# On the Preservation of Local Asymptotic Normality Under Information Loss* 

## by

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## 1. Introduction

In the present paper we consider a situation where there are unobservable random variables $X_{j, n} ; j=1, \cdots, k_{n}, n=1,2, \ldots$ and where what is actually seen are other variables $\mathrm{Y}_{\mathrm{j}, \mathrm{n}}$ that are less informative than the $\mathrm{X}_{\mathrm{j}, \mathrm{n}}$. For instance the $\mathrm{Y}_{\mathrm{j}, \mathrm{n}}$ may be functions of the $X_{j, n}$.

It can then readily happen that the family of probability measures that governs the behavior of the $X_{j, n}$ is simple and easily studied but that the corresponding family for the $Y_{j, n}$ is more untractable. We shall show that if the $\left\{X_{j, n}\right\}$ satisfy certain conditions, such as the LAN conditions, then the $Y_{j, n}$ will also satisfy the same requirements. This means that certain methods of estimation and testing can be carried out with the $Y_{j, n}$ with some assurance of success without having to verify that their distributions satisfy the necessary requirements. Actual computation of the estimates may still be a complex task, but not nearly as difficult as the analytical effort needed in the verification of assumptions.

The results are proved under suitable negligibility conditions imposed on the $X_{j, n}$ and mostly for the case where the $X_{j, n}$ are independent for each fixed $n$. However, we point out that a similar phenomenon can also be expected for certain nonindependent double arrays $\left\{\mathrm{X}_{\mathrm{j}, \mathrm{n}}\right\}$.

The problem was brought to our attention by several special examples. One of them is in a neurophysiological study where the underlying model involved independent variables $U_{v}, X_{1, v}, \ldots$ The variables $U_{v}$ had a binomial distribution $B(m, p)$ and the $\mathrm{X}_{\mathrm{j}, v}$ were positive i.i.d. variables whose common distribution depended smoothly on a one-dimensional parameter $\theta$. Standard textbook results on the asymptotic behavior of the system were readily applicable to the variables $\left\{\mathrm{U}_{\mathrm{v}}, \mathrm{X}_{1, v}, \ldots\right\}$ but what could be observed was only the sums

$$
Y_{v}=\sum_{k}\left\{X_{k, v} ; 0 \leq k \leq U_{v}\right\} \quad v=1,2, \ldots, n .
$$

Another example comes from what is called the use of counting processes for life history data (see [1], [3], [15]). One has a large number $n$ of different individuals behaving independently of each other. During the observation period [0,L] certain events $A_{i} ; i=1,2, \ldots, I$ can happen to them. At time $t$ the instantaneous intensity for the occurrence of event $A_{i}$ for individual $k$ is a variable $\lambda_{k, i}(t)$. As functions of $t$ these $\lambda_{k, i}(t)$ are predictable random processes in the sense that $\lambda_{k, i}(t)$ depends only on the history of individual $k$ in the half open interval $[0, t)$. In such problems there is a well developed theory applicable to the situation where each individual is monitored continuously throughout the entire period [0,L]. Its elaboration, reviewed in [3], relies materially on the fact that likelihood ratios can be written in a simple explicit form.

The question arises whether the theory can be made applicable to a situation where each individual is seen only at a few isolated times, say $0 \leq t_{0}<t_{1}<\cdots<t_{\delta} \leq L$. Even in a purely Markovian system, assumptions made on the intensities $\lambda_{\mathrm{k}, \mathrm{i}}$ do not translate into simple statements on the likelihood ratios for isolated observation times. The fact that our theorems give some information on the behavior to be expected of likelihood ratios allows one to direct numerical computations in an effective manner.

Our main results depend very strongly on the negligibility assumptions placed on the $\mathrm{X}_{\mathrm{j}, \mathrm{n}}$. It is a trivial fact that the results are not valid without some restrictions, but, except for a modification of a result of Robert Davies [8] given in Section 6, we do not know whether extensions to other situations are possible.

Section 2 below gives the essential notations and assumptions with a statement of the problem in a formal mathematical framework. Section 3 shows that contiguity is preserved under loss of information. It also contains related inequalities that will be used in the following sections. Section 4 is about the preservation of asymptotic normality for what we call bounded infinitesimal arrays. These are double arrays of independent experiments where the total information (defined in a suitable manner) remains bounded and where each individual experiment contributes an asymptotically negligible amount of information to the total.

Asymptotic normality has been used mostly in connection with the so-called LAN conditions (see [10], [20], [25], [34]). These involve not only the local asymptotic normality that gave them their name but also linear relations between vector parameters and logarithms of likelihood ratios. It is shown in Section 5 that such linear relations are also preserved. The technique of proof involves the use of quadratic forms that "control" the asymptotic behavior of likelihood ratios. According to Davies, [8], such '"quadratic control" maybe preserved under certain information losses that do not actually preserve the local Gaussian behavior. It is the main relation needed for the improvement of auxiliary estimates, described later in the same section.

Section 5 ends in a description of a method of construction of estimates. One starts with a good auxiliary estimate and adds to it a correction calculated from likelihood ratios. This gives asymptotically efficient estimates. The general method has been described elsewhere by one of the present authors. It is pointed out that, for our infinitesimal arrays, the method admits a number of variants that are often easily applicable in practice.

An application to the neurophysiological problem that motivated us is described in some detail in Example 1, Section 9.

Section 6 departs from the general theme. It does not use independence or negligibility assumptions. What it shows is that joint asymptotic normality of logarithms of likelihood and of estimates available from the restricted information together with a condition that these estimates be "distinguished" leads to preservation of the LAN
conditions. For that LAN situation the result is an improvement of a result of R . Davies who dealt in [8] and [9] with preservation of what we call "quadratic control".

Section 7 is an aside on the independent identically distributed case. It shows that differentiability in quadratic mean (for square roots of likelihood ratios) is preserved under information loss. This is not surprising.

Section 8 sketches a possible extension of the results of Sections 4 and 5 to cases where the observations are not independent but where the information is still acquired by infinitesimal amounts. It is clearly applicable to certain Markov models where the information is lost by observing only at isolated times instead of continuously. Some of the difficulties that arise in the non-Markovian case are pointed out.

Section 9 contains several examples intended as illustrations. The first is the already mentioned neurophysiological problem that started our investigation. The second is intended to show in an i.i.d. situation that one should not expect that good behavior of maximum likelihood estimates would be preserved. It is also a warning against iterative procedures often used to compute approximations to the maximum likelihood estimates.

Example 3 explains how to deal with homogeneous Markov processes. It also contains a warning about loss of identifiability for the instantaneous transition intensities.

Example 4 shows that asymptotic normality may be preserved even though the rates of convergence of estimates are altered. This is to emphasize the caveats of Section 2.

Example 5 discusses the method of moments used in estimation.
Example 6 shows that the techniques proposed here apply to the loss of information incurred by grouping data.

Finally, Example 7 allows one to pass from ordinary data to censored data. Here again the identifiability problems may interfere. Otherwise solutions are reasonably simple. Each section, except Section 2, begins by a short description of its content. The reader who is not interested in the techniques of proofs may skip them. It is possible to read the examples first, referring as needed to the statements of the theorems.

## Acknowledgments

We are indebted to referees who pointed out a gap in our discussion of the method of moments and the possibility of applying the results of Section 8 to the aggregated Markov Chains situation.

## 2. Notations and assumptions

The main results of the present paper refer to a situation describable as follows: For each integer $n$, let $\Theta_{n}$ be a parameter set. For each pair of integers (j,n); $j=1,2, \ldots, k_{n}, n=1,2, \ldots$, let $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ be an experiment indexed by $\Theta_{\mathrm{n}}$, that is a family $\mathbf{E}_{\mathrm{j}, \mathrm{n}}=\left\{\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ of probability measures on a $\sigma$-field $\mathbf{A}_{\mathrm{j}, \mathrm{n}}$. Let $\overline{\mathbf{E}}_{\mathrm{n}}$ be the direct product of the $\mathbf{E}_{j, n}$, that is, $\overline{\mathbf{E}}_{\mathrm{n}}$ is the experiment where one observes independent random elements $X_{j, n}$ in such a manner that, when $\theta$ is true, the distribution of $X_{j, n}$ is given by $\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta)$. The joint distribution is given by the product measure

$$
P_{n}(\theta)=\prod_{j} p_{j, n}(\theta)
$$

Consider another double array of experiments $F_{j, n}=\left\{q_{j, n}(\theta) ; \theta \in \Theta_{n}\right\}$ and the corresponding direct product $\overline{\mathbf{F}}_{\mathrm{n}}$. Assume that each $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$ is less informative than the corresponding $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ in the sense of Blackwell [6], that is for every decision problem any risk function possible on $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$ is also possible in $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$. This can happen for instance if $q_{j, n}(\theta)$ is the restriction of $p_{j, n}(\theta)$ to a sub- $\sigma$-field $\mathbf{B}_{j, n}$ of $\mathbf{A}_{j, n}$.

In such a situation one may inquire whether certain asymptotic properties of the sequence $\bar{E}_{n}$ remain valid for the weaker sequence $\bar{F}_{n}$. For instance, let $s_{n}$ and $t_{n}$ be two points of $\Theta_{n}$. Consider the product measures $P_{n}\left(s_{n}\right)$ and $P_{n}\left(t_{n}\right)$ and the corresponding products $\mathrm{Q}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)$ and $\mathrm{Q}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)$ relative to the experiments $\overline{\mathrm{F}}_{\mathrm{n}}$.
(a) If $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)\right\}$ and $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)\right\}$ are contiguous sequences, is the same true of $\left\{\mathrm{Q}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)\right\}$ and $\left\{Q_{n}\left(\mathrm{t}_{\mathrm{n}}\right)\right\}$ ?
(b) If the pair $\left[\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right), \mathrm{P}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)\right]$ is an asymptotically Gaussian experiment, is the same true of $\left[\mathrm{Q}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right), \mathrm{Q}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)\right]$ ?
(c) If $\Theta_{\mathrm{n}}$ is a Euclidean space $\mathrm{R}^{\mathbf{k}}$ and the $\overline{\mathbf{E}}_{\mathrm{n}}$ satisfy the LAN conditions, is the same true of the $\overline{\mathbf{F}}_{\mathrm{n}}$ ?

The answer to question (a) is easily seen to be positive. This will be proved in Section 3 below. The answer to (b) and (c) is more complex. We shall show that it is also positive under a "negligibility" requirement imposed on the components $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$. However the positivity of the answer must be qualified in an important manner. To indicate the qualifications, note that for such product experiments there are some natural metrics definable as follows. Let

$$
h_{j, n}^{2}(s, t)=\frac{1}{2} \int\left[\sqrt{\mathrm{dp}_{j, n}(s)}-\sqrt{\mathrm{dp}_{j, n}(t)}\right]^{2}
$$

be the square Hellinger distance for $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$. Let $\mathrm{H}_{\mathrm{n}}^{2}(\mathrm{~s}, \mathrm{t})=\sum_{\mathrm{j}} \mathrm{h}_{\mathrm{j}, \mathrm{n}}^{2}(\mathrm{~s}, \mathrm{t})$. Using the alternate measures $\mathrm{q}_{j, \mathrm{n}}(\mathrm{s})$ in a similar manner, define from $\overline{\mathrm{F}}_{\mathrm{n}}$ a distance $\mathrm{K}_{\mathrm{n}}$ analogous to $\mathrm{H}_{\mathrm{n}}$.

For the experiments $\left\{\mathrm{P}_{\mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ the metric $\mathrm{H}_{\mathrm{n}}$ is a very natural one. The local asymptotic theories are often carried out on neighborhoods of the type

$$
V_{n}(\theta, b)=\left\{t ; t \in \Theta_{n}, H_{n}(t, \theta) \leq b\right\}
$$

for a presumed true value $\theta$. This is what one does, more or less, for the LAN assumptions. This is also what one does in the i.i.d. case for many nonparametric studies. (See [31], for instance.)

The corresponding neighborhoods for the $\bar{F}_{\mathrm{n}}$ would be

$$
W_{n}(\theta, b)=\left\{t ; t \in \Theta_{n}, K_{n}(t, \theta) \leq b\right\}
$$

Note that since $K_{n} \leq H_{n}$, one has $V_{n}(\theta, b) \subset W_{n}(\theta, b)$ but that $W_{n}(\theta, b)$ may be immensely larger than $\mathrm{V}_{\mathrm{n}}(\theta, b)$ (in terms of the distance $H_{n}$ ).

What we shall show is that certain properties, such as asymptotic normality, are inherited by the $\left\{\mathrm{Q}_{\mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ on the sets of the type $\mathrm{V}_{\mathrm{n}}$. Nothing can be said in general about the larger $W_{n}$, although Section 7 does contain a result of that nature.

For the same reason, existence of estimates that converge at a given speed on the $\left\{\mathrm{P}_{\mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ is not a property that transfers to the weaker $\left\{\mathrm{Q}_{\mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$, since, for instance, the measures $\mathrm{Q}_{\mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}$ might be all equal to a single one, $\mathrm{Q}_{\mathrm{n}}\left(\theta_{0}\right)$. Section 7 again, contains some further elaborations on this point. The reader may have noted that in the description of the experiments $\bar{E}_{n}=\left\{P_{n}(\theta) ; \theta \in \Theta_{n}\right\}$ the index set is allowed to vary with n . This again because of the local nature of the results. It is sometimes possible to stabilize the sets $\Theta_{\mathrm{n}}$ by various transformations, such as using $\theta=\sqrt{n}\left(t-t_{0}\right)$ or more elaborate matrix multiplication. Then the $\Theta_{n}$ may be replaced by a fixed set $\Theta$ independent of $n$. However such stabilizations are essentially irrelevant to our main purposes and they would make matters more complex since the stabilizing transformation needed for $\overline{\mathbf{F}}_{\mathrm{n}}$ may be different from that needed for $\overline{\mathbf{E}}_{\mathrm{n}}$.

## 3. Preservation of contiguity and some other inequalities

The purpose of this section is to show that contiguity is always preserved under loss of information. This can readily be established directly, but we use a more specific approach yielding inequalities on likelihood ratios. They will be used in the subsequent sections.

One first easy remark is as follows.
PROPOSITION 1. Let $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)\right\}$ and $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)\right.$ be contiguous. Then the sequences $\left\{\mathrm{Q}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)\right\}$ and $\left\{\mathrm{Q}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)\right\}$ are also contiguous.

This is a special case of a more general lemma (see [27], Chapter 6, Section 2).
LEMMA 1. Let $\mathrm{E}_{\mathrm{n}}=\left\{\mathrm{P}_{\mathrm{n}, 0}, \mathrm{P}_{\mathrm{n}, 1}\right\}$ and $\mathrm{F}_{\mathrm{n}}=\left\{\mathrm{Q}_{\mathrm{n}, 0}, \mathrm{Q}_{\mathrm{n}, 1}\right\}$ be two binary experiments. Assume that 1) the deficiency $\delta\left(\mathbf{E}_{\mathrm{n}}, \mathrm{F}_{\mathrm{n}}\right)$ of $\mathbf{E}_{\mathrm{n}}$ with respect to $\mathbf{F}_{\mathrm{n}}$ tends to zero and 2) the sequences $\left\{\mathrm{P}_{\mathrm{n}, 0}\right\},\left\{\mathrm{P}_{\mathrm{n}, 1}\right\}$ are contiguous. Then the sequences $\left\{\mathrm{Q}_{\mathrm{n}, 0}\right\},\left\{\mathrm{Q}_{\mathrm{n}, 1}\right\}$ are contiguous.

This lemma itself is an immediate consequence of an inequality that will be used in Section 4. To state it, let $S$ be a set of probability measures on a $\sigma$-field $\mathbf{A}$. Let $\lambda$ be another positive finite measure on $\mathbf{A}$. Let us say that the set $S$ is $(\alpha, \mathrm{c} \lambda)$ limited if $\|P-P \wedge c \lambda\| \leq \alpha$ for every $P \in S$.

Here the norm is the total variation norm and the measure $P \wedge c \lambda$ is defined in the following manner. One takes the Radon-Nikodym density $\frac{d P}{d \lambda}$ of the part of $P$ that is dominated by $\lambda$. One takes for $P_{\wedge} \subset \lambda$ the measure whose density with respect to $\lambda$ is the minimum of $\frac{d P}{d \lambda}$ and $c$.

LEMMA 2. Let $\mathrm{E}=\left\{\mathrm{P}_{\theta} ; \theta \in \Theta\right\}$ and $\mathrm{F}=\left\{\mathrm{Q}_{\theta} ; \theta \in \Theta\right\}$ be two experiments with the same set of indices $\Theta$. Assume
i) $\left\{\mathrm{P}_{\theta} ; \theta \in \Theta\right\}$ is $(\alpha, \mathrm{c} \lambda)$ limited,
ii) the deficiency $\delta(\mathbf{E}, \mathbf{F})$ does not exceed $\varepsilon / 2$.

Then there is a positive linear transformation T that sends probability measures into probability measures and is such that
iii) $\left\|Q_{\theta}-\mathrm{TP}_{\theta}\right\| \leq \varepsilon$ for all $\theta$,
iv) if $\mu=T \lambda$ then $F$ is $(\alpha+\varepsilon, c \mu)$ limited.

PROOF. The existence of $T$ is part of a general theorem of [21]. Once this is established note that

$$
T\left(\mathrm{P}_{\theta} \wedge \mathrm{c} \lambda\right) \leq\left(\mathrm{TP}_{\theta}\right) \wedge(\mathrm{Tc} \lambda)=\left(\mathrm{TP}_{\theta}\right) \wedge(\mathrm{c} \mu) .
$$

Also

$$
\left\|\mathrm{T}\left(\mathrm{P}_{\theta}-\mathrm{P}_{\theta} \wedge \mathrm{c} \lambda\right)\right\| \leq\left\|\mathrm{P}_{\theta}-\mathrm{P}_{\theta} \wedge \mathrm{c} \lambda\right\| \leq \alpha .
$$

Therefore

$$
\left\|\mathrm{TP}_{\theta}-\left(\mathrm{TP}_{\theta}\right) \wedge(\mathrm{c} \mu)\right\| \leq \alpha .
$$

Now write $\mathrm{Q}_{\theta}-\mathrm{TP}_{\theta}=\mathrm{D}=\mathrm{D}^{+}-\mathrm{D}^{-}$. Then $\mathrm{Q}_{\theta} \leq \mathrm{TP}_{\theta}+\mathrm{D}^{+}$and $\mathrm{Q}_{\theta} \wedge \mathrm{c} \mu \geq$ $\left(\mathrm{TP}_{\theta}-\mathrm{D}^{-}\right) \wedge \mathrm{c} \mu$. Therefore

$$
\left\|\mathrm{Q}_{\theta}-\mathrm{Q}_{\theta} \wedge \mathrm{c} \mu\right\| \leq\left\|\mathrm{TP}_{\theta}-\left(\mathrm{TP}_{\theta}\right) \wedge(\mathrm{c} \mu)\right\|+\|\mathrm{D}\| .
$$

This gives the desired result.
Applying this to the case where $\delta(\mathbf{E}, \mathbf{F})=0$, one sees that truncating the densities of the measures $Q_{\theta}$ results in a modification of these measures that is smaller than the modification made by the same truncation on the $\mathrm{P}_{\boldsymbol{\theta}}$.

