, k} \times \alpha_{i, I}$ $\left(k, I=1, \ldots, d_{i}\right)(i=1, \ldots, v)$.

We assume that $Q$ has dimension $\beta \times \beta$, i.e. that the chain has $\beta$ transient states ( $\beta$ may be 0 ).

Let $d$ be the I.c.m. of $d_{1}, \ldots, a_{v}$. Then $P^{d}$ (see (4.5)) is the transition matrix of a Markov chain with $\gamma=\sum_{i=1}^{V} \alpha_{i}$ recurrent classes and $\beta$ transient states.

For $j=1, \ldots, \beta ; k=1, \ldots, \gamma$ let $q_{j k}$ be the probability of being absorbed in the $k$-th recurrent class of this chain, given that one starts in state $j$.

The experiment sequence $\left\{\xi_{n}\right\}$ defined by $P$ now converges to the experiment $\xi_{A}$ given by the matrix

where $I^{h}$ denotes a column vector of dimension $h$ with all entries equal to 1 .

Conversely, to each matrix $A$ of the above form (with $q_{j k} \geq 0, \sum_{k=1}^{Y} q_{j k}=1$ ) there is a transtition matrix $P$ such that $\left\{\xi_{n}\right\}$ corresponding to this $P$ converges to $\mathcal{G}_{A}$.

Proof: The last assertion is proved by noting that for any $A$ on the form (4.6) $\xi_{A}$ is equivalent to an experiment $\xi_{B}$ with $B$ on the form given in (4.3).

Now $B$ is a transition matrix and it is seen that $B^{n}=B$ for $n=1, \ldots$, Hence $\xi_{1} \sim \xi_{2} \sim \ldots \sim \xi_{B}$ so $\xi_{n} \rightarrow \xi_{B}$.

Remark: Since the states of a Markov chain may be permuted, any matrix which by a permutation of rows becomes a matrix of the form (4.6), is a limit experiment of a sequence $\left\{\xi_{n}\right\}$. The representation of the Limit experiment $\mathcal{G}_{\infty}$ given by (4.6) is minimal sufficient since an addition of columns in $A$ and $a$ succeeding multiplication with a Markov matrix would not lead us back to $A$. Hence the given representation is minimal in the sense that any other minimal sufficient representation is given by a permutation of columns in $A$.

What can be said about the original Markov chain by examining the limit experiment ?

By the discussion in this section it follows that the limit experiment depends on the recurrent classes, the periods and the transient states. Given the limit experiment, however, it turns out that to select say the recurrent classes, we need additional information about the chain. As an example, suppose $r=3$ and $A=(1,1,1)^{\prime}$.

Then $\{1,2\}$ is recurrent and $\{3\}$ transient, or $\{1\}$ is recurrent and $\{2,3\}$ transient and so on.

If $A=\left(\begin{array}{ll}1 & 0 \\ 1 & 0 \\ 0 & 1\end{array}\right)$, then
$\{1,2\}$ and $\{3\}$ may be two recurrent classes, or $\{1,2,3\}$ may be a recurrent class with period 2 etc.

We close this section by taking up the example given in section 3 •

## Example 4.5 .

By Feller [3], the probability $\epsilon_{m}$ of being absorbed into state 0 given that the initial population has $m$ a-genes, is given by

$$
\epsilon_{m}=1-m / 2 \mathbb{N} \quad ; \quad m=0,1, \ldots, 2 N
$$

Hence the limit experiment is defined by

$$
A=\left(\begin{array}{cc}
1 & 0 \\
1-\frac{1}{2 N} & \frac{1}{2 N} \\
1-\frac{2}{2 N} & \frac{2}{2 N} \\
: & \\
1-\frac{2 N-1}{2 N} & \frac{2 N-1}{2 N} \\
0 & 1
\end{array}\right)
$$

## 5. Rate of convergence.

In this section we set $\Delta_{n}=\Delta\left(\xi_{n}, \xi_{\infty}\right)$.
Assume for a moment that the Markov chain under consideration is irreducible and aperiodic. By corollary 4.3 , the limit experiment $\xi_{\infty}$ is the least informative experiment. Intuitively speaking, a decision based on this experiment may be considered as the guessing of a parameter without taking observations. Hence $\Delta_{n}$ measures the maximum loss when "guessing" $X_{o}$ instead of basing a decision about $X_{o}$ on $X_{n}$. It may thus seem reasonable to consider $\Delta_{n}$ as a measure of the information contained in $X_{n}$ -

If the chain is not necessarily irreducible and aperiodic, then $\Delta_{n}$ may be given a similar interpretation. More precisely, $\Delta_{n}$ measures the additional information of observing $X_{n}$ instead of merely observing the recurrent class into which the chain ends up.

We snall make use of the expansion of $p^{n}$ given in (2.2). To begin with, we state some results on eigenvalues of non-negative matrices and in particular of transition matrices. (See e.g. Karlin [5]) . It is tacitly understood that every matrix is of finite dimension.

## Theorem 5.1.

(i) Any non-negative matrix $A \neq 0$ has a real eigenvalue $\mu>0$ such that $|\lambda| \leq \mu$ for any eigenvalue $\lambda$ of $A$. (ii) If $P$ is a transition matrix, then $\mu=1$ and the multiplicity of the eigenvalue 1 is equal to the number of recurrent classes associated with $P$.
(iii) If $P$ is the transition matrix of an irreducible periodic Markov chain with period $d$, then the d'th roots of unity are eigenvalues of $P$, each with multiplicity one, and there are no
other eigenvalues of modulus 1 .
We remark that multiplicity here means multiplicity as a root of the characteristic polynomial, whence the formula (2.2) involves the multiplicity as a root of the mininal polynomial (which may be smaller) .

Let $P$ be given as in theorem 4.4. The set of eigenvalues of $P$ is equal to the collection of eigenvaiues of $P_{1}, \ldots, P_{v}, Q$. Since the chain associated with $P$ has $v$ recurrent classes, it follows from theorem 5.1 that the maximum eigenvalue of $Q$ is $\mu<1$ and that any eigenvalue of $P$ with modulus 1 is a root of unity.

We number the distinct eigenvalues of modulus 1 :

$$
\varphi_{1}=1, \varphi_{2}, \ldots, \varphi_{u} .
$$

Let $d$ be given as in theorem 4.4. Then $\varphi_{i}^{d}=1$ for $i=1,2, \ldots, u$. Since $P^{n}$ is bounded as $n \rightarrow \infty$, (2.2) shows that the index of $\varphi_{1}, \ldots, \varphi_{u}$ are each equal to 1 . (This may also be concluded from theorem 9, ch. VII of [2]). Hence, if $\lambda_{1}, \ldots, \lambda_{s}$ are the eigenvalues of $P$ with modulus less than 1 , and with indexes $m_{1}, \ldots, m_{s}$, respectively, then by (2.2)

$$
\begin{equation*}
P^{n}=\sum_{i=1}^{u} \varphi_{i}^{n_{i}} Z_{i}+\sum_{i=1}^{s} \sum_{j=0}^{\mathrm{m}_{i}-1} \frac{\left(P-\lambda_{i} I\right)^{j}}{j!} n^{(j)} \lambda_{i}^{n-j} C_{i} \tag{5.1}
\end{equation*}
$$

As is proved in section 4 by a probabilistic consideration, (5.1) shows that $P^{m d}$ converges to $\sum_{i=1}^{u} Z_{i}$ as $m \rightarrow \infty$.

More general,

$$
\begin{aligned}
& P^{m d+k} \rightarrow \sum_{i=1}^{u} \varphi_{i}^{k} Z_{i} \quad \text { as } \quad m \rightarrow \infty \\
& \text { for } k=0,1, \ldots, d-1 .
\end{aligned}
$$

We are now in position to derive upper and lower bounds for $\Delta_{n}$.

## Definition 5.2.

If $P$ is a transition matrix, then we define the root of $P$ by

$$
\rho=\max \{|\lambda|: \lambda \quad \text { eigenvalue of } P,|\lambda|<1\}
$$

The maximal index $\tau$ of eigenvalues with modulus $\rho$ is called the index of $P$.

In the following theorems we shall assume that $\rho>0$. If $\rho=0$, then by (5.1) from some $n$ on $P^{n}=\sum_{i=1}^{u} \varphi_{i}^{n} Z_{i}$ so $\Delta_{n} \equiv 0$, i.e. $\xi_{n} \sim \xi_{\infty}$.

