Distinguished Statistics, Loss of information and a theorem of Robert B. Davies (Fourth edition)

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Introduction. This is a sequel to a paper [1] by R. B. Davies and a paper [4] be the authors. In [1] Davies addresses himself to a question that can be paraphrased as follows: Let $\mathbf{E}_{\mathrm{n}}=\left\{\mathrm{P}_{\theta, \mathrm{n}} ; \theta \in \Theta_{\mathrm{n}}\right\}$ be experiments given by measures $\mathrm{P}_{\theta, \mathrm{n}}$ on $\sigma$-fields $\mathbf{A}_{\mathrm{n}}$. Let $\mathbf{F}_{\mathrm{n}}$ be the experiment formed by the restrictions $\mathrm{Q}_{\theta, \mathrm{n}}$ of the $P_{\theta, n}$ to a sub-field $\mathbf{B}_{n} \subset \mathbf{A}_{n}$. Suppose that $\Theta_{n}$ is a subset of a given Euclidean space and that the $E_{n}$ satisfy the so called LAN assumptions.

What other conditions are needed to insure that the $\mathbf{F}_{\mathrm{n}}$ also satisfy the LAN assumptions, or at least have log likelihoods approximable by a sum of two random terms, the first linear and the second quadratic in $\theta$ ?

Davies gives sufficient conditions. They involve an asymptotic normality requirement for the conditional distributions given $B_{n}$ of the $\log$ likelihoods of the $P_{\theta, \mathrm{n}}$.

In [4] we treated cases where the observations $\left\{x_{1}, \cdots, x_{n}\right\}$ are independent or Markovian and where the passage to subfields is carried out on the individual $x_{j}$ 's. No conditional distributions were involved. However the technique could not yield the result obtained by Davies in the application to branching processes of his Theorem 4.1.

Since the study of approximations for conditional distributions is often delicate we took another look at the situation to see, among other things, whether one could use instead approximations for joint distributions. More precisely we looked for conditions under which one could take limits of joint distributions and then look at conditional distributions in these limits.

It turns out that this can be done if the statistics used are "distinguished" according to a definition of [3], Chapter 7. This property is often, but not always, readily verifiable. It held trivially for the likelihood ratio approximations of [4]. It must hold for the pair ( $\mu_{n}, \tau_{n}$ ) of Theorem 4.1 of [1] and, as shown in Section 3 below, can be checked without much effort for the branching process case of [1].

Another aspect of the situation is that both [1] and [4] rely heavily on contiguity restrictions. These may be quite appropriate for the restricted experiments $F_{n}$ on the subfields, but they cannot be defended easily for the parent experiments $\mathbf{E}_{\mathrm{n}}$. A case in point would be the estimation of variances in the branching processes of [1], as will be shown in Section 3.

The paper is organized as follows. Section 2 contains our main result on passages to subfields in the form of Theorem 2. It says that, under appropriate restrictions on the statistics involved, one can take limits $\left\{\mathrm{D}_{\theta, \infty} ; \theta \in \Theta\right\}$ of families of distributions $\left\{\mathrm{D}_{\theta, \mathrm{n}} ; \theta \in \Theta\right\}$ and deduce the asymptotic behavior of the likelihood ratios of the $D_{\theta, n}$ from that of the likelihod ratios for the $D_{\theta, \infty}$. The passage to
subfields can also be carried out on the $D_{\theta, \infty}$. It is also shown (Lemma 2) that a requirement similar to that imposed in Davies' Theorem 4.1 on conditional distributions does imply the requirements of our Theorem 2.

The results do not make any reference to asymptotic normality or contiguity. They can be used more generally.

Section 3 returns to the branching processes. It uses arguments similar to those of Duby and Rouault in [2]. However these authors needed to let the size of the initial population increase to infinity. We show briefly that their results remain valid "locally" with initial population size $M_{0}=1$ and that the statistics of [1] are "distinguished" there. It is also shown that an apparently minor departure from the assumptions makes them lose this property. As already mentioned, the inappropriateness of contiguity restrictions on the parent experiments $\mathbf{E}_{\mathrm{n}}$ is indicated.

Further applications of Theorem 2 will await another publication.

## 2. Stable and distinguished sequences of statistics.

The objects named in the title are described in some detail in [3], chapter 7, pages $99-108$. After recalling their definitions we proceed to give a result relative to the possibility of interchanging two limit procedures as described in the introduction.

Let $\mathbf{E}=\left\{\mathrm{P}_{\theta} ; \theta \in \Theta\right\}$ be an experiment given by probability measures $\mathrm{P}_{\theta}$ on a $\sigma$-field $\mathbf{A}$ of subsets of a set $\Omega$. Let $\mathbf{X}$ be a compact space. By "statistics with values in X"' will be meant a Markov kernel $\omega \rightarrow \mathbf{f}_{\omega}$ that assigns to each $\omega \in \Omega$ a probability measure $f_{\omega}$ on $\mathbf{X}$. These probability measures will be identified with positive linear functionals on the space $C(\mathbf{X})$ of continuous real functions on $\mathbf{X}$. This is sufficient for all purposes if the family $\mathbf{E}$ is dominated. Otherwise we shall extend the definition, admitting as statistics the more general "transitions" of [3]. However the main arguments involve only finite sets $\Theta$. Hence no difficulty will occur.

The restriction to compact $\mathbf{X}$ is a matter of convenience. In applications, $\mathbf{X}$ will often be Euclidean and not compact. One can always compactify as follows:: Take a separating algebra $\Gamma$ of bounded numerical functions on $\mathbf{X}$, imbed $\mathbf{X}$ in the cartesian product $\mathbf{R}^{\Gamma}$ and take its closure there. (The choice of $\Gamma$ does matter).

In the rest of this section the objects $\mathbf{X}$ and $\Theta$ will be kept fixed. One will consider nets or sequences $\left\{\mathbf{E}_{\nu}\right\}$ of experiments $\mathbf{E}_{\nu}=\left\{\mathrm{P}_{\theta, \nu} ; \theta \in \Theta\right\}$ and statistics $S_{\nu}$ with values in $\mathbf{X}$. This gives rise to families of distributions $\mathbf{D}_{\nu}=\left\{\mathrm{D}_{\theta, \nu} ; \theta \in \Theta\right\}$ where $\mathrm{D}_{\theta, \nu}$ is the distribution of $\mathrm{S}_{\nu}$ under $\mathrm{P}_{\theta, \nu}$, that is the image $\mathrm{S}_{\nu} \mathrm{P}_{\theta, \nu}$ of $\mathrm{P}_{\theta, \nu}$ by the kernel $\mathrm{S}_{\nu}$. The expectation of a $\gamma \in \mathrm{C}(\mathbf{X})$ for this distribution will be written $\gamma \mathrm{D}_{\theta, \nu}$.
Definition 1. Let $\left\{\mathbf{D}_{\nu}\right\}$ be indexed by the directed set $\{\nu\}$. The $\mathbf{D}_{\nu}$ converge vaguely to $\mathbf{D}_{\infty}=\left\{\mathrm{D}_{\theta, \infty} ; \theta \in \Theta\right\}$ if for every $\theta \in \Theta$ and $\gamma \in \mathrm{C}(\mathbf{X})$ the expectations $\gamma \mathrm{D}_{\theta, \nu}$ converge to $\gamma \mathrm{D}_{\theta, \infty}$.