One should expect that some other inequalities would also be preserved. Many arguments involving maximum likelihood estimates use bounds on expressions of the type

$$
\phi(x, v)=\sup _{t \in V}\left[\frac{f(x, t)}{f\left(x, \theta_{0}\right)}-1\right]^{+}
$$

for densities $f(x, t)=\frac{d P_{t}}{d \lambda}$. A bound on an integral $\int \phi(x, v) \mathrm{dP}_{\theta_{0}}$ will carry over to the similar expression defined on the weaker $F$. However, some precautions must be taken. See the examples of Section 9.

## 4. Preservation of asymptotic normality.

In this section we consider experiments that are weakly asymptotically normal in the sense of the Gaussian approximability described in Definition 2 below. It is first shown that this property is equivalent to approximability of distributions of log likelihood by multivariate Gaussian distributions at least whenever Hellinger affinities remain bounded away from zero.

This being done we consider independent observations where the individual experiments form what we call bounded infinitesimal arrays (Def. 3 below). For these it is shown that, when loss of information occurs on each individual component, asymptotic normality is preserved.

The approximations by Gaussian experiments involve quadratic forms that "control" the behavior of likelihood ratios (Def. 4 below). The forms that control the weaker experiments are smaller than the initial ones in the sense that the difference is positive semi-definite.

Explicitly, the situation can be described as follows.
An experiment $\mathbf{G}=\left\{G_{\theta} ; \theta \in \Theta\right\}$ is called Gaussian if it satisfies the following two conditions:
i) for any pair (s,t) of elements of $\Theta$ the measures $G_{s}$ and $G_{t}$ are mutually absolutely continuous
ii) Let $\Lambda(\mathrm{t}, \mathrm{s})=\log \frac{\mathrm{dG}_{\mathrm{t}}}{\mathrm{dG}_{\mathrm{s}}}$. The stochastic process $\mathrm{t} \rightarrow \Lambda(\mathrm{t}, \mathrm{s}), \mathrm{t} \in \Theta$ is a Gaussian process for the distributions induced by $\mathrm{G}_{\mathrm{s}}$. Here s is an arbitrary point of $\Theta$.

It is easily checked that, under (i), if the condition (ii) is satisfied for some choice of $s$, it is satisfied for all.

DEFINITION 1. Let $\bar{E}_{n}=\left\{P_{n}(\theta) ; \theta \in \Theta_{n}\right\}$ be a sequence of experiments. One says that $\overline{\mathbf{E}}_{\mathrm{n}}$ admits strong Gaussian approximations if there are Gaussian experiments $\mathbf{G}_{\mathrm{n}}=\left\{\mathrm{G}_{\mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ such that the distance $\Delta\left(\bar{E}_{\mathrm{n}}, \mathbf{G}_{\mathrm{n}}\right)$ tends to zero.
(The distance $\Delta$ is the one defined in [21] or [27].)
DEFINITION 2. One says that $\overline{\mathbf{E}}_{\mathrm{n}}$ admits weak Gaussian approximations if there are Gaussian $G_{n}$ such that for any subsets $S_{n} \subset \Theta_{n}$ that have a cardinality bounded independently of n , the distance between $\overline{\mathbf{E}}_{\mathrm{n}}$ and $\mathbf{G}_{\mathrm{n}}$ restricted to $\mathrm{S}_{\mathrm{n}}$ tends to zero.

The condition of Definition 2 is similar to a relation often used in asymptotics under circumstances where $\Theta_{\mathrm{n}}$ is a fixed set $\Theta$ independent of n . Then one can speak of weak convergence of $\overline{\mathbf{E}}_{\mathrm{n}}$ to a Gaussian limit $\mathbf{G}=\left\{\mathrm{G}_{\theta} ; \theta \in \Theta\right\}$. This is convergence for the distance $\Delta$ for $\overline{\mathbf{E}}_{\mathrm{n}}$ restricted to fixed finite subsets $\mathrm{S} \subset \Theta$. Note that the weak approximability of Definition 2 is considerably stronger than what would be obtained
from weak convergence to a limit G. Our sets $S_{n} \subset \Theta$ have bounded cardinality, just as a fixed finite S , but they are allowed to vary arbitrarily with n . Thus our weak approximability is uniform on all sets of a given finite cardinality. It would be inconvenient here to use weak convergence to limits since our successive $\Theta_{\mathrm{n}}$ need not be related at all.

Definitions 1 or 2 do not put any additional restrictions on the Gaussian $\mathbf{G}_{\mathrm{n}}$. Under some restrictions the weak approximability can be checked on the behavior of distributions of likelihood ratios as follows.

PROPOSITION 2. Let $\overline{\mathbf{E}}_{\mathrm{n}}=\left\{\mathrm{P}_{\mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ be a sequence of experiments. Assume that there is some $\varepsilon>0$ such that

$$
\inf _{\mathrm{n}, \mathrm{~s}, \mathrm{t}} \int \sqrt{\mathrm{dP}_{\mathrm{n}}(\mathrm{~s}) \mathrm{dP}_{\mathrm{n}}(\mathrm{t})} \geq \varepsilon .
$$

Then the $\overline{\mathbf{E}}_{\mathrm{n}}$ admit weak Gaussian approximation if and only if for every fuxed k and for every subset $\left\{\mathrm{s}_{0, \mathrm{n}}, \mathrm{s}_{1, \mathrm{n}}, \cdots, \mathrm{s}_{\mathrm{k}, \mathrm{n}}\right\}$ of elements of $\Theta_{\mathrm{n}}$ the joint distributions $\mathrm{F}_{\mathrm{n}}$ under $\mathrm{s}_{0, \mathrm{n}}$ of the logarithms

$$
\log \frac{d P_{n}\left(s_{j, n}\right)}{d P_{n}\left(s_{0, n}\right)} ; j=1,2, \ldots k
$$

are approximable by Gaussian distributions in the sense that the Lévy (or Prohorov) distance between $\mathrm{F}_{\mathrm{n}}$ and a suitable k -variate Gaussian distribution tends to zero as $\mathrm{n} \rightarrow \infty$.

NOTE. By $\frac{\mathrm{dP}(\mathrm{t})}{\mathrm{dP}(\mathrm{s})}$ is meant the density with respect to $\mathrm{P}(\mathrm{s})$ of the part of $\mathrm{P}(\mathrm{t})$ dominated by $\mathrm{P}(\mathrm{s})$.

PROOF. Consider first two point sets $\left\{\mathrm{s}_{0, \mathrm{n}}, \mathrm{s}_{1, \mathrm{n}}\right\}$ and the corresponding binary experiments $E_{n}^{\prime}=\left\{P_{n}\left(s_{0, n}\right), P_{n}\left(s_{1, n}\right)\right\}$. Taking a subsequence if necessary one can assume convergence to a limit $\mathbf{E}^{\prime}=\left\{\mathrm{R}_{0}, \mathrm{R}_{1}\right\}$. The lower bound imposed on affinities implies that the $\mathrm{R}_{\mathrm{i}}, \mathrm{i}=0,1$ cannot be disjoint. Let $\mathrm{X}_{\mathrm{n}}$ be distributed as

$$
\log \frac{\mathrm{dP}_{\mathrm{n}}\left(\mathrm{~s}_{1, \mathrm{n}}\right)}{\mathrm{dP}_{\mathrm{n}}\left(\mathrm{~s}_{0, \mathrm{n}}\right)}
$$

under $P_{n}\left(s_{0, n}\right)$. There are numbers $a<b$ such that $\operatorname{Pr}\left[a \leq X_{n} \leq b\right]$ remains bounded away from zero. Let $Y_{n}$ be a normal variable with a distribution $N\left(\mu_{n}, \sigma_{n}^{2}\right)$ approximating that of $X_{n}$. Then $\operatorname{Pr}\left[a-1 \leq Y_{n} \leq b+1\right]$ stays bounded away from zero. It follows that $\sigma_{\mathrm{n}}$ must remain bounded away from $+\infty$.

If so, $\left|\mu_{n}\right|$ must also remain bounded. This means that cluster points of the sequence of distributions $L\left(e^{X_{n}}\right)$ cannot have masses at zero. Thus $R_{1}$ dominates $R_{0}$. This implies that the sequences $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{0, n}\right)\right\}$ and $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{1, n}\right)\right\}$ must be contiguous. It then
follows by a standard argument that the conditions given are sufficient to imply the weak Gaussian approximability for an arbitrary k. The result in the opposite direction is also the consequence of a standard argument for which see for instance [22] page 14.

Keeping in mind this result, let us return to the case of product experiments $\overline{\mathbf{E}}_{\mathrm{n}}=\prod_{\mathrm{j}} \mathbf{E}_{\mathrm{j}, \mathrm{n}}$ described in Section 2.

DEFINITION 3. The double array $\left\{\mathbf{E}_{\mathrm{j}, \mathrm{n}}\right\}$ with $\mathbf{E}_{\mathrm{j}, \mathrm{n}}=\left\{\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta), \theta \in \Theta_{\mathrm{n}}\right\}$ will be called bounded and infinitesimal if it satisfies the following two requirements for all pairs ( $\mathrm{s}_{\mathrm{n}}, \mathrm{t}_{\mathrm{n}}$ ) extracted from $\Theta_{\mathrm{n}}$

$$
\begin{align*}
& \sup _{n} \sum_{j} h_{j}^{2}\left(s_{n}, t_{n}\right)<\infty  \tag{A}\\
& \lim _{n} \sup _{j} h_{j}^{2}\left(s_{n}, t_{n}\right)=0
\end{align*}
$$

(where $h_{j}$ is the Hellinger distance defined in Section 2).
Note that when (B) holds condition (A) is equivalent to the affinity restriction used in Proposition 2.

For such bounded infinitesimal arrays weak Gaussian approximability can be shown to be equivalent to any one of a large number of other properties. The one that will be most convenient here is as follows.

PROPOSITION 3. Let $\left\{\mathbf{E}_{\mathrm{j}, \mathrm{n}}\right\}$ be a bounded infinitesimal array. Then the product $\overline{\mathbf{E}}_{\mathrm{n}}$ admits weak Gaussian approximations if and only if for every choice of pairs ( $\mathrm{s}_{\mathrm{n}}, \mathrm{t}_{\mathrm{n}}$ ) of elements of $\Theta_{\mathrm{n}}$ and every $\varepsilon>0$

$$
\sum_{j}\left\|p_{j, n}\left(t_{n}\right)-p_{j, n}\left(t_{n}\right) \wedge(1+\varepsilon) p_{j, n}\left(s_{n}\right)\right\|
$$

tends to zero.
PROOF. This is well-known and easily checked. (Let

$$
\Lambda_{\mathrm{j}, \mathrm{n}}=\log \frac{\mathrm{dp}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)}{\mathrm{dp}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{~s}_{\mathrm{n}}\right)}
$$

The condition says that, under $t_{n}$, the sum of the probabilities $\sum_{j} \operatorname{Pr}\left[\Lambda_{j, n}>\log (1+\varepsilon)\right]$ tends to zero. Then apply the necessary and sufficient conditions for the Central Limit Theorem.)

From this result one can immediately obtain the following:
THEOREM 1. Let $\left\{\mathbf{E}_{\mathrm{j}, \mathrm{n}}\right\}$ be a bounded infinitesimal array of experiments indexed by $\Theta_{\mathrm{n}}$. Let $\left\{\mathrm{F}_{\mathrm{j}, \mathrm{n}}\right\}$ be another array such that $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$ is weaker than the corresponding $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$. Let $\overline{\mathbf{E}}_{\mathrm{n}}=\prod_{\mathrm{j}} \mathbf{E}_{\mathrm{j}, \mathrm{n}}$ and $\overline{\mathbf{F}}_{\mathrm{n}}=\prod_{\mathrm{j}} \mathbf{F}_{\mathrm{j}, \mathrm{n}}$. If $\overline{\mathbf{E}}_{\mathrm{n}}$ admits weak Gaussian approximations, so
does $\overline{\mathbf{F}}_{\mathrm{n}}$.
PROOF. It is clear that $\left\{\mathrm{F}_{\mathrm{j}, \mathrm{n}}\right\}$ is also a bounded infinitesimal array. By Lemma 2, Section 3, there are transformations $T_{j, n}$ such that the measures $q_{j, n}(t)$ of $F_{j, n}$ satisfy $q_{j, n}(t)=T_{j, n} p_{j, n}(t) \quad$ and $\quad$ also $\quad\left\|q_{j, n}\left(t_{n}\right)-q_{j, n}\left(t_{n}\right) \wedge(1+\varepsilon) q_{j, n}\left(s_{n}\right)\right\|$ $\leq\left\|p_{j, n}\left(t_{n}\right)-p_{j, n}\left(t_{n}\right) \wedge(1+\varepsilon) p_{j, n}\left(s_{n}\right)\right\|$. Hence the result.

Note that Theorem 1 refers to weak approximability. At the time of the present writing we do not know whether strong Gaussian approximability is preserved. However, a conjecture stated in [26] would imply that, for bounded infinitesimal arrays, strong and weak Gaussian approximability are equivalent.

The structure of a Gaussian experiment is well defined by certain quadratic forms that may be introduced as follows. Let $G_{n}=\left\{G_{\theta, n} ; \theta \in \Theta_{n}\right\}$ be a Gaussian experiment indexed by $\Theta_{n}$. Let $\mathbf{M}_{\mathrm{n}}$ be the linear space of finite signed measures $\mu$ that have finite support on $\Theta_{\mathrm{n}}$ and are such that $\mu\left(\Theta_{\mathrm{n}}\right)=0$. Let $\Lambda_{\mathrm{n}}^{\prime}(\mathrm{t}, \mathrm{s})=\log \frac{\mathrm{dG}_{\mathrm{t}, \mathrm{n}}}{\mathrm{dG}_{\mathrm{s}, \mathrm{n}}}$ and let $\Gamma(\mu)$ be the variance of $\int \Lambda_{n}{ }^{\prime}(t, s) \mu(d t)$. This is a positive semi-definite quadratic form on $\mathbf{M}_{n}$. It can be used to complete a quotient of $\mathbf{M}_{\mathbf{n}}$ to obtain a Hilbert space.

The quadratic form $\Gamma$ controls the behavior of the process $\Lambda_{\mathrm{n}}^{\prime}$ in the following sense: the integrals $\int\left[\Lambda_{\mathrm{n}}^{\prime}(\mathrm{t}, \mathrm{s})+\frac{1}{2} \Gamma\left(\delta_{\mathrm{t}}-\delta_{\mathrm{s}}\right)\right] \mu(\mathrm{dt})$ are almost surely zero if and only if $\Gamma(\mu)=0$. In other words linearity relations satisfied by the random part of $\Lambda_{\mathrm{n}}{ }^{\prime}$ are those describable through $\Gamma$.

For experiments that are only approximately Gaussian, the situation is not so neat but one can describe an analogue as follows.

Let $\Lambda_{n}(t, s)=\log \frac{\mathrm{dP}_{\mathrm{n}}(\mathrm{t})}{\mathrm{dP}_{\mathrm{n}}(\mathrm{s})}$. Let $\Gamma_{\mathrm{n}}$ be a positive semi-definite quadratic form on $M_{n}$
DEFINITION 4. The sequence of experiments $\overline{\mathbf{E}}_{\mathrm{n}}$ is under the control of the quadratic forms $\Gamma_{n}$ if, for elements $\mu_{n} \in \mathbf{M}_{n}$ that have bounded mass ( $\sup _{\mathrm{n}}\left\|\mu_{\mathrm{n}}\right\|<\infty$ ) and supports with bounded cardinality, convergence of $\Gamma_{n}\left(\mu_{n}\right)$ to zero implies that

$$
\int\left\{\Lambda_{n}\left(t, s_{n}\right)+\frac{1}{2} \Gamma_{n}\left(\delta_{t}-\delta_{s_{n}}\right)\right\} \mu_{n}(d t)
$$

tends to zero in $\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)$ probability.
This is obviously a way to say that, if $\Gamma_{n}\left(\mu_{n}\right)$ tends to zero, the recentered $\Lambda_{n}$ satisfy approximately the corresponding linear relation. The literature contains several conditions meant to express that log-likelihood admit approximate linear-quadratic expansions. This is usually done assuming that $\Theta_{\mathrm{n}}$ is a finite dimensional vector space. One writes that $\Lambda_{n}(t+\theta, \theta)$ is approximately $t^{\prime} V_{n}-\frac{1}{2} t^{\prime} M_{n} t$. In the LAN
conditions $M_{n}$ is non random. In the LAMN conditions $M_{n}$ is random, the pair $\left(V_{n}, M_{n}\right)$ has a limiting distribution $L(V, M)$ where, conditionally given $M$, the variable $\mathrm{t}^{\prime} \mathrm{V}$ is $\mathrm{N}\left[0, \mathrm{t}^{\prime} \mathrm{Mt}\right]$. The case where this restriction on conditional distributions is omitted has been considered by R. Davies [8] as the case "when the amount of information is random". For these conditions see Basawa and Prakasa Rao [4], Basawa and Scott [5], P. Feigin [11], and Jeganathan [17].

Note that all these conditions imply a relation between the linear structure of $\Theta_{\mathrm{n}}$ and that of the log-likelihood. Our quadratic control condition is different. It does not rely on any particular structure, linear or otherwise on $\Theta_{n}$. To understand it let us consider experiments $E=\left\{P_{\theta} ; \theta \in \Theta\right\}$ that do not depend on $n$ at all. Suppose that $\Theta$ is the real line and that $P_{\theta}$ is the distribution of a two dimensional vector $(X, Y)$, jointly normal with covariance matrix the identity and expectations $\mathrm{E}_{\theta} \mathrm{X}=\theta$, $\mathrm{E}_{\theta} \mathrm{Y}=\sinh \theta$. Then $\mathbf{E}$ does not satisfy the LAN or LAMN conditions on $\Theta$, yet it is controlled by a quadratic $\Gamma$ such that $\Gamma\left(\delta_{t}-\delta_{s}\right)=|t-s|^{2}+|\sinh t-\sinh s|^{2}$ for the Dirac masses $\delta_{s}$ and $\delta_{\mathrm{t}}$. This is typical of Gaussian shift experiments where the shift parameters are not linearly related to the original $\theta$.