## Theorem 5.3.

Let $\rho$ and $\tau$ be, respectively, the root and index of $P$. Then there are constants $0<k \leq K<\infty$ such that for any $n=1,2, \ldots$

$$
k n^{\tau-1} \rho^{n} \leq \Delta_{n} \leq K n^{\tau-1} \rho^{n}
$$

Proof: Let $n$ be given. By the Euclidean algorithm there exist non-negative integers $m, k$ such that $n=m d+k$ and $0 \leq \mathrm{k}<\mathrm{d}$. Since $\mathcal{G}_{\infty}$ by lemma 4.1 may be represented by
$\sum_{i=1}^{u} \varphi_{i}^{j} z_{i}$, it follows from (5.1) that

$$
\begin{aligned}
\Delta_{n} & =\underset{M}{\inf \|} \sum_{i=1}^{u} e_{i}^{k} Z_{i} M-P^{n} \| \\
& \leq\left\|\sum_{i=1}^{u} \varphi_{i}^{m d+k_{2}} Z_{i}-P^{n}\right\| \\
& \leq \sum_{i=1}^{S_{i=1}^{m} \sum_{j=0}^{m} n^{(j)}\left|\lambda_{i}\right|^{n-j}\left\|\frac{\left(P-\lambda_{i}\right)^{j}}{j!} c_{i}\right\|} \\
& \leq K n^{T-1} \rho^{n} .
\end{aligned}
$$

Obviously $K$ may be chosen independent of $n$.
It remains to prove the left inequality. Let $n, m$ and $k$ be given as before. Without loss of generality we may assume that $\left|\lambda_{1}\right|=\rho$ and that the index of $\lambda_{1}$ is $\tau$. Then, applying to the inequality $\|A B\| \leq\|A\|\|B\|$, for any $M \in M_{r, r}$ :

$$
\begin{align*}
& \left\|\sum_{i=1}^{u} \varphi_{i}^{k} Z_{i} M-P^{n_{n}}\right\| \geq \frac{1}{\left\|C_{1}\right\|}\left\|C_{1} \sum_{i=1}^{u} \varphi_{i}^{k} Z_{i} M=C_{1} P^{n_{1}}\right\|  \tag{5.2}\\
& =\frac{1}{\left\|C_{1}\right\|}\left\|C_{1} P^{n}\right\|
\end{align*}
$$

since $C_{1} Z_{i}=0$ for $i=1, \ldots, u$.
Since $C_{1}=f_{1}(P)$ for some polynomial $f_{1}(t)$ (see section 2), we have

$$
C_{1} P^{n}=P^{n} C_{1}=\sum_{j=0}^{\tau-1} \frac{\left(P-\lambda_{1} I\right)^{j}}{j!} n(j) \lambda_{1}^{n-j} C_{1}
$$

The right hand side of (5.2) is independent of $M$ so

$$
\begin{align*}
& \frac{\Delta n}{n^{\tau-1} \rho^{n}} \geq \frac{\left\|P^{n} C_{1}\right\|}{\left\|C_{1}\right\| n^{\tau-1} \rho^{n}}  \tag{5.3}\\
= & \frac{1}{\left\|C_{1}\right\|}\left\|\sum_{j=0}^{\tau-1} \frac{\left(P-\lambda_{1} I\right)^{j}}{j!} \frac{n^{n}(j)}{n^{\tau-1}} \frac{\lambda_{1}^{n-j}}{\rho^{n}} C_{1}\right\| \cdot
\end{align*}
$$

The right hand side of (5.3) converges to

$$
\frac{1}{\rho^{\tau-1}\left\|C_{1}\right\|}\left\|\frac{\left(P-\lambda_{1} I\right)^{\tau-1}}{(\tau-1)!} C_{1}\right\|>0 .
$$

But this assures the existence of $k>0$ such that

$$
\frac{\Delta_{n}}{n^{T-1} \rho^{n}} \geq k>0 \quad \text { for all } n
$$

## Corollary 5.4.

$$
\sqrt[n]{\Delta_{n}} \rightarrow \rho \quad \text { as } \quad n \rightarrow \infty
$$

The result of theorem 5.3 may be written

$$
\Delta_{n}=k_{n} n^{\tau-1} \rho^{n}
$$

with $k \leq k_{n} \leq K$ for any $n$.
What can be said about the value of $k_{n}$ ? We shall derive an asymptotic expression for $\Delta_{n}$ in the case of an irreducible, aperiodic chain, and show that in some cases the value of $k_{n}$ is independent of $n$.

Let $M_{0}$ and Ne be defined as in section 2.

## Theorem 5.5.

Let $P$ represent an irreducible, aperiodic chain. Let the eigenvalues of $P$ be $1, \lambda_{1}, \ldots, \lambda_{S}$. Assume that

$$
1>\left|\lambda_{1}\right|=\left|\lambda_{2}\right|=\ldots=\left|\lambda_{p}\right|=\rho>\left|\lambda_{p+1}\right| \geq \cdots \geq\left|\lambda_{s}\right|
$$

and

$$
m_{1}=m_{2}=\ldots=m_{q}=\tau>m_{q+1} \geq \cdots \geq m_{p}(q \leq p)
$$

Let $\quad \lambda_{j}=\rho e^{i \varphi} j ; \quad \varphi_{j} \in[0,2 \pi] ; j=1, \ldots, q$.

Define $\quad T_{n}=\sum_{k=1}^{q} e^{i(n-\tau+1) \varphi_{k}} \frac{\left(P-\lambda_{k} I\right)^{\tau-1}}{(\tau-1)!} C_{k}$.

Then

$$
\Delta_{n}=\inf _{N \in M_{0}}\left\|N-T_{n}\right\|+o\left(n^{\top-1} \rho^{n}\right)
$$

Proof: As is well known from the theory of Markov chains, $\mathrm{P}^{\mathrm{n}}$ converges to a matrix $Z_{1} \in M_{0}$ such that all entries of $Z_{1}$ are strictly positive.

Since $z_{1} \in M_{0} \Rightarrow z_{1} M \in M M_{0}$ for any square Markov matrix M , it is seen that

$$
\Delta_{n}=\inf _{M}\left\|Z_{1} M-P^{n}\right\|=\inf _{M \in M_{0}}\left\|M-P^{n}\right\|
$$

Put $\quad R_{n}=\sum_{k=1}^{q} \frac{\left(P-\lambda_{k} I\right)^{\tau-1}}{(\tau-1)!} n^{(\tau-1)} \lambda_{k}{ }^{n-\tau+1} C_{k}$

$$
=n^{(\tau-1)} \rho^{n-\tau+1} T_{n}
$$

and define $Q_{n}$ so that

$$
\begin{equation*}
P^{n}=Z_{1}+R_{n}+Q_{n} . \tag{5.5}
\end{equation*}
$$

Now

$$
\Delta n=\inf _{M \in M}\left\|M-Z_{1}-R_{n}-Q_{n}\right\|
$$

Let $M$ * be the set of matrices $M-Z_{1}$ obtained as $M$ varies $\operatorname{in} M_{0}$. Then

$$
\Delta n=\inf _{N \in M *}\left\|N-R_{n}-Q_{n}\right\|
$$

By (5.1), (5.4) and (5.5),

$$
\left\|Q_{n}\right\|=o\left(n^{\tau-1} \rho^{n}\right)
$$

and it is thus seen that

$$
\begin{align*}
\Delta_{n} & =\inf _{\mathbb{N}^{\prime} \in \mathbb{M}}\left\|\mathbb{N}-R_{n}\right\|+o\left(n^{\tau-1} \rho^{n}\right) \\
& =n^{(\tau-1)} \rho^{n-\tau+1} \inf _{\mathbb{N}^{\prime} \in M^{*}}\left\|\frac{N}{n^{(\tau-1)} \rho^{n-\tau+1}}-T_{n}\right\| \tag{5.6}
\end{align*}
$$

since $M_{0}$ is a closed set in the set of $r \times r$-matrices, there exists
$N_{0}^{(n)} \in M_{0}$ such that

$$
\inf _{N \in \mathbb{f}}\left\|N-T_{n}\right\|=\left\|N_{0}^{(n)}-T_{n}\right\|
$$

We shall show that from a certain $n$ on, the infimwn in (5.6) is attained by

$$
\begin{equation*}
N^{(n)}=n^{(\tau-1)} \rho^{n-\tau+1} N_{0}^{(n)} \tag{5.7}
\end{equation*}
$$

By definition of $N_{0}(n)$ it suffices to prove that from some $n$ on, $N^{(n)}$ given by (5.7) is a member of $M_{*}$.

The matrices $N_{0}(n)$ are uniformly bounded in $n$. In fact for any $n$

$$
\begin{aligned}
\left\|\mathbb{N}_{0}^{(n)}\right\| & \leq\left\|N_{0}^{(n)}-T_{n}\right\|+\left\|T_{n}\right\| \\
& =\inf _{N \in M_{0}}\left\|\mathbb{N}-T_{n}\right\|+\left\|T_{n}\right\| \\
& \leq 2\left\|T_{n}\right\| \quad \text { by putting } N=0 \\
& \leq 2 \sum_{k=1}^{q}\left\|\frac{\left(P-\lambda_{k} I\right)^{T-1}}{(\tau-1)!} c_{k}\right\|
\end{aligned}
$$

Since all entries of $Z_{1}$ are strictly positive, it is seen that there exists $0<\gamma<1$ such that $M *$ contains any member of $M_{0}$ with all entries of absolute value less than $\gamma$. Hence the fact that $n^{(\tau-1)} \rho^{n-\tau+1} \rightarrow 0$ as $n \rightarrow \infty$ implies that
$N^{(n)} \in \mathbb{N}^{*}$ from a certain $n$ on.
The proof is now complete.

Corollary 5.6.
Assume that $\lambda_{1}$ is real (i.e. $\lambda_{1}= \pm \rho$ ) and that $q=1$. Then there is a constant $c>0$ such that

$$
\frac{\Delta_{n}}{n^{\tau-1} \rho^{n}} \rightarrow c \quad \text { as } n \rightarrow \infty
$$

Proof: By theorem 5.5

$$
\Delta_{n}=k n^{(\tau-1)} \rho^{m-\tau+1}+o\left(n^{\tau-1} \rho^{n}\right)
$$

with $\quad k=\inf _{N \in M_{0}}\left\|\mathbb{N}-\frac{\left(P-\lambda_{1} I\right)^{\tau-1}}{(\tau-1)!} \quad C_{1}\right\| \quad$.
That $k>0$ is a consequence of theorem 5.3.
The corollary follows.