Note that, for the vague convergence, the set of all families of distributions indexed by $\Theta$ is a compact space. Thus every directed set of such families will have cluster points. That is why we have assumed compactness for $\mathbf{X}$.

A family $\mathrm{D}_{\nu}=\left\{\mathrm{D}_{\theta, \nu} ; \theta \in \Theta\right\}$ is also an experiment in its own right. In [3], we introduced a concept of "stability" for nets $\left\{D_{\nu}\right\}$ in terms of a correspondence between vague convergence as families of distributions and weak convergence as experiments. It is a property that is not very strong in appearance, but is shown in [3] to be equivalent to the following approximation property, to be called finite approximability.

Let $\{\mathrm{s}, \mathrm{t}\}$ be a two point subset of $\Theta$ and let $\mu_{\nu}=\mathrm{D}_{\mathrm{s}, \nu}+\mathrm{D}_{\mathrm{t}, \nu}$. Let $\mathrm{f}_{\nu}$ be the Radon-Nikodym density of $\mathrm{D}_{\mathrm{t}, \nu}$ with respect to $\mu_{\nu}$.
Definition 2. The net $\left\{D_{\nu}\right\}$ is stable if for every two point set $\{\mathrm{s}, \mathrm{t}\} \subset \Theta$ and every $\epsilon>0$, there is a finite subset $\left\{\gamma_{j} ; \mathrm{j} \in J\right\}$ of $\mathrm{C}(\mathbf{X})$ and a $\nu(\epsilon)$ such that $\nu \geq \nu(\epsilon)$ implies

$$
\inf _{\mathrm{j} \in \mathrm{~J}} \int\left|\mathrm{f}_{\nu}-\gamma_{\mathrm{i}}\right| \mathrm{d} \mu_{\nu}<\epsilon
$$

A formulation involving finite subsets of $\Theta$ instead of two point sets in given in [3], Theorem 1, page 101, together with several other equivalent conditions. For the two points sets $\{\mathrm{s}, \mathrm{t}\}$ another equivalent statement is as follows
Lemma 1. The net $\left\{\mathrm{D}_{\mathrm{s}, \nu}, \mathrm{D}_{\mathrm{t}, \nu}\right\}$ is stable if and only if for every subnet $\{\nu(\xi)\}$ such that $\left[\mathrm{D}_{\mathrm{s}, \nu(\xi)}, \mathrm{D}_{\mathrm{t}, \nu(\xi)}\right]$ converges vaguely to a limit $\left[\mathrm{D}_{\mathrm{s}, \infty}, \mathrm{D}_{\mathrm{t}, \infty}\right]$ one has

$$
\liminf \int \mathrm{f}_{\nu(\xi)}\left(1-\mathrm{f}_{\nu(\xi)}\right) \mathrm{d} \mu_{\nu(\xi)} \geq \int \mathrm{f}_{\infty}\left(1-\mathrm{f}_{\infty}\right) \mathrm{d} \mu_{\infty}
$$

with $\mu_{\infty}=\mathrm{D}_{\mathrm{s}, \infty}+\mathrm{D}_{\mathrm{t}, \infty}$ and $\mathbf{f}_{\infty}=\frac{\mathrm{dD}_{\mathrm{t}, \infty}}{\mathrm{d} \mu_{\infty}}$.
For a proof, see [3], Chapter 7 and below. (An equivalent statement can also be written in terms of Hellinger affinities $\left.\int \sqrt{\mathrm{f}_{\nu}\left(1-\mathrm{f}_{\nu}\right)} \mathrm{d} \mu_{\nu}\right)$.

To state our next result, it will be convenient to look at a particular pair $\left[\mathrm{D}_{0, \nu}, \mathrm{D}_{1, \nu}\right.$ ] and consider that the elements of $\mathbf{X}$ are random variables $\mathrm{X}_{\nu}$ with possible distributions $D_{i, \nu}, i=0,1$. Each $D_{i, \nu}$ yields a joint distribution for the pair ( $\mathrm{f}_{\nu}, \mathrm{X}_{\nu}$ ). Convergence of such joint distributions will be the vague convergence on $C([0,1] \times \mathbf{X})$.
Theorem 1. Let $\left[\mathrm{D}_{0, \nu}, \mathrm{D}_{1, \nu}\right]$ converge vaguely to $\left[\mathrm{D}_{0, \infty}, \mathrm{D}_{1, \infty}\right]$. Let $\mu_{\nu}=\mathrm{D}_{0, \nu}+\mathrm{D}_{1, \nu}$ and $\mathrm{f}_{\nu}=\mathrm{dD}_{1, \nu} / \mathrm{d} \mu_{\nu}$. The following conditions are all equivalent:

1) $\quad \lim _{\nu} \inf \int \mathrm{f}_{\nu}\left(1-\mathrm{f}_{\nu}\right) \mathrm{d} \mu_{\nu} \geq \int \mathrm{f}_{\infty}\left(1-\mathrm{f}_{\infty}\right) \mathrm{d} \mu_{\infty}$.
2) For every $\epsilon>0$ there is a continuous function $\phi$ such that $0 \leq \phi \leq 1$ and such that $\int\left|\mathrm{f}_{\nu}-\phi\right| \mathrm{d} \mu_{\nu}<\epsilon$ for all sufficiently large $\nu$.
3) The joint distribution of $\left(f_{\nu}, X_{\nu}\right)$ under $D_{i, \nu}, i=0,1$, converges to that of $\left(f_{\infty}, X_{\infty}\right)$ for $D_{i, \infty}$.
If any one of those three conditions hold and if, in addition, the set of discontinuties of $\mathbf{f}_{\infty}$ has $\boldsymbol{\mu}_{\infty}$ measure zero, then $\int\left|\mathrm{f}_{\nu}-\mathrm{f}_{\infty}\right| \mathrm{d} \boldsymbol{\mu}_{\nu}$ tends to zero.