However the quadratic control places a strong restriction on the log-likelihood. For instance, let $P_{\theta}$ be the ordinary gamma distribution with density $\frac{\beta^{\alpha}}{\Gamma(\alpha)} e^{-\beta x} x^{\alpha-1}$. The $\log$-likelihood is a linear function of the vector $\{\beta, \alpha, \log \Gamma(\alpha)\}$ but the experiment $E=\left\{\mathrm{P}_{\theta} ; \theta \in \Theta\right\}$ with $\theta=(\alpha, \beta)$ does not admit quadratic control.

The following results say that linear relations are preserved by going from the product $\prod_{\mathrm{j}} \mathrm{E}_{\mathrm{j}, \mathrm{n}}$ to the weaker $\prod_{\mathrm{j}} \mathrm{F}_{\mathrm{j}, \mathrm{n}}$.

THEOREM 2. Let $\left\{\mathbf{E}_{\mathrm{j}, \mathrm{n}}\right\}$ be a bounded infinitesimal array of experiments $\mathbf{E}_{\mathrm{j}, \mathrm{n}}=\left\{\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$. Let $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$ be weaker than $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$.

If the products $\overline{\mathbf{E}}_{\mathrm{n}}$ admit weak Gaussian approximations they are under the control of quadratic forms $\Gamma_{\mathrm{n}}$ and the weaker products $\overline{\mathrm{F}}_{\mathrm{n}}$ are under the control of quadratic forms $\Gamma_{n}^{*}$ such that $\Gamma_{n}-\Gamma_{n}^{*}$ is positive semi-definite.

PROOF. Select a point $s_{n} \in \Theta_{n}$ and let

$$
X_{j, n}(t)=\frac{\mathrm{dp}_{\mathrm{j}, \mathrm{n}}(\mathrm{t})}{\mathrm{dp}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{~s}_{\mathrm{n}}\right)}-1
$$

Let $Y_{j, n}(t)=1 \wedge X_{j, n}(t)$. It is easy to see that, under the conditions given, for any sequence $t_{n}$, the difference between $\Lambda_{n}\left(t_{n}, s_{n}\right)$ and $\sum_{j} Y_{j, n}\left(t_{n}\right)-\frac{1}{2} \sum_{j}$ variance $Y_{j, n}\left(t_{n}\right)$ tends to zero in $\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)$ probability. (See for instance [22] or [27].)

On the space of measures $\mathbf{M}_{n}$, define $\Gamma_{n}(\mu)$ as the variance under $s_{n}$ of $\int\left[\sum_{\mathrm{j}} \mathrm{Y}_{\mathrm{j}, \mathrm{n}}(\mathrm{t})\right] \mu(\mathrm{dt})$. Then the integral

$$
\int\left\{\Lambda_{n}\left(t, s_{n}\right)+\frac{1}{2} \Gamma_{n}\left[\delta_{t}-\delta_{s_{n}}\right]\right\} \mu_{n}(\mathrm{dt})
$$

differs from $\int\left[\sum_{j} Y_{j, n}(t)\right] \mu_{n}(d t)$ by a quantity that tends to zero in probability as long as $\left\|\mu_{\mathrm{n}}\right\|$ remains bounded and as long as the supports of the $\mu_{\mathrm{n}}$ have bounded cardinality. Thus the assertion that $\overline{\mathbf{E}}_{\mathrm{n}}$ is under the control of the quadratic form $\Gamma_{\mathrm{n}}$ is equivalent to the statement that $\int\left[\sum_{j} Y_{j, n}(t)\right] \mu_{n}(d t)$ tends to zero in probability whenever its variance tends to zero. However the contiguity restrictions involved in the Gaussian approximability imply that the distribution of a term such as $\sum_{j} Y_{j, n}(t)$ is approximated by a Gaussian distribution with expectation zero. Hence the first assertion.

To obtain the second result, let us use again the transformations $\mathrm{T}_{\mathrm{j}, \mathrm{n}}$ with $\mathrm{T}_{\mathrm{j}, \mathrm{n}} \mathrm{p}_{\mathrm{j}, \mathrm{n}}(\mathrm{t}) \equiv \mathrm{q}_{\mathrm{j}, \mathrm{n}}(\mathrm{t})$. Under some technical regularity conditions these $\mathrm{T}_{\mathrm{j}, \mathrm{n}}$ can be represented by Markov kernels. If we look only at finite subsets of $\Theta_{n}$ and only at distributions of likelihood ratios for the $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$, these technical conditions can be assumed to be satisfied. Thus the situation can be described as follows. Let $\mathbf{A}_{\mathrm{j}, \mathrm{n}}$ and $\mathbf{B}_{\mathrm{j}, \mathrm{n}}$ be respectively the $\sigma$-fields of $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ and $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$.

Let $\mathrm{K}_{\mathrm{j}, \mathrm{n}}(\mathrm{B}, \cdot)$ be the Markov kernel of $\mathrm{T}_{\mathrm{j}, \mathrm{n}}$ evaluated at B in $\mathrm{B}_{\mathrm{j}, \mathrm{n}}$. Define measures by $M_{j, n}[t ; A \times B]=\int_{A} K_{j, n}(B, \cdot) d p_{j, n}(t)$. Their marginals on $B_{j, n}$ are the $q_{j, n}(t)$ and their likelihood ratios are $A_{j, n}$ measurable, equivalent to those of the $p_{j, n}(t)$. Now replace our previous $\mathrm{Y}_{\mathrm{j}, \mathrm{n}}(\mathrm{t})$ by the equivalent

$$
\frac{\mathrm{dM}_{\mathrm{j}, \mathrm{n}}^{\prime}(\mathrm{t})}{\mathrm{dM}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{~s}_{\mathrm{n}}\right)}-1
$$

with $M_{j, n}^{\prime}(t)=M_{j, n}(t) \wedge 2 M_{j, n}\left(S_{n}\right)$. Let $Z_{n, j}(t)$ be the conditional expectation (under $\mathrm{s}_{\mathrm{n}}$ ) of $Y_{j, n}(t)$ given $B_{j, n}$. This differs from

$$
\frac{d q_{j, n}(t)}{d q_{j, n}\left(s_{n}\right)}-1
$$

by an amount that can be neglected. It is thus possible to show that

$$
\sum_{\mathrm{j}} \mathrm{Z}_{\mathrm{n}, \mathrm{j}}(\mathrm{t})-\frac{1}{2} \text { variance } \sum_{\mathrm{j}} \mathrm{Z}_{\mathrm{n}, \mathrm{j}}(\mathrm{t})
$$

is an approximation to the logarithm of

$$
\frac{\mathrm{dQ}_{\mathrm{n}}(\mathrm{t})}{\mathrm{dQ}_{\mathrm{n}}\left(\mathrm{~s}_{\mathrm{n}}\right)}
$$

for the product $\overline{\mathbf{F}}_{\mathrm{n}}$. For the same reasons as above, the $\overline{\mathrm{F}}_{\mathrm{n}}$ are under the control of the quadratic forms defined by $\Gamma_{n}^{*}(\mu)=\operatorname{var} \int\left[\sum_{j} Z_{n, j}(t)\right] \mu(d t)$. However, since $Z_{n, j}(t)$ is the conditional expectation of $Y_{j, n}(t)$ the variances are smaller so that $\Gamma_{n}^{*}(\mu) \leq \Gamma_{n}(\mu)$ for all $\mu \in \mathbf{M}_{\mathrm{n}}$. This completes the proof of the Theorem.

REMARK 1. The forms $\Gamma_{\mathrm{n}}$ and $\Gamma_{\mathrm{n}}^{*}$ used here are nonrandom. We shall see in Section 8 experiments that are under the control of random quadratic forms. Under some supplementary restrictions (contiguity for pairs $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)\right\}$, $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)\right\}$, convexity of $\Theta_{\mathrm{n}}$ if imbedded in the Hilbert space) one can show that control by nonrandom quadratic forms is equivalent to weak Gaussian approximability.

REMARK 2. The assumption that we have products with bounded infinitesimal arrays $\left\{\mathrm{E}_{\mathrm{j}, n}\right\}$ and $\left\{\mathrm{F}_{\mathrm{j}, n}\right\}$ and with $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$ weaker than $\mathrm{E}_{\mathrm{j}, \mathrm{n}}$ has been used very forcibly in the above proofs. The results might be extendable to some other cases, as shown for instance in Section 6. However they are not true without restriction. One can have an experiment $\mathbf{F}$ weaker than a Gaussian $\mathbf{E}$ with $\mathbf{F}$ remote from Gaussian. For an instance let $\mathbf{E}$ be the family $\mathbf{N}(\theta, 1)$ of normal distributions with variance unity on the line. If X is the observable variable let Y be the integer closest to X . It does give a family of distributions $F=\left\{Q_{\theta} ; \theta \in \Theta_{n}\right\}$ weaker than $E$. However the log likelihood ratios for the $Q_{\theta}$ have discrete distributions that are far from Gaussian.

## 5. Preservation of the LAN conditions, part I.

The contents of the present section bear on two separate questions. The first is that, for the bounded infinitesimal arrays of Section 4, the LAN conditions are preserved under loss of information if this loss occurs on the individual components. The second refers to the construction of estimates by adding to our auxiliary estimate a correction calculable from likelihood ratios. Here many different procedures may be used. We describe some that can often be easily carried out in practice. An explicit example is given in Section 9, Example 1.

The first recorded instance of the conditions called LAN that has come to our attention is that of [20]. Since then conditions called LAN have been used by many authors. (See [8], [10], [14], [34] for instance.) Unfortunately the statements used by various authors differ in some aspects that are sometimes incidental and sometimes important. Because of this we shall first give two sets of assumptions, and show that they are equivalent if the index sets $\Theta_{\mathrm{n}}$ are sufficiently rich. Then we shall show that they are preserved by passages from a product $\overline{\mathbf{E}}_{\mathrm{n}}$ to the weaker $\overline{\mathbf{F}}_{\mathrm{n}}$ of Section 2 . We end by a description of a method of construction of estimates.

The LAN (for 'locally asymptotically normal') assumptions used in the references listed above differ considerably from the asymptotic Gaussian approximability conditions of Definitions 1 and 2 of Section 4. They involve a finite dimensional vector space V and a sequence of norms $\mid \cdot I_{\mathrm{n}}$ on V . The sequence of norms is used for two purposes: a) to indicate the size of sets on which Gaussian approximability is contemplated and, b) to relate the linear structure of V to that of the Gaussian experiments used as approximations. The relations will be described in detail below.

For the typical LAN assumptions one considers a fixed finite dimensional real vector space V and a sequence $\left\{\mid \cdot \rho_{\mathrm{n}}\right\}$ of norms on V . The parameter set $\Theta_{\mathrm{n}}$ that indexes the experiments $\bar{E}_{n}=\left\{P_{n}(\theta) ; \theta \in \Theta_{n}\right\}$ is mapped into $V$ by some function, say $\tau_{n}$. The assumptions usually refer to some particular $\theta_{0, \mathrm{n}} \in \Theta_{\mathrm{n}}$ called the true value. The norms are often obtained from a single norm $|\cdot|$ on V by using either multiplication by some numerical factors or more complicated matrix manipulation. Think of the fairly common renormalizations $\left|\theta-\theta_{0, n} \|_{n}=\sqrt{n}\right| \theta-\theta_{0, n} \mid$ or $\left|\theta-\theta_{0, n}\right|_{n}=\left|M_{n}\left(\theta-\theta_{0, n}\right)\right|$ for matrices $M_{n}$. These $M_{n}$ are often selected to "stabilize" the image of $\Theta_{n}$ in $V$ making it tend to some limit.

These renormalization and stabilization operations may be convenient in practice but they distract attention from the main statistical arguments. Thus we shall not consider them.

In fact, for simplicity, we shall just assume that $\Theta_{\mathrm{n}}$ is a subset of V and that $\theta_{0, \mathrm{n}}$ is the origin of V .

Consider then experiments $\overline{\mathbf{E}}_{\mathrm{n}}$ subject to the following requirements.
(R1) Each $\Theta_{\mathrm{n}}$ is a subset of the fixed space V and the origin of V belongs to $\Theta_{\mathrm{n}}$.
(R2) If $\mid \theta_{n} l_{n}$ remains bounded the sequences $\left\{\mathrm{P}_{\mathrm{n}}\left(\theta_{\mathrm{n}}\right)\right\}$ and $\left\{\mathrm{P}_{\mathrm{n}}\left(\theta_{0, n}\right)\right\}$ are contiguous.
(R3) For any set $S_{n} \subset \Theta_{n}$ such that $\sup \left\{|s|_{n} ; s \in S_{n}, n=1,2, \ldots\right\}<\infty$ the experiments $\left\{\mathrm{P}_{\mathrm{n}}(\theta) ; \theta \in \mathrm{S}_{\mathrm{n}}\right\}$ admit weak Gaussian approximations according to Definition 2 of Section 4.

To state the remaining conditions proceed as in Section 3 and introduce the set $\mathbf{M}_{\mathrm{n}}$ of finite signed measures $\mu$ carried by finite subsets of $\Theta_{n}$ and satisfying $\mu\left(\Theta_{n}\right)=0$.

For subsets $S_{n} \subset \Theta_{n}$, let $\mathbf{M}_{n}\left(S_{n}\right)$ be the subspace of $\mathbf{M}_{n}$ formed by measures $\mu$ whose support is in $S_{n}$ (with $\mu\left(S_{n}\right)=0$ ).

The existence of Gaussian approximations $G_{n}=\left\{G_{\theta, n} ; \theta \in S_{n}\right\}$ implies the existence of corresponding quadratic forms $\Gamma_{n, S_{n}}$ defined on $\mathbf{M}_{n}\left(S_{n}\right)$. For $\mu \in \mathbf{M}_{n}\left(S_{n}\right)$ and for

$$
\Lambda_{\mathrm{n}}^{*}\left(\mathrm{t}, \mathrm{~s}_{\mathrm{n}}\right)=\log \frac{\mathrm{dG}_{\mathrm{t}, \mathrm{n}}}{\mathrm{dG}_{\mathrm{s}_{\mathrm{n}}, \mathrm{n}}}
$$

the value $\Gamma_{n, S_{n}}(\mu)$ is the variance of $\int \Lambda_{n}^{*}\left(t, s_{n}\right) \mu(d t)$.
The link between the linear structure of V and that of the Gaussian approximations is as follows.
(R4) Consider sets $S_{n} \subset \Theta_{n}$ such that $\sup \left\{|s|_{n} ; s \in S_{n}, n=1,2, \ldots\right\}<\infty$ and such that the cardinality card $S_{n}$ remains bounded. Then, there are Gaussian approximations with the following property:

If $\mu_{n} \in \mathbf{M}_{n}\left(S_{n}\right)$ is such that $\sup _{\mathrm{n}}\left\|\mu_{\mathrm{n}}\right\|<\infty$ and such that $\left|\int \mathrm{s} \mu_{\mathrm{n}}(\mathrm{ds})\right|_{\mathrm{n}}$ tends to zero, then $\Gamma_{n, S_{n}}\left(\mu_{\mathrm{n}}\right) \rightarrow 0$.
It should be clear that the role of the norms in (R3) is to indicate the size of sets on which Gaussian approximability is contemplated. The role of (R4) is to link the vector structure of V to that of the Gaussian approximations.

This set (R1)-(R4) of assumptions is often described in a manner that looks very different as follows:

Let

$$
\Lambda_{\mathrm{n}}(\mathrm{t})=\log \frac{\mathrm{dP}_{\mathrm{n}}(\mathrm{t})}{\mathrm{dP}_{\mathrm{n}}(0)}
$$

(L1) There are random vectors $\mathrm{W}_{\mathrm{n}}$ (with values in V ) and Euclidean norms $\|\cdot\|_{\mathrm{n}}$ such that if $\sup _{\mathrm{n}}\left|\mathrm{t}_{\mathrm{n}}\right|_{n}<\infty$ then

$$
\Lambda_{n}\left(t_{n}\right)-t_{n}^{\prime} W_{n}+\frac{1}{2}\left\|t_{n}\right\|_{n}^{2}
$$

tends to zero in $\mathrm{P}_{\mathrm{n}}(0)$ probability.
(L2) Let $F_{n}$ be the distribution under $P_{n}(0)$ of the vector $W_{n}$. There are joint Gaussian distributions $G_{n}$ such that the dual Lipschitz norm $\left\|F_{n}-G_{n}\right\|_{D}$ induced by $\mid \cdot I_{\text {n }}$ tends to zero. (The dual Lipschitz norm is

$$
\left\|F_{n}-G_{n}\right\|_{D}=\sup _{f}\left|\int f d F_{n}-\int f d G_{n}\right|
$$

for a supremum taken over functions $f$ such that $|f| \leq 1$ and $|f(x)-f(y)| \leq|x-y|_{n}$.)
The sets (R1), (R2), (R3), (R4) and (R1), (R2), (L1), (L2) are technically very different. To relate them one encounters technical difficulties due to the fact that the sets $\Theta_{\mathrm{n}}$ may be very sparse subsets of V . However the two systems are known to be equivalent (see [27] for instance) under the following restriction.
(R5) Let $k$ be the dimension of $V$. There are sets $\left\{s_{0, n}, s_{1, n}, \cdots, s_{k, n}\right\} \subset \Theta_{n}$ with the following properties:
i) $\sup _{\mathrm{i}, \mathrm{n}}\left|\mathrm{s}_{\mathrm{i}, \mathrm{n}}\right|_{\mathrm{n}}<\infty$
ii) if $t_{n} \in \Theta_{n}$ is such that $\sup _{n}\left|t_{n}\right|_{n}<\infty$ then there are numbers $c_{n, i} ; i=1,2, \ldots, k$ such that

$$
t_{n}-s_{0, n}=\sum_{i=1}^{k} c_{n, i}\left(s_{i, n}-s_{0, n}\right)
$$

and such that $\sum_{i=1}^{k}\left|c_{n, i}\right|$ remains bounded.
(If one requires more, for instance that the simplex spanned by the $\mathrm{s}_{\mathrm{i}, \mathrm{n}}$ be contained in $\Theta_{n}$, then condition (L2) is already a consequence of (R1), (R2), (L1).)

As seen in Section 4, it is not true that replacing an experiment $\overline{\mathbf{E}}_{\mathrm{n}}$ by a weaker one will preserve the Gaussian approximability. Thus the LAN conditions (R1)-(R4) are not preserved either. However for product experiments one can assert the following.