If the $Q_{n}$ defined in the proof of theorem 5.5 equals exactly 0 from a certain $n$ on (this happens if $T=1$ and $P$ has no non-zero eigenvalues with modulus less than $\rho$ ), then theorem 5.5 gives rise to exact expressions for $\Delta_{n}$ from a certain $n$ on.

The proof of theorem 5.5 also shows that if $Q_{n}=0$ and $N_{0}(n) \in M *$ for $n=1,2, \ldots$ then we obtain exact expressions for $\Delta_{n}$ for $n=1,2, \ldots$. To give an example, we shall derive $\Delta_{n}$ for a general $P$ of dimension $2 \times 2$.

## Example 5.7.

$$
\text { Let } \quad P=\left(\begin{array}{cc}
1-\alpha & \alpha \\
\beta & 1-\beta
\end{array}\right)
$$

where $0 \leq \alpha, \beta \leq 1$.

To avoid trivial chains, we shall assume that $\alpha+\beta \neq 0,1,2$. The eigenvalues of $P$ are 1 and $1-\alpha-\beta$ and

$$
P^{n}=\frac{1}{\alpha+\beta}\left(\begin{array}{ll}
\beta & \alpha \\
\beta & \alpha
\end{array}\right)+\frac{(1-\alpha-\beta)^{n}}{\alpha+\beta}\left(\begin{array}{cc}
\alpha & -\alpha \\
-\beta & \beta
\end{array}\right)
$$

i.e. by the terminology of theorem 5.5, $\rho=|1-\alpha-\beta|, \tau=1$,

$$
C_{1}=\frac{1}{\alpha+\beta}\left(\begin{array}{rr}
\alpha & -\alpha \\
-\beta & \beta
\end{array}\right)
$$

and $\varphi_{1}=0$ or $\pi$ according to whether $1-\alpha-\beta$ is $>0$ or $<0$. We shall let $N_{0}$ be the $N \in M_{0}$ minimizing $\left\|N-C_{1}\right\|$.
Putting $\quad N=\left(\begin{array}{ll}c & -c \\ c & -c\end{array}\right)$ yields
$\left\|\mathbb{N}-C_{1}\right\|=2\left|c-\frac{\alpha}{\alpha+\beta}\right| \vee\left|c+\frac{\beta}{\alpha+\beta}\right|$.

A geometrical consideration shows that infimum taken over $c$ of this expression is obtained with

$$
c=\frac{\alpha-\beta}{2(\alpha+\beta)}
$$

which yields

$$
\left\|N_{0}=Z_{1}\right\|=1
$$

Furthermore $N_{0} \in \prod^{*}$ so by theorem 5.5 and the remarks preceding this example we have

$$
\Delta_{n}=\rho^{n} \quad \text { for } \quad n=1,2, \ldots
$$

Example 5.8. (Example 3.1 continued).
Feller has computed the eigenvalues of the transition matrix of the chain. In our notations we have

$$
\rho=1-\frac{1}{2} N \quad \text { and } \quad \tau=1
$$

Furthermore, the requirements of corollary 5.6 are satisfied, so that

$$
\Delta_{n} \sim c\left(1-\frac{1}{2 N}\right)^{n}
$$

for some constant $c>0$.
To end this section, let $\left\{\mathscr{G}_{n}\right\}$ be a sequence of experiments determined by the matrix $P$, and let $\left\{\mathcal{F}_{n}\right\}$ be determined by $Q$. Assume further that $\left\{\zeta_{n}\right\}$ and $\left\{f_{n}\right\}$ have a common limit experiment $\mathcal{G}$.

It is obvious that $\Delta\left(\mathcal{G}_{n}, \mathcal{F}_{n}\right) \rightarrow 0$, since $\Delta\left(\xi_{n}, \hat{F}_{n}\right) \leq \Delta\left(\xi_{n}, \mathfrak{Y}\right)+\Delta\left(\tilde{\xi}_{n} \hat{y}\right)$.
We shall prove a result on the rate of convergence of $\Delta\left(\mathcal{\xi}_{n}, \mathcal{F}_{n}\right)$.
Let $\rho, \tau$ be the root and index of $P ; \sigma, \nu$ be the root and index of $Q$.

## Then we have

## Theorem 5.9.

If
(i) $\rho>0$
or (ii) $\rho=\sigma$ and $\tau>\nu$
then
$\frac{\Delta\left(\xi_{n}, \hat{F}_{n}\right)}{n^{\top-1} p^{n}} \quad$ is bounded away from 0 and $\infty$ as
$\mathrm{n} \rightarrow \infty$.
In particular, $\sqrt[n]{\Delta\left(\xi_{n}, f_{n}\right)} \rightarrow \rho$.

If $\rho=\sigma$ and $\tau=\nu$, then the above results may not hold. In fact we may have

$$
\sqrt[n]{\Delta\left(\varepsilon_{n}, \tilde{f}_{n}\right)} \rightarrow \delta \quad \text { with } \quad 0<\delta<\rho
$$

Remark: If $P=Q$, then of course $\Delta\left(\varepsilon_{n}, \tilde{f}_{n}\right)=0$ for all $n$.

Proof: Assume that (i) or (ii) holds.
Then

$$
\frac{\Delta\left(\xi_{n}, \mathcal{\xi}\right)-\Delta\left(\tilde{F}_{n}, \mathcal{G}\right)}{n^{T-1} \rho^{n}} \leq \frac{\Delta\left(\xi_{n}, \hat{F}_{n}\right)}{n^{T-1} p^{n}} \leq \frac{\Delta\left(\xi_{n}, \mathcal{G}\right)+\Delta\left(\mathcal{F}_{n}, \mathcal{G}\right)}{n^{\tau-1} p^{n}}
$$

The first part of the theorem now follows from theorem 5.3, since

$$
\frac{\Delta\left(\tilde{f}_{n}, \mathcal{Y}\right)}{n^{\tau-1}{ }^{n}} \rightarrow 0 \quad \text { as } n \rightarrow \infty
$$

To prove the last part, we shall apply to an example.

Let $P=\left(\begin{array}{ccc}0 & \frac{1}{2} & \frac{1}{2} \\ \frac{1}{4} & \frac{1}{4} & \frac{1}{2} \\ \frac{1}{2} & \frac{3}{8} & \frac{1}{8}\end{array}\right)$
P represents an irreducible, aperiodic chain and $\left\{\mathcal{G}_{n}\right\}$ converges to the minimal informative experiment. The eigenvalues of $P$ are $1,-\frac{3}{8}-\frac{1}{4}$ and

$$
P^{n}=Z_{0}+\left(-\frac{3}{8}\right)^{n} Z_{1}+\left(-\frac{1}{4}\right)^{n} Z_{2}
$$

with

$$
\begin{aligned}
& Z_{0}=\frac{1}{11}\left(\begin{array}{ccc}
3 & 4 & 4 \\
3 & 4 & 4 \\
3 & 4 & 4
\end{array}\right) \\
& Z_{1}=\frac{1}{11}\left(\begin{array}{ccc}
8 & -4 & -4 \\
8 & -4 & -4 \\
-14 & 7 & 7
\end{array}\right) \\
& Z_{2}=\left(\begin{array}{rrr}
0 & 0 & 0 \\
-1 & 1 & 0 \\
1 & -1 & 0
\end{array}\right)
\end{aligned}
$$

Define $Q=Z_{0}-\frac{3}{8} Z_{1}$
i.e. $\quad Q=\left(\begin{array}{ccc}0 & \frac{1}{2} & \frac{1}{2} \\ 0 & \frac{1}{2} & \frac{1}{2} \\ \frac{3}{4} & \frac{1}{8} & \frac{1}{8}\end{array}\right)$

Then $Q^{n}=z_{0}+\left(-\frac{3}{8}\right)^{n} z_{1} \quad(n=1,2, \ldots)$

$$
\text { and } \quad \rho=\sigma=\frac{3}{8}, \quad \tau=V=1 .
$$

Now

$$
\Delta\left(\xi_{n}, \mathcal{F}_{n}\right) \leq\left\|P^{n}-Q^{n}\right\|=\left\|\left(-\frac{1}{4}\right)^{n_{2}} Z_{2}\right\|=2\left(\frac{1}{4}\right)^{n} .
$$

Since $z_{2} z_{0}=z_{2} z_{1}=0, z_{2}^{2}=z_{2}$ we get for any $3 \times 3$ Markov matrix M

$$
\left\|Q^{n_{M}}-P^{n^{n}}\right\| Z_{2}\|\geq\| Z_{2} Q^{n_{M}}-Z_{2} P^{n_{n}}\|=\|\left(-\frac{1}{4}\right)^{n_{2}}\left\|=\left(\frac{1}{4}\right)^{n_{2}}\right\| Z_{2} \|
$$

so that

$$
\Delta\left(\mathcal{G}_{n}, \tilde{f}_{n}\right) \geq \delta\left(\mathcal{F}_{n}, \mathcal{G}_{n}\right)=\inf _{M}\left\|Q^{n} M-P^{n}\right\| \geq\left(\frac{1}{4}\right)^{n}
$$