Proof. Consider a decision problem in which it is desired to select a point $z$ in the interval $[0,1]$ with a loss function $\mathrm{W}_{\mathrm{i}}(\mathrm{z})=(\mathrm{i}-\mathrm{z})^{2}$. A decision procedure is a Markov kernel $K(d z, x)$. For $\left[D_{0, \nu}, D_{1, \nu}\right]$ it has a risk sum equal to

$$
\iint\left[\mathrm{z}^{2}\left(1-\mathrm{f}_{\nu}\right)+(1-\mathrm{z})^{2} \mathrm{f}_{\nu}\right] \mathrm{K}(\mathrm{~d} \mathrm{z}, \mathrm{x}) \mu_{\nu}(\mathrm{dx})=\int \mathrm{f}_{\nu}\left(1-\mathrm{f}_{\nu}\right) \mathrm{d} \mu_{\nu}+\iint\left(\mathrm{z}-\mathrm{f}_{\nu}\right)^{2} \mathrm{~K}(\mathrm{dz}, \mathrm{x}) \mu_{\nu}(\mathrm{dx})
$$

It follows that the minimum risk sum is $\mathrm{r}_{\nu}=\int \mathrm{f}_{\nu}\left(1-\mathrm{f}_{\nu}\right) \mathrm{d} \mu_{\nu}$, attained by the nonrandomized estimate $\mathrm{f}_{\nu}$.

Let $\phi \in C(\mathbf{X})$ be such that $0 \leq \phi \leq 1$. It can be used to estimate $z$ yielding a risk sum

$$
\begin{aligned}
\mathrm{R}_{\nu}(\phi) & =\int \phi^{2} \mathrm{dD}_{0, \nu}+\int(1-\phi)^{2} \mathrm{dD}_{1, \nu} \\
& =\mathrm{r}_{\nu}+\int\left|\mathrm{f}_{\nu}-\phi\right|^{2} \mathrm{~d} \mu_{\nu}
\end{aligned}
$$

Do this also for the limit pair $\left[\mathrm{D}_{0, \infty}, \mathrm{D}_{1, \infty}\right]$ obtaining a risk sum

$$
\mathrm{R}_{\infty}(\phi)=\mathbf{r}_{\infty}+\int\left|\mathrm{f}_{\infty}-\phi\right|^{2} \mathrm{~d} \mu_{\infty}
$$

Note that, according to the first expression for $R_{\nu}(\phi)$ one will have $\lim _{\nu} \mathrm{R}_{\nu}(\phi)=\mathrm{R}_{\infty}(\phi)$ for $\phi \in \mathrm{C}(\mathbf{X})$. Now select a $\phi \in \mathrm{C}(\mathbf{X}), 0 \leq \phi \leq 1$, so that $\int\left|\mathrm{f}_{\infty}-\phi\right|^{2} \mathrm{~d} \mu_{\infty}<\epsilon^{2}$. One can write

$$
\int\left|\mathrm{f}_{\nu}-\phi\right|^{2} \mathrm{~d} \mu_{\nu}=\left[\mathrm{R}_{\nu}(\phi)-\mathrm{R}_{\infty}(\phi)\right]-\left(\mathrm{r}_{\nu}-\mathbf{r}_{\infty}\right)+\int\left|\mathrm{f}_{\infty}-\phi\right|^{2} \mathrm{~d} \mu_{\infty}
$$

Since $R_{\nu}(\phi)-R_{\infty}(\phi) \rightarrow 0$ and since $\liminf \left(r_{\nu}-r_{\infty}\right) \geq 0$, this implies that for $\nu$ large one will have $\left.\int \mid \mathrm{f}_{\nu}-\phi\right)^{2} \mathrm{~d} \mu_{\nu}<\epsilon^{2}$. Thus $(1) \Longrightarrow(2)$. To show that $(2) \Longrightarrow(3)$ it is enough to show that for all $\gamma_{1} \in C[0,1]$ and $\gamma_{2} \in C(\mathbf{X})$, the expectations $\mathrm{E}_{\mathrm{i}, \nu} \gamma_{1}\left(\mathrm{f}_{\nu}\right) \gamma_{2}\left(\mathrm{X}_{\nu}\right)$ for $\mathrm{D}_{\mathrm{i}, \nu}$ converge to the corresponding expectations for $\mathrm{D}_{\mathrm{i}, \infty}$. However, take a $\phi$ as done above and consider $\mathrm{E}_{\mathrm{i}, \nu}\left[\gamma_{1}\left(\mathrm{f}_{\nu}\right)-\gamma_{1}(\phi)\right] \gamma_{2}\left(\mathrm{x}_{\nu}\right)$. This does not exceed

$$
\| \gamma_{2}| | \int\left|\gamma_{1}(\phi)-\gamma_{1}\left(\mathrm{f}_{\nu}\right)\right| \mathrm{d} \mu_{\nu} \leq \mu_{\nu}\left[\left|\mathrm{f}_{\nu}-\phi\right|>\epsilon^{1 / 2}\right]+\operatorname{Osc}\left(\gamma_{1}, \epsilon\right)
$$

where $\operatorname{Osc}\left(\gamma_{1}, \epsilon\right)$ is the maximum oscillation of $\gamma_{1}$ on sets of diameter $\epsilon^{1 / 2}$.
By Chebyshev's inequality, the first term does not exceed $\epsilon$ for all $\nu$ sufficiently large. By taking $\epsilon$ sufficiently small, one can make $\operatorname{Osc}\left(\gamma_{1}, \epsilon\right)$ as small as one please. Thus $\int\left|\gamma_{1}\left(\mathrm{f}_{\nu}\right)-\gamma_{1}(\phi)\right| \mathrm{d} \mu_{\nu}$ can be made arbitrarily small. The same applies to $\int\left|\gamma_{1}\left(\mathrm{f}_{\infty}\right)-\gamma_{1}(\phi)\right| \mathrm{d} \mu_{\infty}$ and the conclusion follows by the triangle inequality.

That (3) implies (1) is immediate. It even implies $\lim r_{\nu}=r_{\infty}$. Hence the equivalence of (1) (2) and (3).