THEOREM 3. Let $\left\{\mathbf{E}_{\mathrm{j}, \mathrm{n}}\right\}$ be a bounded infinitesimal array of experiments and let $\left\{\mathbf{F}_{\mathrm{j}, \mathrm{n}}\right\}$ be another double array all indexed by sets $\Theta_{\mathrm{n}}$. Assume that $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$ is less informative than $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ and consider the products $\overline{\mathbf{E}}_{\mathrm{n}}=\prod_{\mathrm{j}} \mathbf{E}_{\mathrm{j}, \mathrm{n}}$ and $\overline{\mathbf{F}}_{\mathrm{n}}=\prod_{\mathrm{j}} \mathbf{F}_{\mathrm{j}, \mathrm{n}}$.

If $\overline{\mathbf{E}}_{\mathrm{n}}$ satisfies the conditions (R1), (R2), (R3), (R4), so does $\overline{\mathbf{F}}_{\mathrm{n}}$. If in addition the requirement $(R 5)$ is satisfied one can replace the conditions $(R 1),(R 2),(R 3),(R 4)$ by (R1), (R2), (L1), (L2).

PROOF. The preservation of the conditions (R1), (R2), (R3), (R4) is an immediate consequence of Theorems 1 and 2 of Section 4: contiguity, Gaussian approximability are preserved and so are the linear relations satisfied by the Gaussian approximations. Under (R1), (R2), (R5) the pairs (R3), (R4) and (L1), (L2) are always equivalent. Hence the result.

To be complete we should prove the stated equivalence of (R3), (R4) and (L1), (L2). However, for bounded infinitesimal arrays one can use a number of other relations. Since these are important for the possible construction of asymptotically sufficient estimates we shall now proceed to state some of them.

Suppose for instance that (R1)-(R5) are all satisfied. For $i=1,2, \ldots, k$ let $X_{j, n, i}$ be the random variable

$$
\mathrm{X}_{\mathrm{j}, \mathrm{n}, \mathrm{i}}=\left[\frac{\mathrm{dp}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{~s}_{\mathrm{i}, \mathrm{n}}\right)}{\mathrm{dp}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{~s}_{0, \mathrm{n}}\right)}-1\right] \wedge 1
$$

Let $Y_{n, i}$ be the sum $\sum_{j} X_{j, n, i}$. Suppose that $t_{n} \in \Theta_{n}$ is such that $\sup _{\mathrm{n}}\left|t_{n}\right|_{n}<\infty$. Then, by (R5), there are numbers such that

$$
t_{n}-s_{0, n}=\sum_{i} c_{n, i}\left(s_{i, n}-s_{0, n}\right)
$$

and such that $\sup _{\mathrm{n}} \sum_{\mathrm{i}}\left|\mathrm{c}_{\mathrm{n}, \mathrm{i}}\right|<\infty$. Let $\mathrm{D}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)$ be the difference

$$
D_{n}\left(t_{n}\right)=\sum_{j}\left[\frac{d p_{j, n}\left(t_{n}\right)}{d p_{j, n}\left(s_{0, n}\right)}-1\right]-\sum_{i=1}^{k} c_{n, i} Y_{n, i}
$$

We claim that, under (R1) to (R5), the difference $D_{n}\left(t_{n}\right)$ will tend in probability to zero. Indeed, let $\mu_{n}$ be the measure that assigns mass ( $-c_{n, i}$ ) to $s_{i, n}$ for $i=1,2, \ldots, k$, mass 1 to $t_{n}$ and mass $\left(\sum_{i=1}^{k} c_{n, i}\right)-1$ to $s_{0, n}$. By definition of the $c_{n, i}$ one has $\int \mathrm{t} \mu_{\mathrm{n}}(\mathrm{dt})=0$. Consider then the logarithms

$$
\Lambda_{\mathrm{n}}(\mathrm{t})=\log \frac{\mathrm{dP}_{\mathrm{n}}(\mathrm{t})}{\mathrm{dP}_{\mathrm{n}}(0)}
$$

and the integrals $\int \Lambda_{n}(t) \mu_{n}(d t)$. According to (R1)-(R3) these integrals have a distribution that is approximately Gaussian with a certain variance $\Gamma_{n}\left(\mu_{n}\right)$. Since $\int t \mu_{n}(d t)=0$ this variance $\Gamma_{n}\left(\mu_{n}\right)$ must tend to zero, by (R4). However, $D_{n}\left(t_{n}\right)$ is approximately equal to $\int \Lambda_{n}(t) \mu_{n}(d t)$ recentered by subtracting the expectation of the

Gaussian approximation. Hence the assertion.
The same conclusion can be obtained if one uses (L1), (L2) instead of (R3), (R4).
To see how this may be used for the construction of estimates, consider the weaker $F_{j, n}=\left\{q_{j, n}(\theta) ; \theta \in \Theta_{n}\right\}$ and the same sets $S_{n}=\left\{s_{0, n}, s_{1, n}, \cdots, s_{k, n}\right\}$.

An estimation technique that is often successful is to maximize a smooth approximation to the logarithms

$$
L_{n}(t)=\log \frac{\mathrm{dQ}_{\mathrm{n}}(\mathrm{t})}{\mathrm{dQ}_{\mathrm{n}}\left(\mathrm{~s}_{0, \mathrm{n}}\right)}
$$

Note that we said to maximize a smooth approximation not the $\mathrm{L}_{\mathrm{n}}(\mathrm{t})$ themselves. Let

$$
\left.\mathrm{Z}_{\mathrm{j}, \mathrm{n}, \mathrm{i}}=\frac{\mathrm{dq}}{\mathrm{j}, \mathrm{n}} \text { ( } \mathrm{s}_{\mathrm{i}, \mathrm{n}}\right)(1
$$

and let $\mathrm{Z}_{\mathrm{j}, \mathrm{n}, \mathrm{i}}^{*}=\mathrm{Z}_{\mathrm{j}, \mathrm{n}, \mathrm{i}} \wedge 1$.
Under the conditions of Theorem 3 for $t_{n}-s_{0, n}=\sum_{i=1}^{k} v_{n, i}\left(s_{i, n}-s_{0, n}\right)$ a possible approximation to $L_{n}(t)$ is

$$
L_{n}^{(1)}\left(t_{n}\right)=\sum_{j} \sum_{i} v_{n, i} Z_{j, n, i}^{*}-\frac{1}{2} \operatorname{var}\left[\sum_{j} \sum_{i} v_{n, i} Z_{j, n, i}^{*}\right] .
$$

This consists of a random term that is linear in the vector $\mathrm{v}_{\mathrm{n}}{ }^{\prime}=\left\{\mathrm{v}_{\mathrm{n}, \mathrm{i}} ; \mathrm{i}=1, \ldots, \mathrm{k}\right\}$ and a nonrandom term that is quadratic in that vector. In many cases the variances are not readily computable. Then one may prefer to use the approximation

$$
L_{n}^{(2)}\left(t_{n}\right)=\sum_{j} \sum_{i} v_{n, i} Z_{j, n, i}-\frac{1}{2} \sum_{j}\left[\sum_{i} v_{n, i} Z_{j, n, i}\right]^{2}
$$

In matrix notation this may be written as follows. Let $\mathrm{Z}_{\mathrm{n}}$ be the column vector whose coordinates are the $\sum_{j} Z_{j, n, i}$. Let $M_{n}$ be the matrix whose ( $\alpha, \beta$ ) entry is $\sum_{j} Z_{j, n, \alpha} Z_{j, n, \beta}$. Then

$$
L_{n}^{(2)}\left(\mathrm{t}_{\mathrm{n}}\right)=\mathrm{v}_{\mathrm{n}}^{\prime} \mathrm{Z}_{\mathrm{n}}-\frac{1}{2} \mathrm{v}_{\mathrm{n}}^{\prime} \mathrm{M}_{\mathrm{n}} \mathrm{v}_{\mathrm{n}} .
$$

When $\mathrm{M}_{\mathrm{n}}$ is invertible this can also be written in the form

$$
L_{n}^{(2)}\left(t_{n}\right)=-\frac{1}{2}\left\{\left(v_{n}-M_{n}^{-1} Z_{n}\right)^{\prime} M_{n}\left(v_{n}-M_{n}^{-1} Z_{n}\right)-Z_{n}^{\prime} M_{n}^{-1} Z_{n}\right\}
$$

from which it follows that the maximum of this approximation is reached at the point $\hat{v}_{\mathrm{n}}=\mathrm{M}_{\mathrm{n}}^{-1} \mathrm{Z}_{\mathrm{n}}$.

To carry out the above operations one needs to know the sets $\left\{\mathrm{s}_{0, \mathrm{n}}, \mathrm{s}_{1, \mathrm{n}}, \cdots, \mathrm{s}_{\mathrm{k}, \mathrm{n}}\right\}$, (and to be able to obtain good approximations to the likelihood ratios). Of course, if the purpose is to estimate $\theta$, one will not know what set $\left\{s_{0, n}, \cdots, s_{k, n}\right\}$ to use.

The technique proposed in [20] and variously expounded in [22] and [27] is to replace the set $\left\{s_{0, n}, \cdots, s_{k, n}\right\}$ by a suitable estimate of $i t$.

As explained in [20] the validity of the technique seemed to depend on curious cancellations involving the difference between $\mathrm{s}_{0, \mathrm{n}}$ and an estimate $\hat{\mathrm{s}}_{0, \mathrm{n}}$. However the basic reason is very simple: the points that maximize $L_{n}\left(t_{n}\right)$ are the same as those that maximize $L_{n}\left(t_{n}\right)-L_{n}\left(\hat{s}_{0, n}\right)$. This remains approximately true for the approximations of the type $\mathrm{L}_{\mathrm{n}}^{(1)}$ or $\mathrm{L}_{\mathrm{n}}^{(2)}$.

Note however that for the validity of the argument the following must hold:

1) The auxiliary estimate $\hat{\mathrm{s}}_{\mathrm{i}, \mathrm{n}}$ must with large probability be in the range where the approximation holds.
2) Similarly, the correction term $\hat{v}_{n}=M_{n}^{-1} \mathrm{Z}_{\mathrm{n}}$ must be in the range where the approximation holds.

In the present situation this calls for the following comments. The existence of estimates that converge rapidly enough for $\overline{\mathbf{E}}_{\mathrm{n}}$ does not imply existence of similar estimates for $\bar{F}_{n}$. The fact that in Theorem 2, the difference $\Gamma_{n}-\Gamma_{n}^{*}$ is positive semidefinite implies that the matrix $\mathrm{M}_{\mathrm{n}}$ will typically be smaller than the corresponding matrix calculated on $\overline{\mathbf{E}}_{\mathrm{n}}$. The fact that this latter would not degenerate is no guarantee that $\mathrm{M}_{\mathrm{n}}$ will not. These properties must be checked for $\overline{\mathbf{F}}_{\mathrm{n}}$ itself.

The matrix $\mathrm{M}_{\mathrm{n}}$ can be adjusted to a certain extent by modification of the basis $\hat{\mathrm{s}}_{\mathrm{i}, \mathrm{n}}-\hat{\mathrm{s}}_{0, \mathrm{n}} ; \mathrm{i}=1,2, \ldots \mathrm{k}$. It may be wise to check that a change of basis does not modify the end result too substantially.

In spite of all these caveats we shall describe, in Section 7 below, a situation where the necessary verification can be carried out rather easily.

To terminate this section, note that the approximations $L_{n}^{(1)}$ and $L_{n}^{(2)}$ are not the only possible ones. There are many more. In [27] one of the present authors suggested estimating what takes the place of $\mathrm{M}_{\mathrm{n}}$ by parallelogram differences calculated on the log likelihood. For the situation covered by Theorem 3 one can use instead of the

$$
\frac{\mathrm{dq}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{~s}_{\mathrm{i}, \mathrm{n}}\right)}{\mathrm{dq}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{~s}_{0, \mathrm{n}}\right)}
$$

a variety of functions of them, in particular their logarithms or their square roots. For instance, if

$$
\left.\left.\mathrm{W}_{\mathrm{j}, \mathrm{n}, \mathrm{i}}=\left\{\frac{\mathrm{dq}}{\mathrm{j}, \mathrm{n}}{ }^{\mathrm{dq}} \mathrm{~s}_{\mathrm{j}, \mathrm{n}, \mathrm{n}}\right) \mathrm{s}_{0, \mathrm{n}}\right)\right\}^{1 / 2}-1
$$

then

$$
L_{n}^{(3)}\left(t_{n}\right)=\sum_{j} \sum_{i} c_{i, n} W_{j, n, i}-\frac{1}{2} \sum_{j}\left[\sum_{i} c_{i, n} W_{j, n, i}\right]^{2}
$$

is an approximation to $\frac{1}{2} \mathrm{~L}\left(\mathrm{t}_{\mathrm{n}}\right)$.
The estimates $\left\{\hat{s}_{0, \mathrm{n}}, \cdots, \hat{\mathrm{s}}_{\mathrm{k}, \mathrm{n}}\right\}$ used in [20] for instance were assumed to satisfy certain discreteness conditions. This is for technical reasons: They should not seek for singularities in the likelihood functions. That can be achieved by methods other than discretization. If the likelihood functions are sufficiently regular one does not need to worry about this particular point. A referee pointed out that the estimates described here may display a severe lack of robustness. This is particularly true if the approximation called $\mathrm{L}_{\mathrm{n}}^{(2)}$ is used with untruncated variables. Fortunately one can often detect easily forms of misbehavior of the variables $\mathrm{Z}_{\mathrm{j}, \mathrm{n}, \mathrm{i}}$. Possible methods for doing so and for correcting the system are too complex to be described here. We shall return to this matter elsewhere.

Finally instead of using differences one can sometimes use derivatives and the Newton-Raphson procedure. However to show that this particular technique will work one needs assumptions that are much more restrictive than the ones used here.

The reader will note that in Sections 4 and 5 we replaced each component $\mathbf{E}_{\mathrm{j}, n}$ by a weaker $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$. This is essential in the proofs. The next section expounds on a result of R. Davies in which the loss of information occurs in a very different manner.

## 6. Preservation of the LAN conditions, part II.

R. Davies [8], [9] has observed that the LAN conditions are preserved under circumstances that are rather different from the ones described in Section 5: No product structure is assumed. Even if it is there, the loss of information occurs by passage to sub- $\sigma$-fields that do not preserve the product structure. To describe the situation we shall consider a sequence $\left\{E_{n}\right\}$ of experiments $\left\{P_{n}(\theta) ; \theta \in \Theta\right\}$ where the $P_{n}(\theta)$ are measures on some $\sigma$-field $A_{n}$. The weaker experiment $F_{n}=\left\{Q_{n}(\theta) ; \theta \in \Theta_{n}\right\}$ is obtained by taking the restrictions $Q_{n}(\theta)$ of the $P_{n}(\theta)$ to some subfield $B_{n}$.

We shall show that the LAN conditions are preserved if the weaker experiments admit "distinguished" statistics whose joint distributions with the log likelihood of the initial experiment are asymptotically normal. The relations with Davies' work are complex. They are described at the end of the section.

We shall assume that the experiments $\mathbf{E}_{\mathrm{n}}$ satisfy the requirements (R1), (R2) and (L1), (L2) of Section 5. Recall that (L1) involves certain random vectors $\mathrm{W}_{\mathrm{n}}$ and Euclidean norms \|• $\|_{n}$.

Consider also statistics $\mathrm{T}_{\mathrm{n}}$ defined on the experiment $\mathrm{F}_{\mathrm{n}}$ and taking values in a fixed Euclidean space $\mathbf{R}^{\mathrm{q}}$. Recall that the sequence $\left\{\mathrm{T}_{\mathrm{n}}\right\}$ is called distinguished for the sequence $\left\{F_{n}\right\}$ if the following property holds. Let $F_{n}(\theta)$ be the distribution of $T_{n}$ for $Q_{n}(\theta)$. Take pairs ( $s_{n}, t_{n}$ ) of elements of $\Theta_{n}$. This gives a pair $\left\{F_{n}\left(s_{n}\right), F_{n}\left(t_{n}\right)\right\}$ of measures on $\mathbf{R}^{q}$ and a binary experiment $F_{n}{ }^{\prime}=\left\{Q_{n}\left(s_{n}\right), Q_{n}\left(t_{n}\right)\right\}$.

Compactifying $\mathbf{R}^{\mathrm{q}}$ by adjunction of points at infinity one can extract subsequences $\{v\} \subset\{n\}$ such that $i$ ) the experiments $F_{v}{ }^{\prime}$ have a limit $F$ and ii) the measures $F_{v}\left(s_{v}\right)$ and $F_{v}\left(t_{v}\right)$ tend in the usual weak sense to certain limits say $F_{0}$ and $F_{1}$.

The sequence $\left\{\mathrm{T}_{\mathrm{n}}\right\}$ is distinguished if for all such subsequences the experiment formed by the limiting distributions $\left\{\mathrm{F}_{0}, \mathrm{~F}_{1}\right\}$ is as strong as the limit experiment $\mathrm{F}^{\prime}$. (The definition of [24] involves more than pairs, but it is shown in [27] that it is enough to look at pairs.) For the purposes of our next theorem we shall metrize $V \times \mathbf{R}^{q}$ by the square norms $\|\cdot\|_{n}^{2}+\left.1 \cdot\right|^{2}$.

THEOREM 4. Consider experiments $\mathbf{E}_{\mathrm{n}}=\left\{\mathrm{P}_{\mathrm{n}}(\theta) ; \theta \in \Theta\right\}$ given by measures on $a$ $\sigma$-field $\mathbf{A}_{\mathrm{n}}$ and their restrictions $\mathbf{F}_{\mathrm{n}}=\left\{\mathrm{Q}_{\mathrm{n}}(\theta) ; \theta \in \Theta\right\}$ to subfields $\mathbf{B}_{\mathrm{n}} \subset \mathbf{A}_{\mathrm{n}}$.

Assume the following
i) The $\mathrm{E}_{\mathrm{n}}$ satisfy the LAN conditions (R1), (R2), (L1) and (L2).
ii) The joint distributions $\mathrm{L}\left[\mathrm{W}_{\mathrm{n}}, \mathrm{T}_{\mathrm{n}} \mid \mathrm{P}_{\mathrm{n}}\left(\theta_{0, \mathrm{n}}\right)\right]$ admit Gaussian approximations in the sense of condition (L2) and, in this approximation, the variable T is centered at zero.

$$
\log \frac{d Q_{n}\left(t_{n}\right)}{d Q_{n}\left(\theta_{0, n}\right)}
$$

must also tend to zero. Hence the result.
The result presented here is related to a result of Robert Davies ([8] and [9]). His assumptions differ considerably from ours. Davies is interested in the situation where the restricted experiments $\mathrm{F}_{\mathrm{n}}$ do not necessarily satisfy the LAN condition but only what we have called "quadratic control" by random quadratic forms. This corresponds to the approximability of the log likelihood as in condition (L1) of Section 5 but with random norms $\|\cdot\|_{n}$. He assumes that the conditional distributions $L\left\{W_{n} \mid \mathbf{B}_{n}, P_{n}\left(\theta_{0, n}\right)\right\}$ are approximately Gaussian. To prove this in particular cases he uses local limit theorems and an argument previously used by G. Steck [37].