Hence

$$
\begin{aligned}
& \left(\frac{1}{4}\right)^{n} \leq \Delta\left(\mathcal{G}_{n}, \xi_{n}\right) \leq 2\left(\frac{1}{4}\right)^{n} \\
& \text { so } \quad \sqrt[n]{\Delta\left(\xi_{n}, \xi_{n}\right)} \rightarrow \frac{1}{4}<\frac{3}{8}=\rho
\end{aligned}
$$

In [12], Torgersen investigates the situation where $\mathcal{G}$ is a dichotomy (ie. the parameter set contains two points) and where $\mathcal{E}_{0}^{n}$ is the experiment of taking $n$ independent observations of $G$. It is shown that if $\mid$ if is the maximal informative experiment, then there exists a number $c(\xi)$ such that

$$
\sqrt[n]{\Delta\left(\left(M, \varepsilon^{n}\right)\right.} \rightarrow c(\xi) \quad \text { as } n \rightarrow \infty
$$

Further it is shown that if $\mathcal{G}, f$ is a pair of nonequivalent dichotomies, then $\sqrt[n]{\Delta\left(\xi^{n}, \mathcal{F}^{n}\right)} \rightarrow c(\xi) \vee c(\mathcal{F})$ as $n \rightarrow \infty$. The result given in theorem 5.9 is quite analogous, apart from the fact that we may have

$$
\sqrt[n]{\Delta\left(\xi_{n}, \tilde{F}_{n}\right)} \rightarrow \delta<\rho \vee \sigma \text { if } \rho=0 .
$$

## 6. The minimizing Markov matrix.

We shall in this section assume that $P^{n}$ converges to a matrix $Z_{o}$. Then $Z_{o}$ is the form (4.3).

Since the set of Markov-matrices is closed in the topology induced by the matrix norm, we may for each $n$ pick out a Markovmatrix $M_{n}$ such that

$$
\Delta_{n}=\left\|Z_{0} M_{n}-P^{n_{n}}\right\|
$$

Since $P^{n}$ converges to $Z_{o}$, it may seem reasonable that $M_{n}=I$ and thus $\Delta_{n}=\left\|Z_{o}-P^{n^{n}}\right\|$.

That this is not true in general, is easily seen in the case $r=2$. Consider example 5.7. For each $n$

$$
P^{n}=Z_{0}+\frac{(1-\alpha-\beta)^{n}}{\alpha+\beta}\left(\begin{array}{cc}
\alpha & -\alpha \\
-\beta & \beta
\end{array}\right)
$$

so that

$$
\left\|Z_{o}-P^{n}\right\|=\frac{2(\alpha V \beta)}{\alpha+\beta} \rho^{n}>\rho^{n}=\Delta_{n} .
$$

From the expression $\Delta_{n}=\inf _{M}\left\|Z_{o} M-P^{n_{n}}\right\|$ we observe that $M$ is involved only through $Z_{0} M$. Hence $\mathbb{M}_{n}$ may be replaced by any matrix $M_{n}^{\prime}$ such that $Z_{o} M_{n}=Z_{o} M_{n}^{\prime}$. Such $M_{n}^{\prime}$ may be defined by $M_{n}^{\prime}=A+M_{n}$ if $Z_{0} A=0$. Since $Z_{o}^{2}=Z_{0} P=Z_{0} I=Z_{0}$, $A$ may be chosen as the difference between two of the matrices $Z_{o}, P$ or I . (In each case one has to check that $M_{n}^{\prime}$ is non-negative). In example 7.11 of [9] it is shown that with

$$
M_{n}=\left(\begin{array}{cc}
1-a_{n} & a_{n} \\
b_{n} & 1-b_{n}
\end{array}\right)
$$

(in the case $r=2$ from example 5.7) we may choose $a_{n}, b_{n}$ as arbitrary numbers in $[0,1]$ satisfying

$$
a_{n} \beta-b_{n} \alpha=\frac{\beta-a_{0}}{2} \rho^{n}
$$

Taking $a_{n}=b_{n}=\frac{1}{2} \rho^{n}$ yields
$M_{n}=I+\frac{1}{2} \rho^{n}\left(\begin{array}{rr}-1 & 1 \\ 1 & -1\end{array}\right)$
so that $M_{n} \rightarrow I$
If $\alpha<\beta$ we may put
$b_{n}=1, \quad a_{n}=\frac{\beta-\alpha}{2 \beta} \rho^{n}+\frac{\rho}{\beta} \quad$ to get

$$
M_{n}=\left(\begin{array}{cc}
1-\frac{\alpha}{\beta} \\
1 & 0
\end{array}\right)+\frac{\beta-\alpha}{2 \beta} \rho^{n}\left(\begin{array}{cc}
-1 & 1 \\
0 & 0
\end{array}\right)
$$

so that $M_{n}-1$.
However, $M_{n}$ converges to a matrix $A$ satisfying $Z_{0} A=Z_{0}$.
We have the following general results:

Theorem 6.1.
Let $\left\{M_{n}\right\}$ be a sequence of minimizing Markov matrices. Then
(i) $Z_{0} M_{n} \rightarrow Z_{0}$ as $n \rightarrow \infty$ with exponential speed.
(ii) There is a subsequence $\left\{M_{n}^{1}\right\}$ which converges to a matrix $A$ satisfying $Z_{O} A=Z_{0}$.

Proof: Define $R_{n}$ so that $P^{n}=Z_{0}+R_{n}$.
Then $\left\|Z_{0} M_{n}-Z_{0}\right\| \leq\left\|Z_{0} M_{n}-Z_{0}-R_{n}\right\|+\left\|R_{n}\right\|=\Delta_{n}+\left\|R_{n}\right\|$. (i) follows.

To prove (ii), note that by compactness there is a subsequence $\left\{M_{n}^{-}\right\}$which converges to a Markov matrix $A$.

Now $Z_{0} A=Z_{0}$, since
$\left\|Z_{0} A-Z_{0}\right\| \leq\left\|Z_{0}\right\|\left\|A-M_{n}^{\prime}\right\|+\left\|Z_{o} M_{n}^{\prime}-Z_{0}\right\|$
and the right hand side tends to 0 .

The result (ii) may easily be improved in a particular case :

Theorom 6.2.
If the Markov chair has no iiansient states, then we may pick out a minimizing sequence $\left\{M_{n}\right\}$ so that $M_{n} \rightarrow Z_{o}$.

Proof: In this case,

$$
Z_{0}=\operatorname{diag}\left(A_{1}, \ldots, A_{v}\right)
$$

where $A_{1}, \ldots, A_{V}$ are quadratic Markov matrices each with identical row vectors. Let $\mathcal{G}$ be the set of all such matrices.

Furthermore, $P^{n}=$ diag $\left(P_{1}^{n}, \ldots, P_{V}^{n}\right)$.
For any $M \in\left(M_{r, r}, Z_{o} M\right.$ is a matrix of the form

$$
Z_{0} M=\left(\begin{array}{c}
B_{1} \\
: 1 \\
B_{v}
\end{array}\right)
$$

where $B_{1}, \ldots, B_{v}$ each has identical rows. Since $M_{n}$ minimizes
$\left\|Z_{0} M_{n}-P^{n}\right\|$, a glance at $P^{n}$ will make it clear that $M_{n}$ ought to be chosen so that $Z_{o} M_{n} \in \mathcal{G}$. In other words, we have

$$
\Delta_{n}=\inf _{M \in \mathscr{G}}\left\|M-P^{n}\right\|
$$

Let $M_{n}^{\prime}$ be so that $\Delta_{n}=\left\|M_{n}^{\prime}-P^{n^{\prime}}\right\|$.
Since $Z_{0} M=M$ whenever $\mathbb{M} \in \mathscr{\mathcal { G }}$, theorem 6.1 (i) implies that $M_{n}^{\prime} \rightarrow Z_{0}$.
That $\left\{M_{n}^{2}\right\}$ is a sequence of minimizing matrices, follows from the fact that $\left\|Z_{o} M_{n}^{1}-P^{n_{n}}\right\|=\left\|M_{n}^{1}-P^{n_{n}}\right\|$.

## 7. Sufficiency and insufficiency in finite experiments.

Notations:
Let $A=\left\{A_{1}, \ldots, A_{k}\right\}$ be a partition of the space $x=\{1,2, \ldots, s\} \quad(k \leq s)$.

By an A-Markov matrix $W=\left(w_{i j}\right)$ we shall mean a matrix $w \in \mathbb{M}_{k, s}$ such that for $j=1, \ldots, s ; i=1, \ldots, k$

$$
w_{i j}=0 \quad \text { if } \quad j \notin A_{i}
$$

As an example, let $s=3$ and $A=\{\{1,2\},\{3\}\}$.
Then $W=\left(\begin{array}{ccc}\frac{1}{4} & \frac{3}{4} & 0 \\ 0 & 0 & 1\end{array}\right)$ is an A-Markov matrix.
By $V_{A}$ we shall denote the element $\left({ }_{\left({ }_{j i j}\right)}\right.$ of $M M_{s, k}$ defined by

$$
v_{j i}=\left\{\begin{array}{lll}
1 & \text { if } & j \in A_{i} \\
0 & \text { if } & j \notin A_{i}
\end{array}\right.
$$

In the above example, $\mathrm{V}_{\mathrm{A}}=\left(\begin{array}{ll}1 & 0 \\ 1 & 0 \\ 0 & 1\end{array}\right)$.
By $\quad U_{A}$ we shall denote the A-Markov matrix where each nonzero entry of the $i^{\prime}$ th row equals $\left(\# A_{i}\right)^{-1}$.