To obtain the final statement, note that if the set of discontinuities of $f_{\infty}$ has $\mu_{\infty}$ measure zero, then $\int\left|\mathrm{f}_{\infty}-\phi\right| \mathrm{d} \mu_{\nu}$ tends to $\int\left|\mathrm{f}_{\infty}-\phi\right| \mathrm{d} \mu_{\infty}<\epsilon$. Thus $\int\left|\mathrm{f}_{\nu}-\mathrm{f}_{\infty}\right| \mathrm{d} \mu_{\nu} \leq \int\left|\mathrm{f}_{\nu}-\phi\right| \mathrm{d} \mu_{\nu}+\int\left|\phi-\mathrm{f}_{\infty}\right| \mathrm{d} \mu_{\nu}$ will eventually be inferior to $2 \epsilon$.

This concludes the proof of the theorem.
Remark. Life would be more pleasant if, under the equivalent conditions (1) (2) (3) of the theorem, the integrals $\int\left|\mathrm{f}_{\nu}-\mathrm{f}_{\infty}\right| \mathrm{d} \mu_{\nu}$ would always converge to zero. This is not the case as shown by the following example.

Let $\psi$ be a function defined on $[0,1]$, equal to zero at zero and to $\psi(x)=\sin \left(\frac{c}{x}\right)$ for $x \in(0,1]$. Select $c$ so that $\int_{0}^{1} \sin \left(\frac{c}{x}\right) d x=0$. This is possible. There is a value close to $2 \pi$ with that property.

Let $\delta_{0}$ be the probability measure concentrated at zero and let $\lambda$ be the Lebesgue measure on $[0,1]$. Let $\mu_{\infty}=\delta_{0}+\lambda$ and let $\mathrm{f}_{\infty}=\frac{1}{2}+\frac{1}{4} \psi$. For integer values $\nu$, let $\epsilon_{\nu}=\frac{\mathrm{C}}{2 \pi \nu}$. Take $\phi_{\nu}(\mathrm{x})=\frac{1}{2}$ for $0 \leq \mathrm{x} \leq \epsilon_{\nu}$ and let $\phi_{\nu}(\mathrm{x})=\mathrm{f}_{\infty}(\mathrm{x})$ for $\epsilon_{\nu}<\mathrm{x} \leq 1$. Let $\mu_{\nu}=\lambda_{\nu}+\lambda$ where $\lambda_{\nu}$ is the Lebesgue measure on $\left[0, \epsilon_{\nu}\right]$ multiplied by $\epsilon_{\nu}^{-1}$. Then $\int \phi_{\nu} \mathrm{d} \mu_{\nu}$ may differ from unity but one can make an $\mathrm{f}_{\nu}$ such that $\int \mathrm{f}_{\nu} \mathrm{d} \mu_{\nu}=1$ by letting $\mathrm{f}_{\nu}=\phi_{\nu}$ on $\left[0, \epsilon_{\nu}\right]$ and $\mathrm{f}_{\nu}=\mathrm{k}_{\nu} \phi_{\nu}$ in $\left(\epsilon_{\nu}, 1\right]$ for constants $\mathrm{k}_{\nu}$ that tend to unity. Then the conditions (1) (2) (3) of the theorem are satisfied. However

$$
\int\left|\mathrm{f}_{\nu}-\mathrm{f}_{\infty}\right| \mathrm{d} \mu_{\nu} \geq \frac{1}{4} \int_{1}^{\infty}\left|\sin \left(\frac{2 \pi \nu \mathrm{y}}{\mathrm{c}}\right)\right| \frac{\mathrm{dy}}{\mathrm{y}^{2}}
$$

does not tend to zero.
To go further we need an additional definition as follows.
Consider experiments $\mathbf{E}_{\nu}$ as before with statistics $S_{\nu}$ whose distributions from families $\mathbf{D}_{\nu}$.
Definition 3. The statistics $S_{\nu}$ are pairwise asymptotically sufficient if for every two point set $\{\mathrm{s}, \mathrm{t}\} \subset \Theta$ the distance between the binary experiments $\left\{\mathrm{P}_{\mathrm{s}, \nu}, \mathrm{P}_{\mathrm{t}, \nu}\right\}$ and $\left\{\mathrm{D}_{\mathrm{s}, \nu}, \mathrm{D}_{\mathrm{t}, \nu}\right\}$ tends to zero as $\nu$ increases indefinitely. The $\mathrm{S}_{\nu}$ are called distinguished, for $\mathbf{E}_{\nu}$, if they are pairwise asymptotically sufficient and if their distributions are stable.

With this we can pass to our main goal which is to obtain relations between pairs of statistics $\left(\mathrm{S}_{\nu}, \mathrm{T}_{\nu}\right)$.

Consider then experiments $\mathbf{E}_{\nu}=\left\{\mathrm{P}_{\theta, \nu} ; \theta \in \Theta\right\}$, a pair ( $\mathbf{X}, \mathbf{Y}$ ) of compact spaces and pairs $\left(\mathrm{S}_{\nu}, \mathrm{T}_{\nu}\right)$ where $\mathrm{S}_{\nu}$ is a statistic with values in $\mathbf{X}$ and $\mathrm{T}_{\nu}$ is statistic with values in $\mathbf{Y}$. Both are defined on $\mathbf{E}_{\nu}$. Under $\mathrm{P}_{\theta, \nu}$ the pair $\left(\mathrm{S}_{\nu}, \mathrm{T}_{\nu}\right)$ has a joint distribution $D_{\theta, \nu}=\left(S_{\nu}, T_{\nu}\right) \mathrm{P}_{\theta, \nu}$ considered as a positive linear functional on the space $C(\mathbf{X} \times \mathbf{Y})$. Vague convergence of distributions will be pointwise convergence on $\mathrm{C}(\mathbf{X} \times \mathbf{Y})$.

For simplicity we shall consider only the case where $\Theta$ is the two point set $\{0,1\}$. Then one can let $\mathrm{m}_{\nu}=\mathrm{P}_{0, \nu}+\mathrm{P}_{1, \nu}$ and $\mu_{\nu}=\mathrm{D}_{0, \nu}+\mathrm{D}_{1, \nu}$. Each joint measure on $\mathbf{X} \times \mathbf{Y}$ has a marginal on $\mathbf{X}$. This will be noted by a "prime" so that for instance the marginal of $\mu_{\nu}$ on $\mathbf{X}$ is $\mu_{\nu}{ }^{\prime}$. The corresponding marginals on $\mathbf{Y}$ will have double primes, so that the $\mathbf{Y}$-marginal of $\mu_{\nu}$ is $\mu_{\nu}{ }^{\prime \prime}$.