We have limited ourselves here to the case of control by nonrandom quadratic forms. For that particular case our assumptions are weaker than those made by Davies in that joint asymptotic normality of $W_{n}$ and $T_{n}$ does not necessarily imply asymptotic normality of the conditional distributions $L\left\{W_{n} \mid T_{n}, P_{n}\left(\theta_{0, n}\right)\right\}$.

For instance it is highly visible that the vectors $\mathrm{Z}_{\mathrm{n}}$ used in Theorem 3 and further on to describe the construction of estimates are distinguished. Thus, one could prove the required joint asymptotic normality as in Section 4 and then apply our present Theorem 4 to obtain an alternate proof of Theorem 3, Section 5. This does not require examination of conditional distributions.

However it can also be argued that, under the conditions of our Theorem 4, Davies' conditions are almost satisfied. Indeed it can be shown (see [27], Chapter 7, Section 3) that the assumption that the sequence $\left\{\mathrm{T}_{\mathrm{n}}\right\}$ is distinguished is equivalent to an assumption to the effect that the likelihood ratios

$$
\frac{\mathrm{dQ}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)}{\mathrm{dQ}_{\mathrm{n}}\left(\theta_{0, \mathrm{n}}\right)}
$$

can be approximated by smooth functions of $T_{n}$. To state that more specifically, take a fixed $k$ independent of $n$ and sets $\left\{t_{n, i} ; i=1,2, \ldots k\right\}$ in $\Theta_{n}$. Let $Z_{n}$ be the $k$-dimensional vector formed by the densities of the $Q_{n}\left(t_{n, i}\right)$ with respect to their sum $\mu_{n}$. It is shown in [27], Chapter 7, Section 3, that $\left\{T_{n}\right\}$ is a distinguished sequence if and only if for each set $\left\{t_{n, i} ; i=1,2, . . k\right\}$ there is a fixed finite set $\left\{\gamma_{v} ; v=1,2, \ldots r\right\}$ of continuous functions from $\mathbf{R}^{q}$ to the unit simplex of $\mathbf{R}^{\mathbf{k}}$ such that

$$
\inf _{v} \int\left|Z_{n}-\gamma_{v}\left(T_{n}\right)\right| d \mu_{n}
$$

tends to zero. (Here the norm is the maximum coordinate norm.) The proof of that statement is rather involved. Under the contiguity restrictions used here, one can
replace the sum $\mu_{n}$ by $Q_{n}\left(\theta_{0, n}\right)$. Suppose then in addition that the sets $\Theta_{n}$ are sufficiently rich in that, for instance, they contain a fixed open subset of $V$.

It is then possible, using Laplace transforms, to conclude that the conditional distributions of $W_{n}$ given $T_{n}$ are approximable in $Q_{n}\left(\theta_{0, n}\right)$ probability by equicontinuous functions of the $T_{n}$. If the joint distributions converge to normal limits, Steck's theorems will then imply that these equicontinuous approximations will also converge to normal limits. This is all that would be required for the application of Davies' theorem.

Thus, modulo a rather lengthy and complex argument and for sufficiently rich sets $\Theta_{n}$, one can say that our assumptions imply those made by Davies. It should be noted, however, that to derive Theorems 2 and 3 from Theorem 4 one still needs to carry out a good part of the arguments previously described in Sections 4 and 5. As already pointed out above, to derive Theorem 4 from Davies' results one needs a considerable amount of work. This is also true for Theorems 2 and 3.

A further remark is that the proof of Theorem 4 does not actually depend on the fact that the spaces V and $\mathbf{R}^{q}$ have a fixed finite dimensionality. One can carry out a similar proof in infinite dimensional spaces. This is useful in non-parametric situations and will be published elsewhere.

## 7. Differentiability in quadratic mean and the i.i.d. case.

Let $\Theta$ be a subset of a Euclidean space and let $\{p(\theta): \theta \in \Theta\}$ be a family of probability measures on a $\sigma$-field A. A condition often used to study the i.i.d. case is a condition of differentiability in quadratic mean expressible as follows. At a point $\theta_{0}$ let $\xi^{2}(s), \xi(s) \geq 0$, be the density $\frac{d p(s)}{\operatorname{dp}\left(\theta_{0}\right)}$ of the part of $p(s)$ that is dominated by $p\left(\theta_{0}\right)$.
(DQM) (i) There are random vectors X such that as $|\mathrm{t}| \rightarrow 0$

$$
E_{\theta_{0}} \frac{1}{|t|^{2}}\left|\xi\left(\theta_{0}+t\right)-\xi\left(\theta_{0}\right)-t^{\prime} X\right|^{2} \rightarrow 0
$$

and
(ii) Let $\beta\left(\theta_{0}+\mathrm{t}\right)$ be the mass of $\mathrm{p}\left(\theta_{0}+\mathrm{t}\right)$ that is $\mathrm{p}\left(\theta_{0}\right)$ singular. Then

$$
\frac{1}{|t|^{2}} \beta\left(\theta_{0}+t\right) \rightarrow 0 \text { as }|t| \rightarrow 0
$$

Now replace the $\sigma$-field $\mathbf{A}$ by a smaller $\sigma$-field $\mathbf{B}$ and let $q(\theta)$ be the restriction of $p(\theta)$ to $B$. We shall show that, if $\{p(\theta) ; \theta \in \Theta\}$ satisfies DMQ at $\theta_{0}$, so does $\{\mathrm{q}(\theta) ; \theta \in \Theta\}$. For this we shall use an indirect proof, using the asymptotic properties of the likelihood ratios when the number of observations increases indefinitely. A direct proof can be carried out but it is less informative than the indirect proof for the standard i.i.d. case that can be described as follows.

Let $\mathbf{E}$ be the system of measures $\{p(\theta) ; \theta \in \Theta\}$ on a $\sigma$-field $\mathbf{A}$ carried by a set $\mathbf{X}$. For each integer j , let $\left\{\mathbf{X}_{\mathbf{j}}, \mathbf{A}_{\mathbf{j}}, \mathbf{E}_{\mathrm{j}}\right\}$ be a copy of $\{\mathbf{X}, \mathbf{A}, \mathbf{E}\}$. The experiment carried out at stage $n$ is the direct product $\tilde{\mathbf{E}}_{\mathrm{n}}=\prod_{\mathrm{j}=1}^{\mathrm{n}} \mathbf{E}_{\mathrm{j}}$ of n copies of $\mathbf{E}$. That is, $\tilde{\mathbf{E}}_{\mathrm{n}}$ is the experiment that consists of observing $n$ independent variables $\omega_{1}, \omega_{2}, \cdots, \omega_{\mathrm{n}}$ with common distribution equal to a certain $p\left(\theta_{0}\right), \theta_{0} \in \Theta$. As $n$ tends to infinity we assume that the "true" $\theta_{0}$ stays fixed independent of $n$. (For certain purposes, for instance to prove uniformity of convergences, one makes the "true" value of the parameter depend on n . If so, the situation one encounters is more restricted than the general situation covered by the infinitesimal arrays of Section 4. However, the restrictions are not as drastic as one could presume.)

For the experiments $\mathbf{E}_{\mathrm{j}, \mathrm{n}}=\left\{\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ of Sections 2 to 5 we shall take the copy $\mathbf{E}_{\mathrm{j}, \mathrm{n}}=\left\{\mathrm{p}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ with $\theta$ restricted to lie in a subset $\Theta_{\mathrm{n}}$ of $\Theta$. Finally we shall consider only the case where the $\Theta_{n}$ are of the form $\Theta_{n}(b)=\left\{\theta \in \Theta, \sqrt{n}\left|\theta-\theta_{0}\right| \leq b\right\}$ or perhaps of the form $\Theta_{n}^{\prime}(b)=\left\{\theta \in \Theta, a_{n}\left|\theta-\theta_{0}\right| \leq b\right\}$ for some other coefficients $a_{n}$ that tend to infinity.

It is well known that in such situations the condition DQM implies that the products $\tilde{\mathbf{E}}_{\mathrm{n}}$ satisfy the LAN conditions on sets of the form $\Theta_{n}(b)=\left\{\theta \in \Theta ; \sqrt{n}\left|\theta-\theta_{0}\right| \leq b\right\}$. It is also known (see for instance Example 4, Section 9) that cases occur where DQM is not satisfied but where the LAN conditions are satisfied, for instance on sets $\left\{\theta \in \Theta ; \sqrt{n \log n}\left|\theta-\theta_{0}\right| \leq b\right\}$. Even if one sticks with the sets $\Theta_{n}(b)$ and the speed of convergence $\sqrt{n}$, the validity of the LAN conditions does not imply that of DQM. (An example is given in [23], page 816.) Thus, the preservation of the LAN conditions by passage from the $\sigma$-field $\mathbf{A}$ to a smaller $\sigma$-field $\mathbf{B}$ does not logically imply that DQM will also be preserved. We shall now show that it is preserved.

It is known (see [27], Chapter 17, Section 3) that the DQM requirement is equivalent to the following. For a pair ( $\mathrm{s}, \mathrm{t}$ ) let $\Lambda_{\mathrm{n}}(\mathrm{t}, \mathrm{s})$ be the logarithm of likelihood ratio

$$
\log \frac{\mathrm{dP}_{\mathrm{n}}(\mathrm{t})}{\mathrm{dP}_{\mathrm{n}}(\mathrm{~s})}
$$

for the product measures $P_{n}(s)=\prod_{j} p_{j, s}$ where $p_{j, s}$ is the $j^{\text {th }}$ copy of $p_{s}$.
$(\mathrm{DQM})^{\prime}$ There are random vectors $\mathbf{X}_{\mathrm{j}}$ defined on $\left(\mathbf{X}_{\mathrm{j}}, \mathbf{A}_{\mathrm{j}}\right)$ with the following properties.
i) the $\mathrm{X}_{\mathrm{j}}$ are copies of a given $\mathrm{X}_{1}$
ii) $\quad \mathrm{E}_{\theta_{0}} \mathrm{X}_{\mathrm{j}}=0$ and for $\mathrm{t} \in \mathrm{V}, \mathrm{E}_{\theta_{0}}\left(\mathrm{t}^{\prime} \mathrm{X}_{\mathrm{j}}\right)^{2}<\infty$.
iii) If $\left|t_{n}\right|$ remains bounded then

$$
\Lambda_{n}\left(\theta_{0}+\frac{t_{n}}{\sqrt{n}}, \theta_{0}\right)-t_{n}^{\prime}\left(\frac{1}{\sqrt{n}} \sum_{j=1}^{n} X_{j}\right)+\frac{1}{2} E_{\theta_{0}}\left(t_{n}^{\prime} X_{1}\right)^{2}
$$

tends to zero in $\mathrm{P}_{\mathrm{n}}\left(\theta_{0}\right)$ probability.
In other words, for any fixed $b$ let

$$
\Theta_{n}(b)=\left\{\theta \in \Theta ; \sqrt{n}\left|\theta-\theta_{0}\right| \leq b\right\}
$$

and let $\tilde{E}_{\mathrm{n}}$ be the product experiment indexed by $\Theta_{\mathrm{n}}(\mathrm{b})$. The ( DQM$)^{\prime}$ condition is equivalent to the requirement that the $\tilde{\mathbf{E}}_{\mathrm{n}}$ satisfy the LAN conditions (R1), (R2), (L1), (L2) of Section 5 with the further specification that the linear term $t_{n}{ }^{\prime} W_{n}$ of condition (L1) have the special form $t_{n}{ }^{\prime}\left(\frac{1}{\sqrt{n}} \sum_{j=1}^{n} X_{j}\right)$.

Instead of using the logarithms $\Lambda_{n}\left(\theta_{0}+\frac{t_{n}}{\sqrt{n}}, \theta_{0}\right)$ one can use the sum

$$
S_{n}\left(t_{n}\right)=\sum_{j=1}^{n}\left[\frac{d p_{j, \theta_{n}}}{d p_{j, \theta_{0}}}-1\right]
$$

where $\theta_{n}=\theta_{0}+\frac{t_{n}}{\sqrt{n}}$. The condition is then that

$$
S_{n}\left(t_{n}\right)-t_{n}^{\prime}\left(\frac{1}{\sqrt{n}} \sum_{j=1}^{n} X_{j}\right)
$$

tends to zero in probability.
Now let us consider experiments $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$ that are weaker than our $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$. Passage to weaker experiments need not preserve the i.i.d. structure. For instance it could be obtained by passing from the $\sigma$-field $\mathbf{A}_{\mathrm{j}}$ to a $\sigma$-field $\mathbf{B}_{\mathrm{j}, \mathrm{n}} \subset \mathbf{A}_{\mathrm{j}}$ but in such a way that the $\mathbf{B}_{\mathrm{j}, \mathrm{n}}$ are not copies of each other. An example can be readily constructed by looking at the counting processes described in the Introduction. They might be i.i.d. on the total observation period [0,L] but on the weaker experiment the isolated times at which the patients are observed may depend on the patient.

We shall consider only the situation where there is a $\sigma$-field $\mathbf{B} \subset \mathbf{A}$ with copies $\mathbf{B}_{\mathrm{j}} \subset \mathbf{A}_{\mathrm{j}}$ and where the $\mathrm{j}^{\text {th }}$ experiment $\mathbf{F}_{\mathrm{j}}$ is obtained by restricting the $\mathrm{p}_{\boldsymbol{\theta}}$ of $\mathbf{E}$ to $\mathbf{B}$ and copying the result on $\mathbf{B}_{\mathrm{j}}$.

One immediate result is as follows:
PROPOSITION 4. If the $\mathrm{p}_{\theta} ; \theta \in \Theta$ satisfy DQM at $\theta_{0}$ so do their restrictions $\mathrm{q}_{\theta}$ to the $\sigma$-field $\mathbf{B} \subset \mathbf{A}$.

PROOF. Use the variant of (DQM $)^{\prime}$ that says that $S_{n}\left(t_{n}\right)-t_{n}^{\prime}\left(\frac{1}{\sqrt{n}} \sum_{j=1}^{n} X_{j}\right)$ tends to zero. According to the argument of Theorem 2 the corresponding sums $S_{n}^{*}\left(t_{n}\right)$ for the weaker experiment will be such that $S_{n}^{*}\left(t_{n}\right)-t_{n}^{\prime}\left(\frac{1}{\sqrt{n}} \sum_{j=1}^{n} X_{j}^{*}\right) \rightarrow 0$ for $X_{j}^{*}=E_{\theta_{0}}\left(X_{j} \mid B_{j}\right)$.
(The preservation of DQM is no surprise in view of its relation to Lipschitz conditions or to the rate of separation $\sqrt{n}$. See [23].)

In the present case, since the rate of convergence $\sqrt{n}$ for estimates is about as bad as can be, one can also hope that the weaker experiments will still give suitable estimates. This need not be so. However here is a usable result.

Consider the expectation $\mathrm{E}_{\theta_{0}}\left(\mathrm{t}^{\prime} \mathrm{X}_{1}\right)^{2}$ used in ( DQM$)^{\prime}$. It may be written $\mathrm{t}^{\prime} \Gamma\left(\theta_{0}\right) \mathrm{t}$ for a certain matrix $\Gamma\left(\theta_{0}\right)$ called the Fisher information matrix.

PROPOSITION 5. Assume that the $\mathrm{p}_{\theta}$ satisfy the following requirements.
i) if $\theta_{k} \rightarrow \theta$ then the total variation $\left\|\mathrm{p}_{\theta_{\mathrm{k}}}-\mathrm{p}_{\theta}\right\|$ tends to zero
ii) if $\mathrm{t} \neq \mathrm{s}$ then $\left\|\mathrm{p}_{\mathrm{t}}-\mathrm{p}_{\mathrm{s}}\right\|>0$.

Then there exists estimates $\mathrm{T}_{\mathrm{n}}$ based on n observations such that $\sqrt{\mathrm{n}}\left(\mathrm{T}_{\mathrm{n}}-\theta\right)$ stays bounded in $\mathrm{P}_{\mathrm{n}}(\theta)$ probability for any $\theta \in \Theta$ that satisfies ( $D Q M$ ) with a non-singular information matrix $\Gamma(\theta)$.

PROOF. This is a known result. See for instance [22].
Now consider the weaker experiments with measures restricted to the copies $\mathbf{B}_{\mathrm{j}}$ of a $\sigma$-field $\mathbf{B} \subset \mathbf{A}$. Let $q_{\theta}$ be the restriction of $p_{\theta}$ to $\mathbf{B}$.

It is obvious that whenever condition (i) of Proposition 5 is satisfied by the $p_{\theta}$ it is also satisfied by the $q_{\theta}$. Condition (ii) need not be preserved. That $p_{t} \neq p_{s}$ does not necessarily imply $q_{t} \neq q_{s}$. However, that $s \neq t$ implies $q_{t} \neq q_{s}$ is often easy to check.

Thus we are essentially left with the question of non-singularity of Fisher information matrices. A technique that often works is the following. If the $q_{\theta}$ are complicated, there may be $\sigma$-fields $\mathbf{C} \subset \mathbf{B}$ on which their restrictions $\mathrm{r}_{\theta}$ are quite manageable. It is sufficient to verify the nonsingularity of information matrices for the $r_{\theta}$.

## 8. An extension to a dependent situation.

The proofs of Sections 4 and 5 depend very strongly on the independence assumptions made there. However they suggest a possible extension to a dependent situation.

It often happens that an experiment has the following structure: It is given by measures $P_{n}(\theta)$ on a $\sigma$-field $A_{n}$ but there is a filtration $\mathbf{A}_{1, \mathrm{n}} \subset \mathbf{A}_{2, \mathrm{n}} \subset \cdots \subset \mathbf{A}_{\mathrm{j}, \mathrm{n}} \subset \mathbf{A}_{\mathrm{j}+1, \mathrm{n}} \subset \cdots \subset \mathbf{A}_{\mathrm{n}}$ formed by a finite or infinite increasing family of sub- $\sigma$-fields. Since our purpose is not to attain maximum generality, we shall assume that the filtration is a finite sequence with $\mathbf{A}_{\mathrm{n}}$ being the last one of the sequence.

Except for technical details that are annoying but not so important, one can describe such a situation in a more intuitive manner as follows.