In our example $\quad U_{A}=\left(\begin{array}{ccc}\frac{1}{2} & \frac{1}{2} & 0 \\ 0 & 0 & 1\end{array}\right)$.
When no confusion can arise, we may write $V$ and $U$ instead of $V_{A}$ and $U_{A}$.

Let $\mathcal{G}$ be a finite experiment as defined in section 2, and let $P \in M_{r, s}$ be the matrix defining $\mathcal{G}$.

We start this section by giving a criterion of sufficiency in terms of matrices.

A statistic defined on the sample space $X=\{1,2, \ldots, s\}$ of $\mathcal{G}$ may be identified with a partition $A=\left\{A_{1}, \ldots, A_{k}\right\}$ of $X$. (A consists of the subsets of $X$ to which the statistic assigns the same value). Hence we shall deal with the concept sufficient partition instead of sufficient statistic.

The experiment of observing $A_{i}$ is a reduction of the experiment $\mathcal{G}$. It will be denoted by $\mathcal{E}^{A}$ and is defined by $P V_{A}$.

## Theorem 7.1.

Let $\xi$ be given as above. A partition $A$ is sufficient for $P$ if and only if there is a A-Markov matrix $W$ such that

$$
P=P V_{A} W
$$

## Remark:

The criterion (7.1) may be reformulated as follows : Take out the columns corresponding to a certain partition. The resulting matrix has the property that the elements of each row are equally proportioned. As an example, let

$$
P=\left(\begin{array}{ccc}
\frac{1}{2} & \frac{1}{4} & \frac{1}{4} \\
\frac{2}{3} & \frac{1}{3} & 0
\end{array}\right) \cdot \text { Then } A=\{\{1,2\},\{3\}\} \text { is sufficient }
$$

with $W=\left(\begin{array}{ccc}\frac{2}{3} & \frac{1}{3} & 0 \\ 0 & 0 & 1\end{array}\right)$

Proof: $P V_{A}$ is an element of $M_{r, k}$, the $\theta-$ th row vector defining the probability distribution over $A_{1}, \ldots, A_{k}$ when $\theta$ is the "true" parameter. By the definition of $\Delta$-sufficiency (see e.g. Lindqvist and Torgersen [9] ch. 8), A is sufficient if and only if $\xi^{A} \sim \xi$. If the condition (7.1) holds, then this is obviously true.

Conversely, assume that A is sufficient. By definition, the conditional distribution over $X$ given $A$ may be specified independent of $\theta$. Let $W$ be the matrix defining this conditional distribution. Then $W$ is an A-Markov matrix and we have $P=P V_{A} W$.

Le Cam has in [8] introduced a concept called insufficiency, in order to measure the loss of information incurred by restricting ourselves to a function of the observations (statistic). This measure is so defined that a statistic is sufficient if and only if the insufficiency equals 0 . As will be seen, the insufficiency happens to be generaliy larger than the deficiency. We state and prove some results on insufficiency for finite experiments.

## Definition 7.2.

Let $\xi$ be given by the matrix $P$. Let $A$ be a partition of $x$. By the insufficiency of $A$, denoted $\eta(A, \xi)$ we shall mean the quantity

$$
\eta\left(A, \xi^{( }\right)=2 \underset{Q}{\inf \|Q-P\|}
$$

where infimum is taken over all matrices $Q \in M_{s, r}$ such that A is sufficient for $Q$. The set of such matrices is non-void, since it contains any $Q$ with identical rows.

Theorem 7.3.

$$
\Delta\left(\xi^{A}, \xi\right) \leq \eta(A, \xi)
$$

Proof: Assume that $A$ is sufficient for $Q$. By theorem 7.1 there exists an A-Markov matrix $W$ such that

$$
Q=Q V W \text {. Furthermore, }
$$

$\Delta\left(\xi^{A}, \xi\right)=\inf _{M}\|P V M-P\|$

$$
\begin{aligned}
& \leq \inf _{M}\|P V M-Q\|+\|Q-P\| \\
& \leq\|P V W-Q V W\|+\|Q-P\| \\
& \leq\|P-Q\| \cdot\|V\| \cdot\|W\|+\|Q-P\| \\
& =2\|Q-P\| \cdot
\end{aligned}
$$

Hence $\Delta(\xi \mathrm{A}, \xi) \leq \eta(\mathrm{A}, \xi)$.

Theorem 7.3 gives a lower bound for $\eta(A, \xi)$. We shall, however, also need an upper bound.

$$
\text { Define } \quad \gamma(A, \xi)=\inf _{W}\left\|P V_{A} W-P\right\|
$$

where infimum is taken over all A-Markov matrices $W$. Then we have:

Theorem 7.4 .

$$
\gamma(A, \xi) \leq \eta(A, \xi) \leq 2 \gamma(A, \xi) .
$$

Proof: Let $W$ be an A-Markov matrix and define $Q=P V W$. Then $A$ is sufficient for $Q$ by theorem 7.1, since

$$
Q V W=P V W V W=P V W=Q \quad \text { (we have } W V=I) .
$$

Hence

$$
\eta(A, \xi) \leq 2\|P V W-P\| \text { so } \eta(A, \xi) \leq \gamma(A, \xi) \text {. }
$$

Conversely, let $Q$ be such that $A$ is sufficient for $Q$. Let $W$ be given according to theorem 7.1, so that $Q=Q V W$. Now $P V W=Q+(P V-Q V) W$. Subtracting $P$ on either side yields

$$
\gamma(A, \mathcal{E}) \leq\|P V W-P\| \leq 2\|Q-P\| \text {. }
$$

Hence $\gamma(A, \xi) \leq \eta(A, \xi)$ and the proof is complete.

## Example 7.5.

Assume that $A=\{x\}$. This is an extreme case. The experiment $\mathcal{Q}_{\mathrm{A}}$ is now the minimal informative experiment. Thus we have

$$
\Delta(\xi A, \xi)=\inf _{M \in M_{0}}\|M-P\|
$$

On the other hand, $V_{A}=(1,1, \ldots, 1)$ ', so that any $Q$ such that A is sufficient, has identical rows.

Hence

$$
\eta(A, \xi)=\operatorname{iinf}_{M \in M_{0}}\left\|M_{0}-P\right\|=2 \Delta(\xi A, \xi)
$$

Similarly it is seen that $\gamma(A, \xi) \cdot \Delta(\xi A, \xi)$.

The above ideas may be applied to our Markov chain experiments. We have:

## Theorem 7.6.

Consider a Markov chain with v recurrent, aperiodic classes. In other words, we may assume that $P$ is given by

$$
P=\operatorname{diag}\left(P_{1}, \ldots, P_{v}\right) .
$$

Let the partition $A$ of $X$ be the collection of recurrent classes. Then

$$
\eta\left(A, \mathcal{E}_{n}\right)=2 \Delta_{n}
$$

(where $\Delta_{n}=\Delta\left(\xi_{n}, \xi_{\infty}\right)$ as defined in section 5).

Proof: By the proof of theorem 6.2,

$$
\Delta_{n}=\inf _{M \in \varphi}\left\|M-P^{n}\right\|
$$

( 6 is defined in the proof).
Now

$$
\begin{equation*}
\eta\left(A, Q_{n}\right)=2 \inf _{Q}\left\|Q-P^{n}\right\| \tag{7.2}
\end{equation*}
$$

Obviously (see remark of theorem 7.1), A is sufficient for any $Q \in \mathcal{G}$. On the other hand, let $Q$ be an arbitrary matrix so that $A$ is sufficient for $Q$. We contend that there is a $Q^{\prime} \in \mathcal{G}$ such that $\left\|Q^{\prime}-P^{n}\right\| \leq\left\|Q-P^{n}\right\|$. The method of constructing such a $Q^{\prime}$ is illustrated by defining the first row of $Q^{\prime}$.