With these notations one can introduce likelihood ratios such as

$$
\mathrm{f}_{\nu}=\frac{\mathrm{dD}_{1, \nu}}{\mathrm{~d} \mu_{\nu}}, \mathrm{f}_{\nu}^{\prime}=\frac{\mathrm{dD}_{1, \nu}^{\prime}}{\mathrm{d} \mu_{\nu}{ }^{\prime}}, \mathrm{f}_{\nu}^{\prime \prime}=\frac{\mathrm{dD}_{1, \nu}{ }^{\prime \prime}}{\mathrm{d} \mu_{\nu}{ }^{\prime \prime}}
$$

and $\mathrm{f}_{\nu}^{*}=\frac{\mathrm{dP}_{1, \nu}}{\mathrm{dm}_{\nu}}$.
Note that $\mathrm{f}_{\nu}$ is a function of the pair ( $\mathrm{x}, \mathrm{y}$ ) and that, for instance, $\mathrm{f}_{\nu}{ }^{\prime}$ is a function of x only, conditional expectation of $\mathrm{f}_{\nu}$ given x for the measure $\mu_{\nu}$.

If the $\mathrm{D}_{\mathrm{i}, \nu}$ have vague limits $\mathrm{D}_{\mathrm{i}, \infty}$, the pair ( $\mathrm{D}_{0, \infty}, \mathrm{D}_{1, \infty}$ ) defines corresponding densities $f_{\infty}, f_{\infty}{ }^{\prime}, f_{\infty}{ }^{\prime \prime}$. It will be convenient to say that $D_{i, \infty}$ is the distribution of a pair ( $\mathrm{X}, \mathrm{Y}$ ) on $\mathbf{X} \times \mathbf{Y}$.

With this notation one can state the following.
Theorem 2. Assume that $\mathrm{D}_{\mathrm{i}, \nu}$ converges vaguely to $\mathrm{D}_{\mathrm{i}, \infty}$. Assume in addition the following:
(A) For the experiments $\mathbf{E}_{\nu}$ the $\mathrm{S}_{\nu}$ are distinguished.
(B) The net of pairs $\left(\mathrm{D}_{0, \nu}{ }^{\prime \prime}, \mathrm{D}_{1, \nu}{ }^{\prime \prime}\right)$ is stable.

Then for ( $\mathrm{D}_{0, \infty}, \mathrm{D}_{1, \infty}$ ) the X coordinate is sufficient and the joint distributions of ( $\mathrm{f}_{\nu}, \mathrm{f}_{\nu}{ }^{\prime}, \mathrm{f}_{\nu}{ }^{\prime \prime}, \mathrm{S}_{\nu}, \mathrm{T}_{\nu}$ ) for $\mathrm{D}_{\mathrm{i}, \nu}$ converge to that of

$$
\left(\mathrm{f}_{\infty}, \mathrm{f}_{\infty}{ }^{\prime}, \mathrm{f}_{\infty}^{\prime \prime}, \mathrm{X}, \mathrm{Y}\right) \text { for } \mathrm{D}_{\mathrm{i}, \infty} .
$$

Proof. The proof follows closely the argument used for Theorem 1. Thus we shall just give a short sketch. A first remark is that one may as well assume that $S_{\nu}$ is sufficient for $\mathbf{E}_{\nu}$. Indeed let $\overline{\mathrm{f}}_{\nu}=\mathrm{f}_{\nu}{ }^{\prime} \mathrm{S}_{\nu}$ be the image of $\mathrm{f}_{\nu}{ }^{\prime}$ by $\mathrm{S}_{\nu}$. Let $Q_{1, \nu}$ be the measure whose density with respect to $\mathrm{m}_{\nu}$ is $\overline{\mathrm{f}}_{\nu}$. Since $\mathrm{S}_{\nu}$ is asymptotically sufficient the inequalities of [3], page 72 , show that $\int\left|\mathrm{f}_{\nu}^{*}-\overline{\mathrm{f}}_{\nu}\right| \mathrm{d} \mu_{\nu} \rightarrow 0$. Equivalently $\quad\left\|\mathrm{P}_{1, \nu}-\right\| \mathrm{Q}_{1, \nu}\left\|^{-1} \mathrm{Q}_{1, \nu}\right\| \rightarrow 0$. Then replace $\mathrm{P}_{0, \nu}$ by $\mathrm{Q}_{0, \nu}=\mathrm{m}_{\nu}-\left\|\mathrm{Q}_{1, \nu}\right\|^{-1} \mathrm{Q}_{1, \nu}$. For the pair ( $\mathrm{Q}_{0, \nu}, \mathrm{Q}_{1, \nu}$ ) the statistics $\mathrm{S}_{\nu}$ are sufficient and the properties A and B are preserved.

Now, as in the proof of Theorem 1, find functions $\phi^{\prime} \in C(\mathbf{X})$ and $\phi^{\prime \prime} \in C(\mathbf{Y})$ such that $\int\left|\mathrm{f}_{\infty}^{\prime}-\phi^{\prime}\right| \mathrm{d} \mu_{\infty}{ }^{\prime}<\epsilon^{2}$ and $\int\left|\mathrm{f}_{\infty}^{\prime \prime}-\phi^{\prime \prime}\right| \mathrm{d} \mu_{\infty}{ }^{\prime \prime}<\epsilon^{2}$. Then $\int\left|\mathrm{f}_{\nu}{ }^{\prime}-\phi^{\prime}\right| \mathrm{d} \mu_{\nu}{ }^{\prime}<\epsilon^{2}$ and $\int\left|\mathrm{f}_{\nu}{ }^{\prime \prime}-\phi^{\prime \prime}\right| \mathrm{d} \mu_{\nu}{ }^{\prime \prime}<\epsilon^{2}$ for all sufficiently large $\nu$ and one can repeat the argument that yielded property (3) in Theorem 1. This gives the result for the systems ( $\mathrm{f}_{\nu}{ }^{\prime}, \mathrm{f}_{\nu}{ }^{\prime \prime}, \mathrm{S}_{\nu}, \mathrm{T}_{\nu}$ ). The extra part $\mathrm{f}_{\nu}$ need no special attention since $\int\left|\mathrm{f}_{\nu}-\overline{\mathrm{f}}_{\nu}\right| \mathrm{d} \mu_{\nu} \rightarrow 0$ for a function $\overline{\mathrm{f}}_{\nu}$ defined on $\mathbf{X} \times \mathbf{Y}$ by writing $\overline{\mathrm{f}}_{\nu}(\mathrm{x}, \mathrm{y})=\mathrm{f}_{\nu}{ }^{\prime}(\mathrm{x})$.

The result as stated follows.
Remark 1. The theorem can be extended without difficulty, except for more complex notation, to finite sets $\Theta$ instead of two point sets.