One starts by performing the experiment with the $\mathrm{P}_{\boldsymbol{\theta}}$ restricted to $\mathbf{A}_{1, n}$. If everything has been carried out up to and including the $\mathrm{j}^{\text {th }}$ step one performs an additional experiment $\mathbf{E}_{\mathrm{j}+1, \mathrm{n}}$ where the probability measures depend on the observations seen up to and including the $\mathrm{j}^{\text {th }}$ step. One continues till the end of the sequence has been reached.

To retain a notation similar to that used previously, we shall let $\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta)$ be the conditional distribution used at the $\mathrm{j}^{\text {th }}$ stage when the value of the parameter is $\theta$. It is a function of the previous observations, measurable with respect to $\mathbf{A}_{j-1, n}$. Just as in Section 2 , one can define quantities $h_{j, n}$ by

$$
h_{j, n}^{2}(s, t)=\frac{1}{2} \int\left[\sqrt{d p_{j, n}(s)}-\sqrt{d p_{j, n}(t)}\right]^{2}
$$

for integrals taken conditionally given $\mathbf{A}_{j-1, n}$. Now the $h_{j, n}$ are $\mathbf{A}_{j-1, n}$ measurable functions. The joint distributions will be denoted $\mathrm{P}_{\mathrm{n}}(\theta)$ as before.

By analogy with the requirements imposed in Section 4, we shall assume that the following conditions are satisfied for all possible choices of pairs ( $s_{n}, t_{n}$ ) in $\Theta_{n}$
(A) As $n \rightarrow \infty$ the quantities $\sup _{j} h_{j, n}\left(s_{n}, t_{n}\right)$ tend to zero in $P_{n}\left(s_{n}\right)+P_{n}\left(t_{n}\right)$ probability.
(B) There is a number $b<\infty$, independent of $\left(s_{n}, t_{n}\right)$ such that $\sum_{j} h_{j, n}^{2}\left(s_{n}, t_{n}\right) \leq b$.
(C) For every fixed $\varepsilon>0$ the sums

$$
\sum_{j}\left\|p_{j, n}\left(t_{n}\right)-\left[(1+\varepsilon) p_{j, n}\left(s_{n}\right)\right] \wedge p_{j, n}\left(t_{n}\right)\right\|
$$

tend to zero in probability.
Note that here the conditions (A), (B), (C) are no longer sufficient to imply asymptotic normality of the logarithms

$$
\Lambda_{n}\left(t_{n}, s_{n}\right)=\sum_{j} \log \frac{\operatorname{dp}_{j, n}\left(t_{n}\right)}{d p_{j, n}\left(s_{n}\right)}
$$

As far as we have been able to determine, necessary and sufficient conditions for the asymptotic normality of $\Lambda_{n}\left(\mathrm{t}_{\mathrm{n}}, \mathrm{s}_{\mathrm{n}}\right)$ are not known, even under the restrictive conditions (A), (B). In particular it is not known whether (C) is necessary. We have written it by analogy with the conditions that are necessary and sufficient in the independent case. The literature does contain necessary and sufficient conditions for the asymptotic normality of the likelihood process considered as a function of $\theta$ and of a time parameter corresponding to the index $\mathbf{j}$ of $\mathbf{A}_{\mathrm{j}, \mathrm{n}}$. These are the so-called "invariance principles" for which see Aalen [1], Rebolledo [33] or Shiryayev [35]. However that is a different matter. Here we are interested only on the behavior of the likelihood ratios on the final $\sigma$-field $\mathbf{A}_{\mathrm{n}}$.

Under the conditions (A), (B), (C) the Taylor expansion argument that gives an approximation for the $\log$ likelihood $\Lambda_{n}\left(t_{n}, s_{n}\right)$ is still valid. Therefore it is still true that the difference

$$
\Lambda_{n}\left(t_{n}, s_{n}\right)-\sum_{j}\left[\frac{d p_{j, n}\left(t_{n}\right)}{d p_{j, n}\left(s_{n}\right)}-1\right]+\frac{1}{2} \sum_{j}\left[\frac{d p_{j, n}\left(t_{n}\right)}{d p_{j, n}\left(s_{n}\right)}-1\right]^{2}
$$

tends in probability to zero. This suggests the introduction of quadratic forms $\Gamma_{\mathrm{n}}$ as follows. Let $\mathbf{M}_{\mathrm{n}}$ be the space of finite signed measures with finite support on $\Theta_{\mathrm{n}}$ and total algebraic mass zero as before. Select some $\mathrm{s}_{\mathrm{n}} \in \Theta_{\mathrm{n}}$ and let

$$
\Gamma_{\mathrm{n}}(\mu)=\sum_{\mathrm{j}}\left|\int\left[\frac{\mathrm{dp}_{\mathrm{j}, \mathrm{n}}(\mathrm{t})}{\mathrm{dp}_{\mathrm{j}, \mathrm{n}}\left(\mathrm{~s}_{\mathrm{n}}\right)}-1\right] \mu(\mathrm{dt})\right|^{2}
$$

Consider also the forms

$$
\mathrm{K}_{\mathrm{n}}(\mu)=-4 \Sigma \int \sum_{\mathrm{j}} \int \mathrm{~h}_{\mathrm{j}, \mathrm{n}}^{2}(\mathrm{~s}, \mathrm{t}) \mu(\mathrm{ds}) \mu(\mathrm{dt}) .
$$

THEOREM 5. Let the $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ satisfy the conditions $(A),(B),(C)$. Then the experiment $\overline{\mathrm{E}}_{\mathrm{n}}$ defined by the $\mathrm{E}_{\mathrm{j}, \mathrm{n}}$ is controlled by the quadratic forms $\mathrm{K}_{\mathrm{n}}$. Also if the cardinality of the supports of measures $\mu_{\mathrm{n}}$ remains bounded and if $\sup _{\mathrm{n}}\left\|\mu_{\mathrm{n}}\right\|<\infty$ then $\mathrm{K}_{\mathrm{n}}\left(\mu_{\mathrm{n}}\right)-\Gamma_{\mathrm{n}}\left(\mu_{\mathrm{n}}\right)$ tends to zero in probability.

PROOF. According to condition (C) there is no real loss of generality in assuming that the conditional distributions $\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta), \theta \in \Theta$ are mutually absolutely continuous. A standard truncation argument shows also that one can proceed as if the martingale differences

$$
X_{j, n}\left(t_{n}\right)=\frac{d p_{j, n}\left(t_{n}\right)}{d p_{j, n}\left(s_{n}\right)}-1
$$

were all bounded by unity. The integrals $\int \mathrm{X}_{\mathrm{j}, \mathrm{n}}(\mathrm{t}) \mu_{\mathrm{n}}(\mathrm{dt})$ are also martingale
differences under $\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)$. Their sum will tend to zero in probability if their quadratic variation $\Gamma_{\mathrm{n}}\left(\mu_{\mathrm{n}}\right)$ tends to zero. This quadratic variation can also be replaced by the sum

$$
\sum_{j} E_{j-1}\left[\int X_{j, n}(t) \mu_{n}(d t)\right]^{2}
$$

where $\mathrm{E}_{\mathrm{j}-1}$ denotes an expectation taken conditionally given the past at $\mathrm{j}-1$.
(For these relations, see Neveu [29].)
A Taylor expansion argument shows that if $Y_{j, n}(t)=\sqrt{1+X_{j, n}(t)}-1$ then $X_{j, n}(t)$ can be replaced by $2 \mathrm{Y}_{\mathrm{j}, \mathrm{n}}(\mathrm{t})$ in the preceding formula. Thus we are led to consider the expressions

$$
\begin{aligned}
4 \mathrm{E}_{\mathrm{j}-1}\left[\int Y_{j, n}(t) \mu_{\mathrm{n}}(\mathrm{dt})\right]^{2} & =4 \mathrm{E}_{\mathrm{j}-1} \iint \mathrm{Y}_{\mathrm{j}, \mathrm{n}}(\mathrm{~s}) \mathrm{Y}_{\mathrm{j}, \mathrm{n}}(\mathrm{t}) \mu_{\mathrm{n}}(\mathrm{ds}) \mu_{\mathrm{n}}(\mathrm{dt}) \\
& =4 \iint\left[1-h_{j, n}^{2}(\mathrm{~s}, \mathrm{t})\right] \mu_{\mathrm{n}}(\mathrm{ds}) \mu_{\mathrm{n}}(\mathrm{dt}) \\
& =K_{\mathrm{n}}\left(\mu_{\mathrm{n}}\right)
\end{aligned}
$$

yielding the desired result.
The fact that the experiments $\overline{\mathbf{E}}_{\mathrm{n}}$ are quadratically controlled can often be used to construct estimates that possess asymptotic minimaxity or asymptotic sufficiency properties. A theory to that effect has been expounded in [27], Chapter 11. The technique of [27] requires for its validity a number of additional restrictions such as existence of well behaved auxiliary estimates and dimensionality restrictions on the parameter spaces $\Theta_{\mathrm{n}}$. For such conditions we can only refer the reader to [27] and to the work of Jeganathan [16] [17].

Now let us pass to the situation where there is another filtration $\mathbf{B}_{\mathrm{j}, \mathrm{n}}, \mathrm{j}=1,2, \ldots$ with $\mathbf{B}_{\mathrm{j}, \mathrm{n}} \subset \mathbf{B}_{\mathrm{j}+1, \mathrm{n}} \subset \cdots \subset \mathbf{B}_{\mathrm{n}}$. If $\mathbf{B}_{\mathrm{j}, \mathrm{n}} \subset \mathbf{A}_{\mathrm{j}, \mathrm{n}}$ for every j it may happen that the conditional experiment $F_{j, n}$ carried out at the $j^{\text {th }}$ stage is always weaker than the corresponding $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$. This is far from automatic. It is also plain that conditions such as (A) or (C) have no reason to be inherited by the conditional distributions $q_{j, n}(\theta)$ used on the weaker filtration. Note that the conditional measures $\mathrm{q}_{\mathrm{j}, \mathrm{n}}(\theta)$ are $\mathbf{B}_{\mathrm{j}-1, \mathrm{n}}$-measurable while the $\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta)$ are only $\mathbf{A}_{\mathrm{j}-1, \mathrm{n}}$ measurable. We shall restrict our attention to cases where the experiments $\left\{\mathrm{q}_{\mathrm{j}, \mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ are always weaker than the $\left\{\mathrm{p}_{\mathrm{j}, \mathrm{n}}(\theta) ; \theta \in \Theta_{\mathrm{n}}\right\}$ no matter what may have happened in the passage from $\mathbf{B}_{\mathrm{j}-1, \mathrm{n}}$ to $\mathbf{A}_{\mathrm{j}-1, \mathrm{n}}$. This may seem to be a very strong restriction, and it really is. However there do exist interesting situations of that type.

For an example, consider a finite state Markov process $\{\mathrm{X}(\mathrm{t}): \mathrm{t} \in[0, \mathrm{~L}]\}$ observed during the entire interval $[0, \mathrm{~L}]$. One can obtain the type of structure described above by dividing the interval [0,L] by times $\mathrm{t}_{\mathrm{j}}$ such that
$0=\mathrm{t}_{0}<\mathrm{t}_{1}<\cdots<\mathrm{t}_{\mathrm{j}}<\mathrm{t}_{\mathrm{j}+1}<\cdots<\mathrm{t}_{\mathrm{k}}=\mathrm{L}$. The conditional experiment $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ is then obtained by observing the process during the entire interval $\left(\mathrm{t}_{\mathrm{j}-1}, \mathrm{t}_{\mathrm{j}}\right]$. One obtains a weaker $F_{j, n}$ by observing $X(t)$ only at the end point $t_{j}$.

Another example is that of "aggregated Markov chains" for which see [19] and the references given there. It occurs as follows. Let $\left\{X_{s}\left(t_{j}\right) ; j=1,2, \ldots k ; s=1,2, \ldots m\right\}$ be $m$ independent identically distributed finite state Markov chains.

At time $t_{j}$ the fully informed observer sees the actual states for all the individual processes $X_{s}\left(t_{j}\right)$. The restricted observer looks at each possible state $i$ and is given the number $N\left(t_{j}, i\right)$ of processes that are in state $i$. Here, in spite of the fact that $\mathbf{A}_{j-1, n}$ can be considerably larger than $\mathbf{B}_{j-1, n}$ the fully informed observer can carry out his conditional experiment $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ using only the information in $\mathbf{B}_{\mathrm{j}-1, \mathrm{n}}$ but he keeps track of who goes where while the restricted observer is given only summary information.

In view of such examples the following result may be useful.
THEOREM 6. Let the $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ satisfy conditions (A), (B) and (C). Assume that no matter what was observed in the $\mathbf{E}_{\mathrm{r}, \mathrm{n}}, \mathrm{r} \leq \mathrm{j}-1$ the conditional experiment $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$ is always weaker than $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$. Then the $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$ also satisfy the conditions $(\mathrm{A})$, ( B ) and ( C ). The global experiment $\overline{\mathbf{F}}_{\mathrm{n}}$ defined by the $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$ is under the control of quadratic forms $\mathrm{K}_{\mathrm{n}}^{*}$ such that the differences $\mathrm{K}_{\mathrm{n}}-\mathrm{K}_{\mathrm{n}}^{*}$ are positive semidefinite.

PROOF. This can be proved exactly as for the independent case, using Lemma 2 of Section 3 on the conditional distributions.

A corollary is that if the $\overline{\mathbf{E}}_{\mathrm{n}}$ satisfy the LAMN conditions, then, under the restrictions of Theorem 6, the log likelihoods for $\overline{\mathbf{F}}_{\mathrm{n}}$ will posses the same kind of linearquadratic expansions. All the arguments of Sections 4 and 5 remain applicable, except that we have said nothing at all about asymptotic normality. The construction of estimates described in Section 5 remains feasible. The resulting estimates will still be asymptotically Bayes, asymptotically minimax, etc. as shown by P. Jeganathan [16]. Together with estimates of the quadratic forms, they will be asymptotically sufficient.

As to the matter of asymptotic normality, standard martingale limit theorems (see [33] or [35]) show that the sums $\sum_{j} \int X_{j, n}(t) \mu_{n}(d t)$ used in Theorem 5 will be asymptotically normally distributed if the quadratics $\Gamma_{n}\left(\mu_{n}\right)$ or $K_{n}\left(\mu_{n}\right)$ differ from non random quantities by amounts that tend to zero. It would be convenient if that property was inherited by the weaker $\mathbf{F}_{\mathrm{j}, n}$. However it need not be. To give an example consider the case where one has $n$ independent identically distributed variables $\xi_{1}, \xi_{2}, \cdots, \xi_{n}$ and where $\mathrm{E}_{\mathrm{j}, \mathrm{n}}$ consists of observing $\xi_{\mathrm{j}}$. Let $\mathrm{N}_{\mathrm{n}}$ be a stopping time of the sequence $\xi_{1}, \xi_{2}, \cdots, \xi_{n}$ with $\mathbf{N}_{\mathrm{n}} \leq \mathrm{n}$. Let $\mathbf{F}_{\mathrm{j}, \mathrm{n}}=\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ if $\mathrm{j} \leq \mathbf{N}_{\mathrm{n}}$. If $\mathrm{j}>\mathbf{N}_{\mathrm{n}}$ let $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$ be the trivial experiment where nothing is observed.

Suppose that the $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ satisfy the differentiability in quadratic mean condition of Section 7 with non degenerate derivatives. Then if the $\frac{\mathbf{N}_{n}}{n}$ have distributions that do not degenerate as $\mathrm{n} \rightarrow \infty$, the $\overline{\mathbf{F}}_{\mathrm{n}}$ will not be asymptotically normal.

We have already observed that conditions such as (A) and (C) have no reason to be inherited in the passage from $\sigma$-fields $\mathbf{A}_{\mathrm{j}, \mathrm{n}}$ to smaller $\mathbf{B}_{\mathrm{j}, \mathrm{n}}$. Condition (B) is of a different nature. Something like it is inheritable. Indeed conditions (A) (B) (C) imply the contiguity of pairs $\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{s}_{\mathrm{n}}\right)\right\},\left\{\mathrm{P}_{\mathrm{n}}\left(\mathrm{t}_{\mathrm{n}}\right)\right\}$. It has been shown by P . Greenwood and A . Shiryayev ([12] Theorem 4, page 48). That contiguity is equivalent to the conjunction of two conditions one of which is that the sum $\sum_{j} h_{j, n}^{2}\left(s_{n}, t_{n}\right)$ remains bounded in $P_{n}\left(s_{n}\right)+P_{n}\left(t_{n}\right)$ probability. Since contiguity is inherited by the weaker experiments, if (A) (B) (C) holds for the $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ then the corresponding sum of conditional square Hellinger distances for the weaker experiments $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$ will still be bounded in probability whether or not they satisfy (A), (C) or the conditions of Theorem 6.

In Section 6, we mentioned the work of Robert Davies [8]. He considers a supercritical Galton-Watson branching process where the observer sees for each $\mathrm{j} \leq \mathrm{n}$ the size $Z_{j}$ of the $j^{\text {th }}$ generation. This can correspond to our experiment $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$. Davies studies the properties of the experiment $\mathbf{F}_{\mathrm{n}}$ obtained from the $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$ by introducing stronger experiments $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ as follows. Label the $\mathrm{Z}_{\mathrm{j}-1}$ individuals of the $(\mathrm{j}-1)^{\text {st }}$ generation by integers $\mathrm{i}=0,1,2, \ldots, \mathrm{Z}_{\mathrm{j}-1}$. For individual i , let $\xi_{\mathrm{i}}$ be the size of his progeny in the $\mathrm{j}^{\text {th }}$ generation. Record the $\xi_{i}$ for $\mathrm{i}=0,1, \ldots, \mathrm{Z}_{\mathrm{j}-1}$. Letting j vary from 1 to n this yields our experiment $\overline{\mathbf{E}}_{\mathrm{n}}$. The situation so described is analogous to the situation we mentioned earlier for Markov processes except that here the loss of information from $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ to $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$ is due to a different cause.

However, the theorems given in the present section do not apply directly to Davies' situation. The negligibility condition called (A) here is not satisfied. A fortiori (C) is not satisfied. To see this it is enough to note that for $n$ large the passage from the $(\mathrm{n}-1)^{\text {st }}$ generation to the $\mathrm{n}^{\text {th }}$ one gives a large fraction of the information contained in the entire set of observations from 1 to n . This is so because of the exponential increase in the size of the population.

The reasons for the validity of Davies' procedure were described in Section 6.

## 9. Illustrative examples.

Example 1. The problem that motivated us to write the present paper arose from stochastic modelling of the activity of nerve cells [39]. The membrane of such a cell contains numerous channels responsible for exchanges between the interior of the cell and the surrounding medium. In the experiments reported in [2] a nerve cell was electrically stimulated every second for a total of approximately 500 seconds and the behavior of sodium channels was monitored on a microscopic patch of the membrane. Opening of a sodium channel results in an influx of sodium ions and in a measurable change of the electric potential between the interior and the exterior of the cell.