Assume that $P_{1}$ has dimension $\alpha \times \alpha$ and let $\left(q_{1}, \ldots, q_{r}\right)$ be the first row of $Q$. We define the first row
$\left(q_{1}^{\prime}, q_{2}^{\prime}, \ldots, q_{\alpha}^{\prime}, 0, \ldots, 0\right)$ of $Q^{\prime}$ by sharing the amount $q_{\alpha+1}+\ldots+q_{r}$ between $q_{1}, \ldots, q_{\alpha}$ and adding in such a way that the proportions are not altered, ie.

$$
q_{1}: q_{2}: \ldots: q_{\alpha}=q_{1}^{\prime}: q_{2}^{\prime}: \ldots: q_{\alpha}^{\prime} .
$$

This is obtained by defining

$$
q_{j}^{\prime}=\frac{q_{j}}{q_{1}+\ldots+q_{\alpha}} ; j=1, \ldots, \alpha .
$$

We remark that $q_{1}+\ldots+q_{\alpha}$ may always be assumed to be $>0$, since $q_{1}+\ldots+q_{\alpha}=0$ would imply $\left\|Q-P^{n}\right\|=2$ and we have (see Torgersen [10])

$$
\inf _{Q \in \zeta_{0}}\left\|Q-P^{n}\right\| \leq 2-\frac{1}{r}
$$

The remaining rows of $Q^{\prime}$ are defined in a similar way, so that we end up with

$$
Q^{\prime}=\operatorname{diag}\left(Q_{1}^{\prime}, \ldots, Q_{V}^{\prime}\right)
$$

Since $A$ is sufficient for $Q$, it follows from the remark of theorem 7.1 that $A$ is sufficient for $Q^{\prime}$. Furthermore, since for any $j$ the elements of the rows of $Q_{j}$ are equally proportioned and have sum $1, Q_{j}^{\prime}$ has identical rows. Hence $Q^{\prime} \in \mathscr{G}$. That $\left\|Q^{\prime}-P^{n}\right\| \leq\left\|Q-P^{n_{1}}\right\|$ follows from the fact that $P^{n}$ is of the same diagonal form as $Q^{\prime}$ and has zero's outside the diagonal. The theorem follows from (7.2).

## 8. Lumping of states.

Kemeny and Snell have in $\S 6.3$ of [6] introduced the concept of lumping states of a Markov chain.

There is given a Markov chain $\left\{\mathrm{X}_{n}\right\}$ on the state space $x=\{1,2, \ldots, r\}$ with transition matrix $P$. Let $A=\left\{A_{1}, \ldots, A_{k}\right\}$ be a partition of the state space.

Define a stochastic process $\left\{Y_{n}\right\}$ on $A$ by taking $Y_{n}$ as the set $A_{j}$ containing $X_{n}$. In general, $\left\{Y_{n}\right\}$ is no Markov chain; in a number of situations, however, it will in fact be. For examples, we refer to [6].

According to Kemeny and Snell, a Markov chain is said to be Iumpable w.r.t. a partition $A$ (in short: A-Iumpabie) if for each initial probability vector $\pi,\left\{Y_{n}\right\}$ constitutes a Markov chain with transition probabilities independent of $\pi$ •

If the chain given by $P$ is A-lumpable, then the transition matrix of the lumped chain is given by $\hat{P}=U_{A} P V_{A}$ (for notations, see section 6).

Theorem 6.3.5 of [6] states that a chain is A-lumpable if and only if

$$
\begin{equation*}
V_{A} U_{A} P V_{A}=P V_{A} \tag{8.1}
\end{equation*}
$$

In this case we have

$$
\mathrm{P}^{\mathrm{n}}=U P^{\mathrm{n}} \mathrm{~V} .
$$

As iṣ seen from section 4, the information contained in a Markov chain depends on the eigenvalues of the transition matrix. We shall state and prove a result on the eigenvalues of lumped chains:

## Theorem 8.1.

Let $P$ represent a chain which is lumpable w.r.t. a partition $A$. Then any eigenvalue of $\hat{P}=U_{A} P V_{A}$ is an eigenvalue of $P$. In other words, 令 "inherits" eigenvalues from $P$.

Proof: If $\lambda$ is an eigenvalue of $\hat{P}$, then there exists a vector x such that $\mathrm{P}_{\mathrm{x}}=\lambda \mathrm{x}$. This is equivalent to $U P V \mathrm{x}=\lambda \mathrm{x}$, which implies $V U P V x=\lambda V x$. By (8.1) this is equivalent to $P V x=\lambda V x$ so that $\lambda$ is an eigenvalue of $P$ with corresponding eigenvector $V x$

Iumping of the states leads us to a reduction of the experiment $\xi_{n}$. We shall now consider the experiment $\mathcal{F}_{n}$ with sample space $A=\left\{A_{1}, \ldots, A_{k}\right\}$ and consisting in observing $Y_{n}$. S is defined by the $r \times k$ Markov matrix. $V \hat{P}{ }^{n}$, the $\theta$-th row vector of which is the row vector of $\hat{\mathbb{P}}^{n}$ corresponding to the $A_{j}$ to which $\theta$ belongs. An easy verification shows that $\hat{V P}^{n}=P^{n} V$. The experiment $\mathcal{F}_{n}$ has the same parameter set as $\mathcal{\xi}_{n}$, so that $\mathcal{G}_{n}$ and $\mathcal{F}_{n}$ are comparable as regards information.

In some practical situations one may find it convenient observing $Y_{n}$ instead of $X_{n}$. This may happen if $r$ is large and it is difficult (or expensive) to get an accurate value for $X_{n}$. As seems reasonable, $\mathcal{E}_{n} \geq \mathcal{F}_{n}$. This is seen from the randomization criterion, since $\hat{V P}^{n}=P^{n}$.

Next, it will be of interest to estimate how much information will get lost by observing $Y_{n}$ instead of $X_{n}$. This may be measured by $\Delta\left(\mathcal{\zeta}_{n}, \tilde{\mathcal{F}}_{n}\right)$, which will be discuissed later. We describe now an alternative approach :

For simplicity, assume that $\left\{\xi_{n}\right\}$ converges to the least informative experiment $O$, and let $\left\{\mathcal{F}_{n}\right\}$ be defined as before. Then $\left\{\hat{f}_{n}\right\}$ converges to the least informative experiment, since $\xi_{n} \geq f_{n} \geq \theta$.

As is noted in section 5, $\Delta_{n}=\Delta\left(\xi_{n}, \theta\right)$ measures the information contained in $X_{n}$.

Define $\gamma_{n}=\Delta\left(\mathcal{S}_{n}, \theta\right)$. Examining $\hat{V P}^{n}$ we will see that $\gamma_{n}=\Delta\left(\tilde{\mathcal{F}}_{n}, \tilde{\mathscr{\theta}}\right)$, where $\mathcal{F}_{n}$ is the experiment with parameter set $\{1, \ldots, k\}$ defined by $\widehat{\mathrm{P}}^{\mathrm{n}}$ and $\tilde{\theta}$ is the least informative experiment over the same parameter set.

Hence, $\rho$ being the root of $P$ and $\mu$ being the root of $\hat{P}$,

$$
\sqrt[n]{\frac{\Delta\left(f_{n}, \theta\right)}{\Delta\left(\xi_{n}, \theta\right)}} \rightarrow \mu / \rho \quad \text { as } n \rightarrow \infty
$$

The quantity $\mu / \rho$ measures the "goodness" of our lumping as regards the information we obtain about $\theta$.

By theorem 8.1, $\mu / \rho \leq 1$ and $\mu / \rho=1$ if and only if $\stackrel{A}{P}$ inherits an eigenvalue with modulus $\rho$. Hence, we should always lump in such a way that $\hat{P}$ inherits the root of $P$.

## Example 8.2.

We shall consider the chain studied in [6] p.29, ex. 8 . $X_{n}$ represents the weather in the Land of $O z$ on the $n ' t h$ day ; the possible values being

$$
1=\text { rain }, 2=\text { sunshine }, \quad 3=\text { snow }
$$

The transition matrix is assumed to be

$$
P=\left(\begin{array}{lll}
\frac{1}{2} & \frac{1}{4} & \frac{1}{4} \\
\frac{1}{2} & 0 & \frac{1}{2} \\
\frac{1}{4} & \frac{1}{4} & \frac{1}{2}
\end{array}\right)
$$

The eigenvalues of $P$ are $1, \frac{1}{4},-\frac{1}{4}$. A computation (see Lindqvist [10], p.59) shows that $\Delta_{n}=\frac{4}{3}\left(\frac{1}{4}\right)^{n}$. This chain is Iumpable w.r.t. $A=\{\{1,3\},\{2\}\}$, the lumped chain being given by the transition matrix

$$
\hat{P}=U P V=\left(\begin{array}{ll}
\frac{3}{4} & \frac{1}{4} \\
1 & 0
\end{array}\right)
$$

The root of $\hat{P}$ equals $\frac{1}{4}$, so by example $5.7, \quad \gamma_{n}=\left(\frac{1}{4}\right)^{n}$. Considering $\Delta_{n}$ and $\gamma_{n}$ as measures of information, this shows that very little information is lost by lumping the states.

Of course states may be lumped regardless of the resulting chain being a Markov chain or not.

Let $\left\{X_{n}\right\}, A$ and $\left\{Y_{n}\right\}$ have the same meaning as in the beginning of this section.

The Markov-matrix defining the probability distribution of $Y_{n}$ is obtained from the one defining $X_{n}$ by adding the columns corresponding to each element of our partition. Hence $\mathcal{F}_{n}$ (the experiment of observing $Y_{n}$ ) is given by $P^{n} V_{A}$. As is seen in the case of lumpability, $\mathcal{G}_{n} \geq \mathcal{F}_{n}$.

To ask if $\mathcal{G}_{n} \sim \mathcal{F}_{n}$ is the same as asking if $Y_{n}$ is sufficient for $X_{n}$. Applying to theorem 7.1 we may prove by induction that if $Y_{1}$ is sufficient for $X_{1}$, then $Y_{n}$ is sufficient for $X_{n}$ for each $n$. Hence, if $Y_{1}$ is sufficient for $X_{1}$, then it makes sense to say that $A$ is sufficient for $\left\{X_{n}\right\}$.