Remark 2. Note well that the result is not only that the ( $\mathrm{f}_{\nu}, \mathrm{f}_{\nu}{ }^{\prime}, \mathrm{f}_{\nu}{ }^{\prime \prime}, \mathrm{S}_{\nu}, \mathrm{T}_{\nu}$ ) have limiting distributions but that the distribution are those given by the likelihood ratios in the limits $\left(D_{0, \infty}, D_{1, \infty}\right)$.
Remark 3. The result does not remain valid under condition (A) alone, even if $\mathrm{T}_{\nu}$ is also exactly sufficient. The stability of the pairs $\left(\mathrm{D}_{0, \nu}, \mathrm{D}_{1, \nu}\right)$ is not inherited by their marginals ( $\mathrm{D}_{0, \nu}{ }^{\prime \prime}, \mathrm{D}_{1, \nu}^{\prime \prime}$ ).

An example can be constructed as follows. Let $X$ be the two point set $\{0,1\}$ and let $\mathbf{Y}$ be the interval $[0,1]$. Let $\Omega$ be $[0,1]$ with its Lebesgue measure $\lambda$. Let $\lambda_{i}$ be the Lebesgue measure on the segment $i \times Y$ in $\mathbf{X} \times \mathbf{Y}$. Let $h_{\nu}(\omega)$ be the $\nu^{\text {th }}$ binary digit of $\omega$ in $\Omega$. Consider an experiment $\mathbf{E}_{\nu}$ where for $\theta=0$ the distribution of $\omega \in \Omega$ is $\lambda$ and where for $\theta=1$ the distribution of $\omega$ has a density $2 h_{\nu}$ with respect to $\lambda$. Map the experiment on $\mathbf{X} \times \mathbf{Y}$ by the statistic $\omega \rightarrow\left[h_{\nu}(\omega), \omega\right]$. Here our $\mathrm{D}_{\mathrm{i}, \nu}$ have limits on $\mathbf{X} \times \mathbf{Y}$ in the sense that $\int \gamma \mathrm{dD}_{\mathrm{i}, \nu} \rightarrow \int \gamma \mathrm{dD}_{\mathrm{i}, \infty}$ for every bounded measurable function $\gamma$ on $\mathbf{X} \times \mathbf{Y}$. Each coordinate map is sufficient. Yet both $\mathrm{D}_{0, \infty^{\prime \prime}}$ and $\mathrm{D}_{1, \infty} \prime \prime$ are identical to the Lebesgue measure on $[0,1]$. Condition $A$ is satisfied. Condition $B$ is not. The systems ( $\mathrm{f}_{\nu}, \mathrm{f}_{\nu}{ }^{\prime}, \mathrm{f}_{\nu}{ }^{\prime \prime}$ ) have a limiting distribution, but it is not that of ( $f_{\infty}, f_{\infty}{ }^{\prime}, f_{\infty}^{\prime \prime}$ ).

A natural example of a similar situation will be encountered in Section 3.
If one can establish that our condition (B) holds, our Theorem 2 can be used to bypass the requirement of convergence of conditional distributions used by Davies' in Theorem 4.1 of [1]. There the role of our $S_{\nu}$ is played by Davies' $\tilde{\Delta}_{n}$ and the role of $T_{\nu}$ is played by the pair $\left(\mu_{n}, \tau_{n}\right)$.

In the set up of Theorem 2 conditional distributions can be used to establish (B) as follows.

Consider the joint distributions $D_{i, \nu}$ as distributions of the pair of coordinates $\left(X_{\nu}, Y_{\nu}\right)$ in $\mathbf{X} \times \mathbf{Y}$. Let us say that the conditional distributions of $X_{\nu}$ given $Y_{\nu}$ are finitely approximable if for every $\gamma=\mathrm{C}(\mathbf{X})$ and every $\epsilon>0$ there is a finite set $\left\{\psi_{j} ; j \in J\right\}$ of elements of $\mathrm{C}(\mathrm{Y})$ such that, taking conditional expectations for $\mu_{\nu}$ one has

$$
\inf _{\mathrm{j} \in \mathrm{~J}} \mathrm{E}\left|\mathrm{E}\left[\gamma\left(\mathrm{X}_{\nu}\right) \mid \mathrm{Y}_{\nu}\right]-\psi_{\mathrm{j}}\left(\mathrm{Y}_{\nu}\right)\right|<\epsilon
$$

for all sufficiently large $\nu$.
Lemma 2. Consider the situation described for Theorem 2. Let (A) be satisfied and assume that the conditional distributions of $\mathrm{X}_{\nu}$ given $\mathrm{Y}_{\nu}$ for $\mu_{\nu}$ are finitely approximable. Then condition ( $B$ ) holds.
Proof. As in Theorem 2, one can assume that $X_{\nu}$ is exactly sufficient. Let then $\overline{\mathrm{f}}_{\nu}(\mathrm{x}, \mathrm{y})=\mathrm{f}_{\nu}{ }^{\prime}(\mathrm{x})$. For $\mu_{\nu}$ the conditional expectations satisfy

$$
\mathrm{E}\left(\overline{\mathrm{f}}_{\nu} \mid \mathrm{Y}_{\nu}\right)=\mathrm{E}\left(\mathrm{f}_{\nu}^{\prime} \mid \mathrm{Y}_{\nu}\right)=\mathrm{f}_{\nu}^{\prime \prime}\left(\mathrm{Y}_{\nu}\right)
$$

By assumption (A) and Theorem 1, for each $\epsilon>0$ there is a $\phi \in \mathrm{C}(\mathbf{X})$ such that $\int\left|\mathrm{f}_{\nu}{ }^{\prime}-\phi\right| \mathrm{d} \mu_{\nu}{ }^{\prime}<\epsilon$ for all large $\nu$ this will give

$$
\mathrm{E}\left|\mathrm{Ef}_{\nu}^{\prime \prime}\left(\mathrm{Y}_{\nu}\right)-\mathrm{E}\left(\phi \mid \mathrm{Y}_{\nu}\right)\right|<\epsilon
$$

By the assumption of the lemma, the $\mathrm{E}\left(\phi \mid \mathrm{Y}_{\nu}\right)$ can be approximated by elements of $\mathrm{C}(\mathrm{Y})$. Therefore the finite approximability also holds for $\mathrm{f}_{\nu}{ }^{\prime \prime}$ and the $\mathrm{f}_{\nu}{ }^{\prime \prime}$ satisfy the condition given in Definition 2. This gives the desired result.