Experimenters try to isolate patches containing very few channels. The actual number $m$ of such channels on a patch is a number that can be ascertained by neurophysiological means.

The channels can be in several different states. In the experiments described they always start in a closed resting state and end in an inactivated absorbing state. The actual states are unobservable, however every passage through an open state results in a measurable voltage change. For each stimulus, the number of passages through the open state is recorded. The actual activity lasts about 15 milliseconds. It has been modelled by Markov processes as follows.

It is assumed that the various sodium channels in the patch behave independently of one another. For a given stimulus, one channel has a probability p of responding. If it responds, it moves through various states according to a homogeneous Markov process, the total number of passages through the open state being random, distributed according to a geometric distribution. This gives a variable $W$ such that $P[W=0]=1-p$ and $P[W=k]=p \theta^{k-1}(1-\theta)$ for $k \geq 1$. Since there are $m$ channels on the patch the result actually observed at the $j^{\text {th }}$ stimulus is a sum $Y_{j}=\sum_{v=1}^{m} W_{v, j}$ of $m$ independent replicates of W .

The relevant density is given by $P\left[Y_{j}=0\right]=(1-p)^{m}$ and for $k=1,2, \ldots$,

$$
\mathrm{P}\left[\mathrm{Y}_{\mathrm{j}}=\mathrm{k}\right]=\sum_{l=1}^{\mathrm{m}}\binom{\mathrm{~m}}{l}\binom{\mathrm{k}-1}{\mathrm{k}-l} \mathrm{p}^{l}(1-\mathrm{p})^{\mathrm{m}-l} \theta^{\mathrm{k}-l}(1-\theta)^{l}
$$

which is a mixture of negative binomials with binomial weights.
There are three parameters $\mathrm{p}, \boldsymbol{\theta}$ and m . However, for the data described in [2], m was supposed to be equal to 4 . In addition p had been estimated from other considerations and could be taken equal to .4. Thus we were left with the problem of estimating $\theta$. In view of technical complications that arose in the experiments it was decided
to use the conditional probabilities

$$
\mathrm{q}(\mathrm{k}, \theta)=\mathrm{P}\left[\mathrm{Y}_{\mathrm{j}}=\mathrm{k} \mid \mathrm{Y}_{\mathrm{j}} \neq 0\right]
$$

It is a simple matter to obtain an auxiliary estimate $\theta_{\mathrm{n}}^{*}$ by the method of moments and to compute $q\left(k, \theta_{n}^{*}\right)$ and $q\left(k, \theta_{n}^{*}+\frac{1}{\sqrt{n}}\right)$ for all values of $k$ actually observed in the sample. In the particular data set we used, k never exceeded 7.

Let $\mathrm{N}_{\mathrm{k}}$ be the number of times the value k was recorded and let $\mathrm{Z}_{\mathrm{k}}$ be equal to

$$
\mathrm{Z}_{\mathrm{k}}=\frac{\mathrm{q}\left(\mathrm{k}, \theta_{\mathrm{n}}^{*}+\frac{1}{\sqrt{n}}\right)}{\mathrm{q}\left(\mathrm{k}, \theta_{\mathrm{n}}^{*}\right)}-1
$$

Compute the sums $S_{1}=\sum_{k} N_{k} Z_{k}$ and $S_{2}=\sum_{k} N_{k} Z_{k}^{2}$. The recipe suggested in Section 5 yields the estimated value $T_{n}=\theta_{n}^{*}+\frac{1}{\sqrt{n}} \frac{S_{1}}{S_{2}}$.

The operation was duly carried out. Unfortunately a chi-square test showed that the fit was not acceptable. The fit from maximum likelihood was even worse. Thus we carried out the procedure described in Section 5 estimating both p and $\theta$ from a data set of $n=560$ observations. The method of moments gave starting values $\mathrm{p}^{*}=.457$ and $\theta^{*}=.159$. To compute the approximation $\mathrm{L}_{\mathrm{n}}^{(2)}$ of Section 5 one needs to select incremented values $p^{*}+h_{1} / \sqrt{n}$ and $\theta^{*}+h_{2} / \sqrt{n}$. This was done for various values from $h_{i}=-1$ to $h_{i}=+1$. The estimates $\tilde{p}$ so obtained varied only in their fourth decimals. The estimates $\tilde{\theta}$ varied in their third decimals, from . 1635 to .1665. This value $\tilde{\theta}=.1665$ with $\tilde{p}=.4535$ gave the best fit, as judged by chi-square.

For comparison, the maximum likelihood estimates were $\hat{\mathrm{p}}=.4538$ and $\hat{\theta}=.1649$, giving a slightly larger, but acceptable chi-square.

The technique also provides automatically approximate values of the covariance matrix of ( $\tilde{\mathrm{p}}, \tilde{\theta}$ ) from the estimated $\mathrm{M}_{\mathrm{n}}^{-1}$ of Section 5. Here the estimated covariance matrix is

$$
10^{-4}\left\{\begin{array}{rr}
1.88 & -1.42 \\
-1.42 & 4.32
\end{array}\right\}
$$

Example 2. This is similar to Example 1 but intended to show that, often, one must proceed with due caution. Consider a two dimensional parameter $\theta=(\mu, \sigma)$ with $\mu \in(-\infty,+\infty)$ and $\sigma \in(0, \infty)$. Let $f(x, \theta)$ and $g(x, \theta)$ be two densities on the real line and let $\alpha$ be a known number $\alpha \in(0,1)$, say $\alpha=1 / 10$. (The case where $\alpha$ is also unknown is very important. It can be treated by extensions of the arguments given here. However many complications can occur. They would take us too far afield.)

Let us suppose that a refined experimenter could observe $n$ independent pairs $\left(\mathrm{X}_{\mathrm{j}}, \mathrm{I}_{\mathrm{j}}\right)$ where $\mathrm{I}_{\mathrm{j}}$ takes value zero with probability $1-\alpha$ and value 1 with probability $\alpha$.

Assume that if $I_{j}=0$ then $X_{j}$ has the density $f(x, \theta)$. If $I_{j}=1$, then $X_{j}$ has the density $g(x, \theta)$. For this observer, the likelihood function takes the form

$$
L_{n}(\theta)=\prod_{j=1}^{n}\left[(1-\alpha) f\left(x_{j}, \theta\right)\right]^{1-I_{\mathrm{j}}}\left[\alpha g\left(x_{j}, \theta\right)\right]^{\mathrm{I}_{\mathrm{j}}}
$$

If both $f$ and $g$ satisfy the $D Q M$ condition of Section 7 with respect to $\theta$, so does the density of $\left(\mathrm{X}_{\mathrm{j}}, \mathrm{I}_{\mathrm{j}}\right)$.

Now consider a restricted observer who sees only the variables $\mathbf{X}_{\mathrm{j}}$. His likelihood function takes the form

$$
L_{n}^{*}(\theta)=\prod_{j=1}^{n}\left\{(1-\alpha) f\left(x_{j}, \theta\right)+\alpha g\left(x_{j}, \theta\right)\right\}
$$

According to Section 7 if DQM is satisfied for the refined observer it is also satisfied for the restricted one. Thus one can proceed to the construction of estimates as in Section 5 at least whenever the Fisher information matrix does not become singular.

Now consider the two special cases where $f(x, \theta)$ is the ordinary normal

$$
f(x, \theta)=\frac{1}{\sqrt{2 \pi}} \exp \left\{-\frac{1}{2}(x-\mu)^{2}\right\}
$$

and where, for case 1

$$
g_{1}(x, \theta)=\frac{1}{\sigma \sqrt{2 \pi}} \exp \left\{-\frac{1}{2 \sigma^{2}}(x-\mu)^{2}\right\}
$$

while for case $2, \mathrm{~g}_{2}(\mathrm{x}, \theta)$ is the Cauchy

$$
g_{2}(x, \theta)=\frac{1}{\sigma \pi} \frac{1}{1+\sigma^{-2}(x-\mu)^{2}} .
$$

In both cases $\sup _{\theta} L_{n}(\theta)<\infty$ unless $\sum_{j=1}^{n} I_{j}<2$. Thus $\sup _{\theta} L_{n}(\theta)<\infty$ except for cases having a total probability at most

$$
[1+(n-1) \alpha](1-\alpha)^{n-1}
$$

By contrast (see Kiefer and Wolfowitz [18]) $\sup _{\theta} L_{n}^{*}(\theta)=\infty$ always. The infinite value can be achieved by taking $\mu$ equal to any one of the observed $x_{j}$ and taking $\sigma=0$. This does not contradict the inequalities of Section 3 since $E\left\{\sup _{\theta} L_{n}(\theta) \mid x_{1}, \ldots, x_{n}\right\}=\infty$ even though the probability that $\sup _{\theta} L_{n}(\theta)$ be large tends to zero exponentially fast as $\mathrm{n} \rightarrow \infty$.

As a result of this state of affairs one can assert the following: Any procedure that attempts to select, iteratively or not, a value $\hat{\theta}_{n}$ such that $L_{n}^{*}\left(\hat{\theta}_{n}\right)=\sup _{\theta} L_{n}^{*}(\theta)$ will either get trapped at or near a local maximum or it will be such that the $\partial_{n}$ of $\hat{\theta}_{n}=\left(\hat{\mu}_{n}, \hat{\sigma}_{n}\right)$ tends to zero, or it will not achieve any maximum local or not.

This is particularly visible in the Cauchy contamination case (called case 2 above). There, an appealing auxiliary estimate of $\mu$ is the median $\mathrm{m}_{\mathrm{n}}$ of the observations. However, if $n$ is odd, $\sup _{\sigma} L_{n}^{*}\left[\left(m_{n}, \sigma\right)\right]=\infty$. Thus, an iterative procedure that starts with an estimate $\left(m_{n}, \sigma_{n}^{*}\right)$ of $\theta=(\mu, \sigma)$ may readily be trapped into a path where $L_{n}^{*}$ goes to infinity. In spite of this, the arguments of Section 5 remain applicable. One can prevent the auxiliary estimates from searching for singularities in the likelihood function by a variety of procedures: Discretization, smoothing, etc. If so, the one step estimates of Section 5 are guaranteed to work "asymptotically". (Contrary to common advice, it may not be safe to iterate!) This, of course, does not guarantee good behavior for a fixed finite $n$. As is the case for most statistical problems, we do not have any general recipe that always avoids all troubles for all finite $n$. The only general recommendation that comes to mind is to exert due caution and at least check that the estimated model fits the data.

The general theme of the present example can be applied to many cases where the densities are mixtures of smooth families. This remains true even if the mixing proportions are unknown. However, in this more general situation, due care must be exercised in the selection of auxiliary estimates. Also identifiability problems may make life even more complicated (as is the case in the mixture of Gaussian distributions of our case 1. If $\sigma=1$, the mixing proportion $\alpha$ is estimable for the refined observer but not for the restricted one. Neighborhoods of the form $\left\{\alpha|\sqrt{n}| \alpha-\alpha_{0} \mid \leq b\right\}$ are then inappropriate. We shall return elsewhere to what happens on the neighborhoods described in Section 2).

Example 3. Consider, for each integer $\mathrm{j} \leq \mathrm{n}$, an homogeneous Markov process $\left\{Z_{j}(t), t \in[0,1]\right\}$ with three states $S_{1}, S_{2}, S_{3}$ and an infinitesimal generator A. Let the experiment $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ consist of observing the $\mathrm{Z}_{\mathrm{j}}(\mathrm{t})$ for the entire interval $[0,1]$, the $\mathrm{Z}_{\mathrm{j}}$ being mutually independent. Let $F_{j, n}$ be the experiment in which the values of $Z_{j}(t)$ are seen only at times 0 and 1 . Assume for instance that the initial distribution of $Z_{j}(0)$ is uniform over the three states. The problem is to estimate the generator $A$.

The transition probabilities $P\left[Z_{j}(1)=S_{k} \mid Z_{j}(0)=S_{l}\right]$ are the entries in the matrix exponential $\exp \{A\}$. For any fixed matrix $A$ it is fairly easy to compute the exponential $\exp \{\mathrm{A}\}$, but carrying out a maximum likelihood procedure is another matter. However one can look at the matrix $\mathrm{M}_{\mathrm{n}}$ that gives the observed frequencies of transitions, take its logarithm and take the generator $A_{n}^{*}$ closest to that logarithm. If $A$ is
unrestricted, there are six parameters. Application of the method of Section 5 will necessitate computation of $\exp \left\{A^{*}\right\}$ and of six exponentials of the type $\exp \left\{A_{n}^{*}+u_{n, i}\right\}$ for suitably selected deviations $u_{n, i}, i=1, \ldots, 6$.

Often the generator A is assumed to belong to a parametrized family. For instance one could consider a one dimensional family $A(\theta)$ of the type $A(\theta)=(1-\theta) A_{0}+A_{1}$ for fixed generators $A_{0}$ and $A_{1}$ and for $\theta \in[0,1]$. In such a case one would take the $A\left(\theta_{n}^{*}\right)$ closest to $\log M_{n}$ in the model. It would be enough to compute two exponentials of the type $\exp \left\{\mathrm{A}\left(\theta_{\mathrm{n}}^{*}\right)\right\}$ and $\exp \left\{\mathrm{A}\left(\theta_{\mathrm{n}}^{*}+u_{\mathrm{n}}\right)\right\}$.

Here several complications may occur.
The matrices A may have complex eigenvalues. Thus the determination of $\log \mathrm{M}$ is an uncertain affair. It may also be that several different generators yield the same transition probabilities. (See [7] and [36].) Here is an example of a parametrized family $A(\theta)=(1-\theta) A_{0}+A_{1}$ where $\exp \left\{A_{0}\right\}=\exp \left\{A_{1}\right\}$ but where $\exp \{A(r)\} \neq \exp \{A(s)\}$ for any pair ( $r, s$ ), $r \neq s$ in ( 0,1 ). Let $A_{0}=-a I+a M-\frac{3 \pi}{2} \Phi$ and $A_{1}=-a I+a M+\frac{3 \pi}{2} \Phi$ where $a$ is the number $a=\frac{3 \pi \sqrt{3}}{2}$, the matrices $I$ and $M$ are respectively the identity matrix and the matrix whose entries are all equal to $1 / 3$ and the matrix $\Phi$ is

$$
\Phi=\left\{\begin{array}{ccc}
0 & 1 / 2 & -1 / 2 \\
-1 / 2 & 0 & 1 / 2 \\
1 / 2 & -1 / 2 & 0
\end{array}\right\}
$$

The same type of example can be used to illustrate the results of Section 8. Instead of observing $n$ independent replicates $Z_{j}(t), t \in[0,1]$, observe just one process $\mathrm{Z}(\mathrm{t}), \mathrm{t} \in[0, \mathrm{n}]$. Let $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ be the observation of the process in $(\mathrm{j}-1, \mathrm{j}]$ conditionally on the value at $\mathrm{j}-1$. The restricted experiment $\mathrm{F}_{\mathrm{j}, \mathrm{n}}$ consists of observing only $\mathrm{Z}(\mathrm{j})$. The same kind of analysis will apply as long as all the states communicate with one another.

Example 4. Let $\mathrm{X}_{\mathrm{j}}, \mathrm{j}=1,2, \ldots, \mathrm{n}$, be i.i.d. with density $[1-|\mathrm{x}-\theta|]^{+}$on the line. Let $U_{j}, j=1, \ldots, n$, be i.i.d. independent of the $X_{j}$ with a fixed known distribution, say $N(0,1)$. Let $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ consist of observing $\mathrm{X}_{\mathrm{j}}$ and let $\mathbf{F}_{\mathrm{j}, \mathrm{n}}$ consist of observing $\mathrm{X}_{\mathrm{j}}+\mathrm{U}_{\mathrm{j}}$.

For the $\mathrm{E}_{\mathrm{j}, n}$ the LAN conditions are satisfied in neighborhoods of the type $V_{n}=\left\{\sqrt{n \log n}\left|\theta-\theta_{0}\right| \leq b\right\}$. Thus, by Section 5 , the $F_{j, n}$ will also satisfy the LAN conditions in these same neighborhoods. However, this is no great consolation since in these $\mathrm{V}_{\mathrm{n}}$ the product $\overline{\mathrm{F}}_{\mathrm{n}}$ degenerates. The $\overline{\mathrm{F}}_{\mathrm{n}}$ do satisfy the LAN conditions in neighborhoods of the type $\left\{\sqrt{n}\left|\theta-\theta_{0}\right| \leq b\right\}$. This is immediate from the fact that the distributions satisfy the DQM condition of Section 7. However it is not a consequence of the results of Sections $4,5,6$ or 7 . It could be deduced from the arguments of Section 6 if one knew already of the existence of distinguished sequence of estimates for the
local problems.
Example 5 (Method of moments). Let $\omega_{1}, \omega_{2}, \cdots, \omega_{\mathrm{n}}$ be independent random elements with values in a set $\Omega$ and a common distribution $p_{\theta} ; \theta \in \Theta_{n}$. The method of moments consists of selecting a measurable function $\phi_{\mathrm{n}}$ from $\Omega$ to a Euclidean space $\mathbf{R}^{\mathbf{m}}$ and estimating $\theta$ by minimizing the Euclidean norm

$$
\left\|\sum_{j=1}^{n}\left[\phi_{n}\left(\omega_{j}\right)-E_{\theta} \phi_{n}\left(\omega_{j}\right)\right]\right\|
$$

We shall assume that $\Theta_{\mathrm{n}}$ is a convex subset of a Euclidean space $\mathbf{R}^{\mathbf{k}}$ and that the origin of $\mathbf{R}^{k}$ belongs to $\Theta_{\mathrm{n}}$.

The passage from $\omega_{j}$ to the vector $\phi_{\mathrm{n}}\left(\omega_{\mathrm{j}}\right)$ can result in a loss of information. It is subject to the results of Sections 4 and 5 . Let

$$
\begin{aligned}
& P_{\theta, \mathrm{n}}=\mathbf{L}\left[\omega_{1}, \omega_{2}, \cdots, \omega_{\mathrm{n}} \mid \theta\right] \\
& \mathrm{Q}_{\theta, \mathrm{n}}=\mathbf{L}\left[\phi_{\mathrm{n}}\left(\omega_{1}\right), \cdots, \phi_{\mathrm{n}}\left(\omega_{\mathrm{n}}\right) \mid \theta\right] .
\end{aligned}
$$

If $\mathbf{E}_{n}=\left\{P_{\theta, n} ; \theta \in \Theta_{n}\right\}$ is asymptotically Gaussian, so is $F_{n}=\left\{Q_{\theta, n} ; \theta \in \Theta_{n}\right\}$. If $E_{n}$ satisfies the LAN conditions, so will $\mathbf{F}_{\mathrm{n}}$.