It is noted that sufficiency of a partition $A$ and lumpability of A are quite different concepts. Lumpability will in general not imply sufficiency and conversely.

We recall that in the case of lumpability we had $\mathcal{F}_{1} \geq \mathcal{F}_{2} \geq \ldots$ (this follows from the randomization criterion since $\quad n$ is given by $\frac{\hat{V P}^{n}}{}$ ).

If $A$ is sufficient, then $\mathcal{G}_{1} \sim \mathcal{F}_{1}, \varepsilon_{2} \sim \mathcal{F}_{2}, \ldots$ so $\mathscr{F}_{1} \geq \mathcal{F}_{2} \geq \ldots$ since $\xi_{1} \geq \xi_{2} \geq \ldots$.

If $A$ is neither sufficient nor the chain is lumpable, then, however, this may not be true. For example, define

$$
P=\left(\begin{array}{ccc}
\frac{1}{2} & \frac{1}{4} & \frac{1}{4} \\
0 & \frac{1}{2} & \frac{1}{2} \\
\frac{1}{2} & 0 & \frac{1}{2}
\end{array}\right) \text { and let } A=\{\{1,2\},\{3\}\} .
$$

$\operatorname{Then} P V=\left(\begin{array}{cc}\frac{3}{4} & \frac{1}{4} \\ \frac{1}{2} & \frac{1}{2} \\ \frac{1}{2} & \frac{1}{2}\end{array}\right) \quad P^{2} V=\left(\begin{array}{cc}\frac{5}{8} & \frac{3}{8} \\ \frac{1}{2} & \frac{1}{2} \\ \frac{5}{8} & \frac{3}{8}\end{array}\right)$

By the $\psi$-criterion (see Torgersen [11]) $\hat{f}_{1} \geq \mathfrak{F}_{2}$ if and only if $\psi\left(\mathcal{F}_{1}\right) \geq \Psi\left(\mathcal{F}_{2}\right)$ for all $\psi \in \psi$. Defining $\psi\left(x_{1}, x_{2}, x_{3}\right)=$ $\left|x_{2}-x_{3}\right|$, however, we get $\Psi\left(\xi_{1}\right)=0, \Psi\left(\tilde{f}_{2}\right)=\frac{1}{4}$ (for computaction of $\Psi\left({ }_{i}\right)$, see [9]).

If the relation $\mathcal{F}_{1} \geq \mathcal{F}_{2} \geq \ldots$ does not hold, then we shall see that a limit experiment $\mathcal{f}_{\infty}$ such that $\hat{f}_{n} \rightarrow \mathcal{F}_{\infty}$ may not exist. If $P^{n}$ converges, then of course $\mathcal{F}_{\infty}$ exists and is given by $\lim P_{n} V_{A}$. Thus a sequence $\left\{\mathcal{F}_{n}\right\}$ having no limit necessarily springs from a periodic chain. A very simple example is given on $X=\{1,2,3\}$ by

$$
P=\left(\begin{array}{lll}
0 & \frac{1}{2} & \frac{1}{2} \\
1 & 0 & 0 \\
1 & 0 & 0
\end{array}\right)
$$

Here $\{1\}$ and $\{2,3\}$ are cyclic classes.
We have $P=\left(\begin{array}{ccc}1 & 0 & 0 \\ 0 & \frac{1}{2} & \frac{1}{2} \\ 0 & \frac{1}{2} & \frac{1}{2}\end{array}\right) \quad$ and it is
seen that

$$
P=P^{3}=P^{5}=\ldots, \quad P^{2}=P^{4}=P^{6}=\ldots .
$$

Hence the limit experiment $\mathcal{F}_{\infty}$, if it exists, may be given by either $P V$ or $P^{2} V$. These experiments are, however, not equivalent. Consequently, the sequence $\left\{\tilde{F}_{n}\right\}$ has no limit.

As is noted earlier, the effectiveness of a restriction to an experiment $\mathcal{F}_{n}$ ought to be measured by $\Delta\left(\mathscr{G}_{n}, \mathcal{F}_{n}\right)$. However, since $Y_{n}$ is a function of $X_{n}$, the insufficiency of $Y_{n}$ may
be considered (more precisely the insufficiency of the partition A defined by $\left.Y_{n}\right), \eta\left(A, \mathcal{G}_{n}\right)$. For simplicity, put $\varepsilon_{n}=\Delta\left(\mathcal{G}_{n}, \mathcal{F}_{n}\right)$, $\eta_{n}=\eta\left(A, E_{n}\right)$.

To avoid the difficulties occurring if $\left\{\hat{\jmath}_{n}\right\}$ does not converge, we will in the following assume that $p^{n}$ converges. Further we shall assume that $\left\{\oint_{n}\right\}$ and $\left\{\mathcal{f}_{n}\right\}$ have a common limit experiment $G$."If this assumption is not satisfied, then lumping seems to have little practical interest.

$$
\begin{array}{r}
\text { Of course } \varepsilon_{n} \rightarrow 0 \text { as } \dot{n} \rightarrow \infty \text {, since } \\
\epsilon_{n} \leq \Delta\left(\xi_{n}, \mathcal{G}\right)+\Delta\left(\mathcal{F}_{n}, \mathcal{G}\right) .
\end{array}
$$

Theorem 7.3 implies that $\epsilon_{n} \leq \eta_{n}$. The rest of this section will be devoted to showing that $\sqrt[n]{\epsilon_{n}}$ and $\sqrt[n]{\eta_{n}}$ have a common limit as $\mathrm{n} \rightarrow \infty$, i.e. that $\epsilon_{\mathrm{n}}$ and $\eta_{\mathrm{n}}$ behave asymptotically in the same manner.

In general, however, it may happen that the $\Delta$-distance is small and the $\eta$-insufficiency large. An illustrative example is given in Le Cam [8] p. 46, (where the deficiency behaves like $\frac{C_{1}}{n}$ and the insufficiency behaves like $\frac{C_{2}}{\sqrt{n}}$. We shall need the succeeding lemma :

## Lemma 8.3.

Let $A_{1}, A_{2}, \ldots, A_{z}, B_{1}, B_{2}, \ldots, B_{z} \in M_{r, r}$ and assume that there exists no $M \in M_{r, r}$ such that

$$
A_{i} M=B_{i} \quad \text { for } \quad i=1, \ldots, z
$$

Then there is a number $\varphi>0$ such that for any $M \in M_{r, r}$

$$
\left\|A_{i} M=B_{i}\right\| \geq \varphi \text { for at least one } i=1, \ldots, z \text {. }
$$

Proof: The real function $f$ on $M_{r, r}$ defined by $f(M)=\sum_{i=1}^{Z}\left\|A_{i} M-B_{i}\right\|$ is a continuous function on a compact space. Since $f(M)>0$ on $M_{r, r}$, there exists $a>0$ so that $f(M) \geq a$ for all $M$. Put $\varphi=a / z$.

## Theorem 8.4.

$$
\text { There exists } \nu \geq 1,0 \leq \delta \leq 1,0<k \leq K<\infty
$$

so that

$$
\mathrm{cn}^{\nu-1} \delta^{n} \leq \epsilon_{n} \leq C_{n}^{\nu-1} \delta^{n} \quad \text { for } n=1,2, \ldots .
$$

Proof : We order the eigenvalues of $P$ according to decreasing modulus. Note that since $P^{n}$ is required to converge, 1 is the only eigenvalues of $P$ with modulus 1 . Furthermore it has index 1 .

Let $1<\rho_{1}<\rho_{2}<\ldots<\rho_{t}<0$ be the different moduli of eigenvalues of $P$. We have $\varepsilon_{n}=\inf _{M}\left\|P^{n} N M-P^{n}\right\|$, the infimum taken over $\left(M_{k, r}\right.$. Let $Z_{o}=\lim P^{n}$. Then

$$
\begin{equation*}
P^{n} V M-P^{n}=Z_{0}(V M-I)+\sum_{i=1}^{s} \sum_{j=0}^{m_{i}^{-1}\left(P-\lambda_{i} I\right)^{j}} \frac{j!}{}(j)_{\lambda_{i}^{n-j}}^{Z_{i}}(V M-I) \tag{8.2}
\end{equation*}
$$

The value of $\delta$ is found by the following procedure : Since $\left\{\mathscr{\xi}_{n}\right\}$ and $\left\{\tilde{\xi}_{n}\right\}$ are assumed to have the same limit experiment, and since these limit experiments must be given by $Z_{0}$ and $Z_{0} V$, respectively, it follows that there exists at least
one $M$ so that $Z_{o} V M=Z_{0}$. Let $M M_{0}$ be the set of such M's $; M_{0} \subseteq M_{r, r}$. Define $M_{1} \subseteq M_{0}$ as the set of $M \in M_{0}$ so that $Z_{i} V M=z_{i}$ for all $Z_{i}$ corresponding to eigenvalues with modulus $\rho_{1}$. In the same manner define $M_{2} \subseteq M_{1}$ so that $z_{i} V M=z_{i}$ for all $Z_{i}$ corresponding to eigenvalues with modulus $\rho_{2}$. Continuing, one gets a chain

$$
m_{r, r} \supseteq M_{0} \supseteq M_{1} \supseteq \cdots
$$

If $M_{t} \neq \varnothing$, then it is clear that $\epsilon_{n}$ equals 0 from some $n$ on . This is the case if our lumping is sufficient.
If $M_{t}=\varnothing$, let $j$ be the smallest number $(1 \leq j \leq t)$ so that

$$
M M_{j}=\varnothing
$$

We claim that $\delta=\rho_{j}$.
For each $M \in M_{j-1}$, let $\nu_{M}$ be the maximum index of eigenvalues $\lambda_{i}$ with modulus $\rho_{j}$ corresponding to $Z_{i}$ 's with $Z_{i} V M \neq Z_{i}$. We claim that

$$
\nu=\min \left\{\nu_{\mathrm{M}}: M \in \mathbb{M}_{j-1}\right\}
$$

By inserting an $M^{\prime} \in M_{j-1}$ with $v_{M}=v$, we get

$$
\epsilon_{n}=\inf _{M}\left\|P^{n^{n}} V M-P^{n^{n}}\right\| \leq\left\|P^{n} V M^{\prime}-P^{n^{n}}\right\|
$$

which by (8.2) is seen to be $\leq C_{n}^{\nu-1} \delta^{n}$ for some number $C>0$, independent of $n$.