According to this lemma, the conditions (A) (B) of Theorem 2 are implied by those used by Davies in his Theorem 4.1. Note, however, that under the other restrictions imposed there, the type of convergence of conditional distributions required by Davies is a necessary requirement for his conclusion. This is so because of the properties of Laplace transforms defined on open sets. It would not be necessary if one would work with a finite $\Theta$, as done here.

In conclusion, Theorem 2 offers an alternate approach to that of Davies in a fairly general framework. The stability property (B) is often easy to verify (as in the cases treated in [4], where it is entirely obvious). In the branching process example used by Davies in [1], the stability is obvious once one has established the required asymptotic sufficiency. Some additional remarks on this case follow.

## 3. Some remarks on branching processes

In this section we consider on branching process constructed as follows. There is a probability measure p on the integers $\{1,2,3, \ldots\}$. If the $\mathrm{n}^{\text {th }}$ generation contains $m_{n}$ individuals, each one of them, independently of the rest, produces a progeny whose size is taken at random from the distribution $p$, leading to a total population $M_{n+1}$ for the $(n+1)^{\text {st }}$ generation.

If $\xi$ is a variable with distribution p , it will be assumed that $\mathrm{E} \xi=\mu>1$ and that $\sigma^{2}=\mathrm{E}(\xi-\mu)^{2}$ is finite and non zero. There is a huge literature on inference for such processes. See for instance the references in [1] and [2]. Thus we shall content ourselves with a few simple remarks.

In [1] Davies assume that $p$ depends on a single parameter $\theta=\mu$ and studies the asymptotic behavior of the experiment $f_{n}$ in which one observes $M_{0}=1$, $M_{1}, \cdots, M_{n}$. This is done through his Theorem 4.1 by introducing a larger experiment $E_{n}$ in which one observes not only the total size of the population but also the entire genealogical tree for each individual in it. Note that here the initial population has size $M_{0}=1$. One can also start with a population of size $z$ as done in [2] and let both $n$ and $z$ go to infinity. Naturally, the asymptotics will be different for such a case.

Davies did include a short remark about estimating parameters other than the mean $\mu$. As will be explained below a technique different from that of his Theorem 4.1 becomes necessary.

Technically it should be noted that the results of Davies depend very strongly on the fact that he works in contiguous neighborhoods of a fixed distribution $p$. By contrast the results of [2] are uniform over a certain class $\mathbf{P}$ of measures $p$, one of the important features being that it is proved that, asymptotically, the experiment can be parametrized through the pair ( $\mu, \sigma^{2}$ ) independently of all other features of $p$.

As said above the results of [2] depend materially on the fact that the initial population size $z$ goes to infinity. Here are some simple remarks on what happens if $M_{0}$ is kept equal to unity.

To begin consider numbers $\mu, \sigma^{2}, \mathrm{~b}_{0}$ and $\epsilon_{0}$ with $0<\sigma^{2} \leq \mathrm{b}_{0}<\infty$ and $\epsilon_{0}>0$ with $\mu>1$. Let $\mathbf{P}_{0}$ be the class of all probability measures on the integers $\{1,2, \ldots\}$ that satisfy the following requirements

1) $\mathrm{E} \xi=\mu$
2) $\mathrm{E}(\xi-\mu)^{2}=\sigma^{2}$
3) $\mathrm{E}|\boldsymbol{\xi}-\mu|^{3} \leq \mathrm{b}_{0}$
4) $\sup _{\mathrm{k}} \mathrm{p}(\mathrm{k}) \mathrm{p}(\mathrm{k}+1) \geq \epsilon_{\mathrm{o}}$.

Lemma 3. Let $\mu, \sigma^{2}, \mathrm{~b}_{0}$ and $\epsilon_{0}$ be fixed. Consider an experiment $\mathrm{f}_{\nu}$ in which one observes the entire process $\mathrm{M}_{0}=1, \mathrm{M}_{1}, \mathrm{M}_{2}, \cdots, \mathrm{M}_{\mathrm{n}} \cdots$. Assume that the possibilities for the distributions of progeny are either $\mathrm{p}_{\nu}$ or $\mathrm{q}_{\nu}$, both in $\mathbf{P}_{\mathbf{0}}$. Let $\mathrm{P}_{\nu}$ and $Q_{\nu}$ be the corresponding distributions for the process.

Then if $\left\|\mathrm{p}_{\nu}-\mathrm{q}_{\nu}\right\| \rightarrow 0$ so does $\left\|\mathrm{P}_{\nu}-\mathrm{Q}_{\nu}\right\|$.
Proof. Consider the behavior of $M_{n+1}$ conditionally given the size $M_{n}$ of the $n^{\text {th }}$ generation. According to local limit theorems, for which see for instance Petrov
 that for any $\mathrm{p} \in \mathbf{P}_{\mathbf{0}}$ one has

$$
\begin{aligned}
\sum_{k} \mid \mathrm{P}\left[\mathrm{M}_{\mathrm{n}+1}\right. & \left.=\mathrm{k} \mid \mathrm{M}_{\mathrm{n}}\right] \left.-\frac{1}{\sigma \sqrt{2 \pi \mathrm{M}_{n}}} \exp \left\{-\frac{1}{2} \frac{\left(\mathrm{k}-\mu \mathrm{M}_{\mathrm{n}}\right)}{\sigma^{2} \mathrm{M}_{\mathrm{n}}}\right\} \right\rvert\, \\
& \leq \mathrm{b}\left[\frac{1}{\mathrm{M}_{\mathrm{n}}^{1 / 3}}+\mathrm{M}_{\mathrm{n}}^{2 / 3} c^{\mathrm{M}_{n}}\right] .
\end{aligned}
$$

For any $\mathrm{p} \in \mathbf{P}_{0}$ the series $\sum_{\mathrm{n}} \mathrm{M}_{\mathrm{n}}^{-1 / 3}$ and $\sum_{\mathrm{n}} \mathrm{M}_{\mathrm{n}}^{2 / 3} \mathrm{c}^{M_{n}}$ are almost surely convergent, with uniform bounds as $p$ varies in $\mathbf{P}$. Thus, for pairs $(p, q)$ in $\mathbf{P}_{\mathbf{0}}$ the tail of the process $\left\{\mathrm{M}_{\mathrm{n}}\right\}$ will be uninformative. For our $\mathrm{f}_{\nu}$, if $\left\|\mathrm{p}_{\nu}-\mathrm{q}_{\nu}\right\|$ tends to zero, the distance between the two possible distributions of $\left\{\mathrm{M}_{1}, \cdots, \mathrm{M}_{\mathrm{n}}\right\}$ also tends to zero for each fixed $n$. The result follows.
Remark 1. With $M_{0}=1$ if $\left\|p_{\nu}-q_{\nu}\right\|$ does not tend to zero the experiment $f_{\nu}$ can retain a great deal of information. For instance it could happen that $\mathrm{p}_{\nu}$ and $\mathrm{q}_{\nu}$ are disjoint so that the first generation allows one to find out which one of the two is operating.