The passage from the $n$-tuple $\left\{\phi_{n}\left(\omega_{1}\right), \cdots, \phi_{n}\left(\omega_{n}\right)\right\}$ to the sum $T_{n}=\sum_{j=1}^{n} \phi_{n}\left(\omega_{j}\right)$ can also result in a further loss of information. Here the results of Sections 4 and 5 do not apply. The situation is complex. We shall discuss only a special case for which we need an appropriate affine invariant distance between measures on $\mathbf{R}^{m}$. We shall use the half-space distance. A half-space is any set of the form $\left\{x: t^{\prime} x \geq a\right\}$ or $\left\{x: t^{\prime} x>a\right\}$ with a real and with $t^{\prime}$ in the dual of $\mathbf{R}^{m}$. The half-space norm of a signed measure $\mu$ is

$$
\|\mu\|_{\mathrm{H}}=\sup _{\mathrm{H}}|\mu(\mathrm{H})|
$$

for a supremum taken over all half-spaces.
We shall use the following restrictions
a) The experiments $\mathbf{F}_{\mathrm{n}}$ satisfy the LAN conditions.

Let us call "outliers" any one of the $\phi_{n}\left(\omega_{j}\right)$ selected by the statistician after inspection of the $n$-tuple $\left\{\phi_{\mathrm{n}}\left(\omega_{\mathrm{j}}\right) ; \mathrm{j}=1,2, \ldots, \mathrm{n}\right\}$ by any measurable criterion he pleases.
b) Let $\mathrm{T}_{\mathrm{n}}^{*}$ be the sum $\mathrm{T}_{\mathrm{n}}$ with one outlier removed. Then no matter the choice of outlier the difference

$$
L\left(T_{n}^{*} \mid 0\right)-L\left(T_{n} \mid 0\right)
$$

tends to zero for the half-space distance.

This is a commendable form of robustness. By a result of $P$. Lévy it is equivalent to the assertion that there exist Gaussian measures $G_{0, n}$ such that $\left\|L\left(T_{n} \mid 0\right)-G_{0, n}\right\|_{H}$ tends to zero.

The following condition can always be insured by restricting oneself to an affine subspace.
c) The support of $\mathrm{L}\left(\mathrm{T}_{\mathrm{n}} \mid 0\right)$ is not contained in any proper affine subspace of $\mathbf{R}^{\mathrm{m}}$

If (b) and (c) are satisfied one can renormalize $T_{n}$ so that $G_{0, n}$ is the standard $\mathbf{N}(0, \mathrm{I})$ distribution on $\mathbf{R}^{\mathrm{m}}$. We shall do so.

Let $\Lambda_{n}(\theta)=\log \frac{d Q_{\theta, n}}{d Q_{0, n}}$. Condition (a) implies the existence of random variables $V_{n, j}$, functions of $\phi_{n}\left(\omega_{j}\right)$, and the existence of matrices $M_{n}$ such that

$$
\Lambda_{\mathrm{n}}(\theta)-\theta^{\prime} \sum_{j=1}^{\mathrm{n}} \mathrm{~V}_{\mathrm{n}, \mathrm{j}}+\frac{1}{2} \theta^{\prime} M_{\mathrm{n}} \theta
$$

tends to zero in probability. Furthermore, let $\bar{V}_{n}=\sum_{j=1}^{n} V_{n, j}$. Under conditions (a) (b) (c) with the renormalization indicated the joint distribution $L\left[\bar{V}_{n}, T_{n} \mid 0\right]$ admits a Gaussian approximation for the Lévy or Prokorov distance. Indeed $L\left(\bar{V}_{n} \mid 0\right)$ and $L\left(T_{n} \mid 0\right)$ admit separate Gaussian approximations and the joint distribution $L\left(\bar{V}_{n}, T_{n} \mid 0\right)$ admits an infinitely divisible approximation. It follows that there are matrices $B_{n}$ such that if $\theta_{\mathrm{n}} \in \Theta_{\mathrm{n}}$ then

$$
\left\|L\left[T_{n} \mid \theta_{n}\right]-N\left[B_{n} \theta_{n}, I\right]\right\|_{H}
$$

tends to zero. These matrices can be obtained as covariance matrices in the normal approximations to $L\left[\bar{V}_{n}, T_{n} \mid 0\right]$.

This suggests that the experiments $\bar{F}_{n}=\left\{\bar{Q}_{\theta, n} ; \theta \in \Theta_{n}\right\}$ with $\bar{Q}_{\theta, n}=L\left[T_{n} \mid \theta\right]$ will also be asymptotically Gaussian. Indeed, we do not presently have any examples to the contrary. It does not follow however that if the $\bar{F}_{\mathrm{n}}$ are asymptotically Gaussian they are approximable by the $G_{n}=\left\{G_{\theta, n} ; \theta \in \Theta\right\}$ with $G_{\theta, n}=N\left[B_{n} \theta, I\right]$, as will be shown below.

To enforce such a property one needs to add further restrictions. To state them consider the following possible modifications of $T_{n}$. Let $\varepsilon_{n}$ be independent random vectors independent of $T_{n}$. Let $T_{n}{ }^{\prime}=T_{n}+\varepsilon_{n}$. Also, consider an integer $v(n)$ and let $\mathrm{T}_{\mathrm{n}}^{\prime \prime}$ be obtained as follows. Write the coordinates of $\mathrm{T}_{\mathrm{n}}$ in decimal expansions. Then $T_{n}^{\prime \prime}$ has the same decimals as $T_{n}$ up to $v(n)$. The remaining ones are all put equal to zero. Consider binary experiments $B_{n}=\left\{\bar{Q}_{0, n}, \bar{Q}_{\theta_{\mathrm{n}}, \mathrm{n}}\right\}$, $\quad B_{\mathrm{n}}{ }^{\prime}=\left\{\mathrm{Q}_{0, \mathrm{n}}{ }^{\prime}, \mathrm{Q}_{\theta_{\mathrm{n}}, \mathrm{n}}{ }^{\prime}\right\}$, and $B_{n}^{\prime \prime}=\left\{Q_{0, n}^{\prime \prime}, Q_{\theta_{n}, n}^{\prime \prime \prime}\right\}$ with $\theta_{n} \in \Theta_{n}$ and with $Q_{\theta, n}{ }^{\prime}=L\left[T_{n}^{\prime} \mid \theta\right], Q_{\theta, n}^{\prime \prime}=L\left[T_{n}^{\prime \prime}, \theta\right]$.
d) For any arbitrary choice of noises $\varepsilon_{\mathrm{n}}$ if $\left\|\overline{\mathrm{Q}}_{0, \mathrm{n}}-\mathrm{Q}_{0, \mathrm{n}}{ }^{\prime}\right\|_{\mathrm{H}}$ tends to zero so does the Lévy distance between the distributions $\mathrm{L}\left[\left.\frac{\mathrm{d}_{\theta_{\mathrm{n}}, \mathrm{n}}}{\mathrm{d} \overline{\mathrm{Q}}_{0, \mathrm{n}}} \right\rvert\, 0\right]$ and $\mathrm{L}\left[\left.\frac{\left.\mathrm{dQ}_{\theta_{\mathrm{n}, \mathrm{n}}{ }^{\prime}} \mid 0\right] \text {. }}{\mathrm{dQ}_{0, \mathrm{n}}{ }^{\prime}} \right\rvert\,\right.$
$\mathrm{d}^{\prime}$ ) Assume that $\mathrm{L}\left(\mathrm{T}_{\mathrm{n}} \mid 0\right)$ has be normalized as mentioned above. Then for any choice of cut-off $v(n)$ such that $v(n) \rightarrow \infty$ the Lévy distance between $L\left[\left.\frac{d \overline{\mathrm{Q}}_{\theta_{n}, n}}{d \overline{\mathrm{Q}}_{0, \mathrm{n}}} \right\rvert\, \theta\right]$ and $\mathrm{L}\left[\left.\frac{\mathrm{dQ}_{\theta_{\mathrm{n}}, \mathrm{n}}{ }^{\prime \prime}}{\mathrm{dQ}_{0, \mathrm{n}}^{\prime \prime}} \right\rvert\, 0\right]$ tends to zero.

Conditions (d) or ( $\mathrm{d}^{\prime}$ ) can be taken as expressing the fact that the information contained in $\mathrm{T}_{\mathrm{n}}$ occurs there in a robust sort of a way.

Under conditions (a) (b) (c) (d) the $\mathrm{T}_{\mathrm{n}}$ form a distinguished sequence for the experiments $\overline{\mathbf{F}}_{\mathrm{n}}$. Thus Theorem 4 of Section 6 becomes applicable. The same holds true if (d) is replaced by ( $\mathrm{d}^{\prime}$ ).

This is not difficult to prove. If (d) is used one can find noises $\varepsilon_{\mathrm{n}}$ such that $L\left[T_{n}{ }^{\prime} \mid 0\right]-N(0, I)$ tends to zero in total variation. The same will then be true of $L\left[T_{n}{ }^{\prime} \mid \theta_{n}\right]-N\left(B_{n} \theta_{n}, I\right)$. If ( $d^{\prime}$ ) is used instead one can obtain a similar type of effect by considering discretizations of $\mathrm{N}(0, \mathrm{I})$ that match $\mathbf{L}\left(\mathrm{T}_{\mathrm{n}}^{\prime \prime} \mid 0\right)$ in total variation norm except for a difference that tends to zero as $n \rightarrow \infty$. The argument is given in [27] Theorem 2, Section 3, Chapter 7. Condition (d') suggests a construction to obtain sums $\mathrm{T}_{\mathrm{n}}$ that satisfy (a) (b) (c) and are such that the corresponding experiments $\overline{\mathrm{F}}_{\mathrm{n}}$ are asymptotically Gaussian but significantly stronger than the $\mathbf{G}_{\mathrm{n}}$.

To do this suppose for simplicity that the $\phi_{\mathrm{n}}$ are real valued and that they and the $\mathrm{V}_{\mathrm{n}, \mathrm{j}}$ have been truncated so that $\left|\phi_{\mathrm{n}}\left(\omega_{\mathrm{j}}\right)\right| \leq 1$ and $\left|\mathrm{V}_{\mathrm{n}, \mathrm{j}}\right| \leq 1$. Under (a) (b) (c) and the normalization described earlier, this can be done without modifying the asymptotic behavior of the distributions or experiments concerned.

Let $\tilde{\phi}_{n}\left(\omega_{j}\right)$ be $\phi_{n}\left(\omega_{j}\right)$ with all its decimals beyond the $n^{\text {th }}$ one put equal to zero. Let $\phi_{n}^{+}\left(\omega_{j}\right)=\tilde{\phi}_{n}\left(\omega_{j}\right)+10^{-2 n} V_{n, j}$.

Then the $\operatorname{sum} \mathrm{T}_{\mathrm{n}}^{+}=\sum_{\mathrm{j}=1}^{\mathrm{n}} \phi_{\mathrm{n}}^{+}\left(\omega_{\mathrm{j}}\right)$ defines an experiment $\mathrm{F}_{\mathrm{n}}^{+}$that admits a Gaussian approximation $\mathbf{G}_{\mathrm{n}}^{+}$such that the difference between $\mathbf{F}_{\mathrm{n}}$ and $\mathbf{G}_{\mathrm{n}}^{+}$tends to zero. Indeed the sum $\bar{V}_{n}$ is asymptotically sufficient and distinguished for $F_{n}$. The asymptotic behavior of distributions is not changed for the half-space distance. However $\mathbf{G}_{\mathrm{n}}^{+}$can be considerably more informative than $\mathbf{G}_{\mathrm{n}}$. This can be seen for instance in the gamma density $[\Gamma(\alpha)]^{-1} e^{-x} x^{\alpha-1}$ if one estimates $\alpha$ by $\frac{1}{n} \sum_{j=1}^{n} X_{j}$.

A further remark is as follows. In the method of moments one estimates $\theta$ globally while our conditions (a) and (b) are meant to be applied locally in small sets $\Theta_{n}$. In the local situation it is tempting to replace the expectation $E_{\theta} T_{n}$ by the centers $B_{n} \theta$ of the Gaussian approximations. This can be done under suitable integrability requirements. Then minimizing the Euclidean norm $\left\|T_{n}-B_{n} \theta\right\|$ becomes a linear problem. It will have a point solution $\hat{\theta}_{n}$ only if the dimension $m$ of $T_{n}$ is at least equal to the dimension $k$ of $\theta$ and if $B_{n}$ has full rank. The solution is then $\hat{\theta}_{n}=\left(B_{n}{ }^{\prime} B_{n}\right)^{-1} B_{n}{ }^{\prime} T_{n}$. If $m>k$, passage from $T_{n}$ to $\hat{\theta}_{n}$ will often lose information. If $m=k$ no information is lost. However, unless conditions such as (d) or ( $\mathrm{d}^{\prime}$ ) are satisfied, information may be lost in the passage to approximate solutions. This can happen in particular in the minimization of $\left\|T_{n}-E_{\theta} T_{n}\right\|$.

Example 6 (Grouping data). Let $\mathrm{X}_{\mathrm{j}}, \mathrm{j}=1,2, \ldots \mathrm{n}$, be i.i.d. with densities $\mathrm{f}(\mathrm{x}, \theta)$, $\theta \in \Theta$, on the real line. For $E_{j, n}$ one observes $X_{j}$ itself. For $F_{j, n}$ one sees only the integer $\mathrm{Y}_{\mathrm{j}}$ nearest to $\mathrm{X}_{\mathrm{j}}$. (If $\mathrm{X}_{\mathrm{j}}$ happens to be equal to $\mathrm{k}+\frac{1}{2}$ one takes k or $\mathrm{k}+1$ by tossing a coin.)

Even for fairly simple functions $f(x, \theta)$, integrals of the type $\int_{k-1 / 2}^{k+1 / 2} f(x, \theta) d x$ are usually obnoxious to deal with. However, let us assume for instance that $\Theta=\mathbf{R}^{1}$ and that the $f(x, \theta)$ satisfy DQM of Section 7. It is often fairly easy to verify that the $F_{j, n}$ satisfy the identifiability conditions of Proposition 5, Section 7. Assuming this, and assuming that one can find suitable auxiliary estimates $\theta_{k^{7} 1 / 2}^{*}$ the method described in Section 5 will require only computation of the integrals $\int_{k-1 / 2}^{k+1 / 2} f(x, s) d x$ for $q+1$ values of s and for those k that have been seen in the sample. Note that the results remain applicable to situations where the grouping classes are made dependent on the number n of observations.

For grouping in classes that are random, dependent on the observations, other arguments are required.

Example 7 (Censoring). Consider independent pairs ( $\mathrm{X}_{\mathrm{j}}, \mathrm{Y}_{\mathrm{j}}$ ) of positive random variables with a joint density $f(x, y, \theta)$ that depends on a Euclidean parameter $\theta$. The refined experiment $\mathbf{E}_{\mathrm{j}, \mathrm{n}}$ consists of observing the pair ( $\mathrm{X}_{\mathrm{j}}, \mathrm{Y}_{\mathrm{j}}$ ). In the restricted experiment $F_{j, n}$ one sees only $\min \left(X_{j}, Y_{j}\right)$ and an indicator variable $I_{j}$ equal to zero if $X_{j}>Y_{j}$ and to unity if $X_{j} \leq Y_{j}$. Assume that the $E_{j, n}$ yield a product experiment $\tilde{E}_{\mathrm{n}}$ that satisfies the LAN conditions on sets of the form $\left\{\theta: \sqrt{n}\left|\theta-\theta_{0}\right| \leq b\right\}$. Then the restricted product $\tilde{\mathbf{F}}_{\mathrm{n}}$ will also satisfy the same LAN condition. If in addition the densities $f(x, y, \theta)$ satisfy the differentiability in quadratic mean condition of Section 7, the same will be true of the distributions of $\left\{\min \left(X_{j}, Y_{j}\right), I_{j}\right\}$.

In such a situation it is well known (see [30] and [38]) that a $\theta$ that is identifiable on the complete experiment may become unidentifiable on the restricted one.

However most of the examples discussed in the literature satisfy conditions analogous to the following.

1) The densities $f(x, y, \theta)$ have the form $f(x, y, \theta)=g(x, \theta) h(y, \theta)$ so that $X_{j}$ and $Y_{j}$ are independent.
2) The supports of $X_{j}$ and $Y_{j}$ are the same.
3) Let $G(x, \theta)=\int_{x}^{\infty} g(t, \theta) d t$ and $H(y, \theta)=\int_{y}^{\infty} h(t, \theta) d t$ be the survival functions attached respectively to $g$ and $h$. There are two positive numbers $\alpha$ and $\beta$ such that $\left|\theta_{1}-\theta_{2}\right|>\alpha$ implies $\sup _{\mathrm{x}}\left|\mathrm{G}\left(\mathrm{x}, \theta_{1}\right)-\mathrm{G}\left(\mathrm{x}, \theta_{2}\right)\right|>\beta$.
4) The family $\{\mathrm{g}(\mathrm{x}, \theta) ; \theta \in \Theta\}$ satisfies DQM with nonsingular covariance matrices.

Under such conditions, $\theta$ is identifiable. Also one can get auxiliary estimates $\theta_{\mathrm{n}}^{*}$ that converge at the $\sqrt{n}$ rate. A possibility is to start with some nonparametric estimate $\hat{\mathrm{G}}_{\mathrm{n}}$ of the survival function, for instance the Kaplan-Meier estimate or the NelsonAalen estimate from cumulative hazard functions. Then apply a minimum distance technique as in the procedure that yields Proposition 5 of Section 7.

A direct verification that nonsingularity of Fisher information matrices is inherited by the weaker experiment seems awkward. However, this nonsingularity is implied by the properties of the estimate just described.
(In the most standard case, where $h$ does not depend on $\theta$, verification is trivial.)
Thus the results of Section 5 can be applied here. The density of $\mathrm{Z}_{\mathrm{j}}=\min \left(\mathrm{X}_{\mathrm{j}}, \mathrm{Y}_{\mathrm{j}}\right)$ has the form

$$
g(z, \theta) H(z, \theta) \text { on } I_{j}=1
$$

and

$$
\mathrm{h}(\mathrm{z}, \theta) \mathrm{G}(\mathrm{z}, \theta) \text { on } \mathrm{I}_{\mathrm{j}}=0
$$

Assuming that these formulas can be computed for the observed values $\left(\mathrm{Z}_{\mathrm{j}}, \mathrm{I}_{\mathrm{j}}\right)$ and a few values of the parameter $\theta$, Section 5 will yield asymptotically efficient estimates.

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