To prove the left inequality of the theorem, assume that a is the largest label for eigenvalues with modulus $\delta=\rho_{j}$. We multiply each side of (8.2) by $\sum_{i=1}^{a} Z_{i}$ to get (for any $M$ )

$$
\begin{align*}
& \left\|\sum_{i=1}^{a} Z_{i}\right\|\left\|P^{n} V M-P^{n}\right\| \\
& \geq\left\|Z_{0}(V M-I)+\sum_{i=1}^{a} \sum_{j=0}^{m_{i}-\neq\left(P-\lambda_{i} I\right)^{j}} n^{j!}(j)_{\lambda_{i}^{n}-j} Z_{i}(V M-I)\right\| \tag{8.3}
\end{align*}
$$

Applying to lemma 8.3, there exists a number $\varphi>0$, independent of $M$, so that

$$
\left\|Z_{i}(V M-I)\right\| \geq \varphi \quad \text { for at least one } i=0,1, \ldots, k .
$$

Considering the terms occurring on the right hand side of (8.3) we thus conclude that

$$
\left\|P^{n} V M-P^{n}\right\| \geq c n^{\nu-1} \delta^{n}
$$

where $c>0$ is independent of $M$.
This completes the proof.

We are now in position to prove our main result.

Theorem 8.5.
Let $\delta$ be as in theorem 8.4.
Then $\sqrt[n]{\eta_{n}} \rightarrow \delta$ as $n \rightarrow \infty$.

Proof: Let $\gamma_{n}=\gamma\left(A, \bigotimes_{n}\right)$, where $\gamma$ is defined in section 7 . It suffices to prove that there exists $K>0$ and $w \geq 0$ such that $\gamma_{n} \leq K_{n}{ }_{\delta}{ }^{n}$, $k$ since by theorem 7.3 and 7.4 we have $\epsilon_{n} \leq \eta_{n} \leq 2 \gamma_{n}$.

Now $\quad \gamma\left(A, G_{n}\right)=\inf _{W}\left\|P^{n} V W-P^{n}\right\|$
where infimum is taken over all A-Markov matrices $W$.

We return to the proof of theorem 8.4 . Let $b$ be the largest label of eigenvalues with modules $\rho_{j-1}$. Then there exists $M \in M_{r, r}$ such that

$$
\begin{equation*}
Z_{o} V M=Z_{0}, \ldots, Z_{b} V M=Z_{b} \tag{8.5}
\end{equation*}
$$

$Z_{0}$ is of the form $B$ given in (4.3), where the matrices $A_{1}, \ldots, A_{V}$ have strictly positive entries.

Hence from some $n$ on (which of course may be $n=1$ )

$$
\begin{equation*}
Q_{n}=Z_{0}+\sum_{i=1}^{b} \sum_{j=0}^{m_{i}-1} \frac{\left(P-\lambda_{i} I\right)^{j}}{j!} n(j)_{\lambda_{i}} n-j_{Z_{i}} \tag{8.6}
\end{equation*}
$$

will have no negative entries and hence define an experiment.
By (8.5) there exists $M$ so that $Q_{n} V M=Q_{n}$ for $n=1,2, \ldots$. Hence $A$ is sufficient for $Q_{n}$, so by theorem $7.1 M$ may for each $n$ be replaced by an A-Markov matrix $W_{n}$ such that $Q_{n} V W_{n}=Q_{n}$. But then, by (8.2), (8.4) and (8.6), $Y_{n} \leq\left\|P^{n} V w_{n}-P^{n}\right\| \leq C_{n} w^{n}{ }^{n}$ and the proof is complete.

We end up this section by a couple of examples.
Example 8.6.

$$
\text { Let } P=\left(\begin{array}{lll}
\frac{1}{2} & \frac{1}{2} & 0 \\
1 & 0 & 0 \\
\frac{1}{4} & 0 & \frac{3}{4}
\end{array}\right)
$$

The eigenvalues of $P$ are $1, \frac{3}{4},-\frac{1}{2}$.
Let $A=\{\{1,2\},\{3\}\}$. The chain is A-Iumpable.

We have $P^{n}=Z_{0}+\left(\frac{3}{4}\right)^{n} Z_{1}+\left(-\frac{1}{2}\right)^{n} Z_{2}$ where

$$
\begin{aligned}
& Z_{0}=\left(\begin{array}{ccc}
\frac{2}{3} & \frac{1}{3} & 0 \\
\frac{2}{3} & \frac{1}{3} & 0 \\
\frac{2}{3} & \frac{1}{3} & 0
\end{array}\right) \\
& Z_{1}=\left(\begin{array}{ccc}
0 & 0 & 0 \\
0 & 0 & 0 \\
-\frac{3}{5} & -\frac{2}{5} & 1
\end{array}\right) \\
& z_{2}=\left(\begin{array}{rrr}
\frac{1}{3} & -\frac{1}{3} & 0 \\
-\frac{2}{3} & \frac{2}{3} & 0 \\
-\frac{1}{15} & \frac{1}{15} & 0
\end{array}\right)
\end{aligned}
$$

A simple computation shows that $M_{0}$ consists of all matrices of the form

$$
\left(\begin{array}{ccc}
\frac{2}{3} & \frac{1}{3} & 0 \\
m_{1} & m_{2} & m_{3}
\end{array}\right)
$$

$m_{1}, m_{2}, m_{3} \geq 0 ; m_{1}+m_{2}+m_{3}=1$.
With $M \in M_{0}$ we get
$Z_{1}(V M-I)=\left(\begin{array}{ccc}0 & 0 & 0 \\ 0 & 0 & 0 \\ m_{1}-\frac{1}{15} & m_{2}+\frac{1}{15} & m_{3}-1\end{array}\right)$
and no allowable choice of $m_{1}, m_{2}, m_{3}$ will give $Z_{1}(V M-I)=0$.

Hence, $M_{1}=\varnothing$, and by the preceding theorems,

$$
\sqrt[n]{\epsilon_{n}} \text { and } \sqrt[n]{\eta_{n}} \rightarrow \frac{3}{4}
$$

i.e. we have convergence to the root of $P$. That this will not happen in general, is seen from the next example.

## Example 8.7.

Let $P=\left(\begin{array}{lll}\frac{7}{16} & \frac{1}{4} & \frac{5}{16} \\ \frac{1}{2} & 0 & \frac{1}{2} \\ \frac{5}{16} & \frac{1}{4} & \frac{7}{16}\end{array}\right)$

The eigenvalues of $P$ are $1,-\frac{1}{4}, \frac{1}{8}$. Let $A=\{\{1,3\},\{2\}\}$.

Now $P^{n}=Z_{0}+\left(-\frac{1}{4}\right)^{n} Z_{1}+\left(\frac{1}{8}\right)^{n} Z_{2}$
with

$$
\begin{aligned}
& Z_{0}=\left(\begin{array}{ccc}
\frac{2}{5} & \frac{1}{5} & \frac{2}{5} \\
\frac{2}{5} & \frac{1}{5} & \frac{2}{5} \\
\frac{2}{5} & \frac{1}{5} & \frac{2}{5}
\end{array}\right), \quad z_{1}=\left(\begin{array}{ccc}
\frac{1}{10}-\frac{2}{10} & \frac{1}{10} \\
-\frac{2}{5} & \frac{4}{5} & -\frac{2}{5} \\
\frac{1}{10}-\frac{2}{10} & \frac{1}{10}
\end{array}\right) \\
& Z_{2}=\left(\begin{array}{ccc}
\frac{1}{2} & 0 & -\frac{1}{2} \\
0 & 0 & 0 \\
-\frac{1}{2} & 0 & \frac{1}{2}
\end{array}\right)
\end{aligned}
$$

A computation shows that $M_{1} \neq \varnothing$, more precisely, $M_{1}$ has one element

$$
M=\left(\begin{array}{lll}
\frac{1}{2} & 0 & \frac{1}{2} \\
0 & 1 & 0
\end{array}\right)
$$

However, $M_{2}=\varnothing$ since $Z_{2} V=0$. Hence $\sqrt[n]{\epsilon_{n}}$ and $\sqrt[n]{\eta_{n}} \rightarrow \frac{1}{8}$, whence the root of $P$ equals $\frac{1}{4}$.

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