However, Lemma 3 or a slight modification can still be used to show that the asymptotic parametrizability by $\mu$ and $\sigma^{2}$ only, described in Duby and Rouault [2] or Swensen [6], remains, valid "locally", see below.

If one maintains condition (1) (2) and (3) but remove condition (4) of the definition of $\mathbf{P}_{\mathbf{0}}$ the situation can be entirely different. There are pairs $(p, q)$ satisfying (1) (2) and (3) such that $\mathrm{p}(\mathrm{k})>0$ for all $\mathrm{k} \geq 1$ but $\mathrm{q}(2 \mathrm{k}+1) \equiv 0$ for all integers $k$. (For the measure $q$ every birth is a twin birth). What happens then in the process $M_{0}=1, M_{1}, \cdots, M_{n}$ may be complex, but what happens in one generation is easy to see.

Consider experiments $\mathbf{G}_{\mathrm{n}}$ where one observes $\mathrm{M}_{\mathrm{n}+1}$ conditionally given the size $M_{n}=m_{n}$ of the $n^{\text {th }}$ generation. Assume that the distribution of the progeny is one of two possibilities $\mathrm{f}_{\mathrm{n}}=\left(1-\alpha_{\mathrm{n}}\right) \mathrm{q}+\alpha_{\mathrm{n}} \mathrm{p}$ or q , where $\alpha_{\mathrm{n}}$ is some number $\alpha_{\mathrm{n}} \in[0,1]$.

If $m_{n} \alpha_{n} \rightarrow \infty$, then the (conditional) distribution of $M_{n+1}$ under $f_{n}$ differs little from what it would be under $p$, even if $\alpha_{n} \rightarrow 0$. The conditional experiment $\mathbf{G}_{\mathrm{n}}$ does not degenerate, nor does it become perfect. The same is true if $m_{n} \rightarrow \infty$ and $m_{n} \alpha_{n}$ tends to a finite non zero limit. However if $m_{n} \rightarrow \infty$ and $m_{n} \alpha_{n} \rightarrow 0$, the experiment $\mathbf{G}_{\mathrm{n}}$ degenerates.

This suggests, but does not quite prove, that experiments $f_{n}$ in which one observes $\mathrm{M}_{0}=1, \mathrm{M}_{1}, \cdots, \mathrm{M}_{\mathrm{n}}$ and where the distribution of the progeny is either $f_{n}$ or $q$ will become perfect if $\mu^{r_{n}} \alpha_{n} \rightarrow \infty$ for some sequence $r_{n}$ such that $n-r_{n} \rightarrow \infty$. This would be so even if $\alpha_{n} \rightarrow 0$ so that $\left\|f_{n}-q\right\| \rightarrow 0$.

A proof that the $f_{n}$ do tend to the perfect experiment is easy. Assume $\mathrm{M}_{\mathrm{n}}=\mathrm{m}_{\mathrm{n}}$ and let $\xi_{\mathrm{j}} ; \mathrm{j}=1,2 \ldots, \mathrm{~m}_{\mathrm{n}}$ be the sizes of the individual progenies from $\mathrm{M}_{\mathrm{u}}$. Let $\eta_{\mathrm{j}}=\xi_{\mathrm{j}}$ modulo 2. Then $\Sigma \eta_{\mathrm{j}}$ is a binomial variable with $\mathrm{m}_{\mathrm{n}}$ trials and probability of success $\alpha_{\mathrm{n}} \pi$ with $\pi=\mathrm{p}\left\{\xi\right.$ is odd). If $\mathrm{m}_{\mathrm{n}} \alpha_{\mathrm{n}} \pi$ is large the probability that $\sum_{j} \eta_{j}$ is odd will be close to $1 / 2$. Thus, if $\mu^{r_{n}} \alpha_{n} \rightarrow \infty$, with $n-r_{n} \rightarrow \infty$, the probability that some $M_{k}, k \leq n$ will be odd will tend to unity under $f_{n}$ but will be exactly zero under $q$.

Under such conditions, note that the conditional distribution of $\frac{M_{n+1}-\mu M_{n}}{\sigma \sqrt{M_{n}}}$ is asymptotically normal $\mathbb{N}(0,1)$ for all $f_{n}$. Thus the statistics $\left(M_{n+1}-\mu m_{n}\right) / \sigma \sqrt{m_{n}}$ cannot be stable for the experiments $G_{n}$. This is very much the same thing as the behavior of the $Y_{\nu}$ in the example of Remark 3, following Theorem 2.

In the class $\mathbf{P}_{0}$ the expectation $\mu$ and variance $\sigma^{2}$ of the progeny distribution were kept constant. One can use the local limit theorems to obtain results in which $\mu$ and $\sigma^{2}$ vary. To do this, consider fixed numbers $0<\epsilon_{0}<\mathrm{b}_{0}<\infty$ and a sequence $\left\{\epsilon_{n}\right\}, \epsilon_{n}>0 \epsilon_{n} \rightarrow 0$. Let $\mathbf{P}_{1}$ be the class of distributions $p=\mathbf{L}(\xi)$ such that

1) $1+\epsilon_{0} \leq \mathrm{E} \xi \leq \mathrm{b}_{0}$
2) $\epsilon_{0} \leq \mathrm{E}(\xi-\mu)^{2} \leq \mathrm{b}_{0}$
3) $\mathrm{E}|\xi-\mu|^{3} \leq \mathrm{b}_{0}$
4) $\sup _{k} p(k) p(k+1) \geq \epsilon_{0}$.

Let $\mathbf{P}_{\mathrm{n}}$ be any subset of $\mathbf{P}_{1}$ that has a total variation norm diameter at most $\epsilon_{\mathrm{n}}$.
The approximation to likelihood ratios described by Duby and Rouault in [2] is still valid on $\mathbf{P}_{\mathrm{n}}$ if one observes $\mathrm{M}_{0}=1, \mathrm{M}_{2}, \cdots, \mathrm{M}_{\mathrm{n}}$. That approximation involves only the parameters $\mu$ and $\sigma^{2}$ and the expressions

$$
\sum_{\mathrm{k}=0}^{\mathrm{n}-1} \frac{\left(\mathrm{M}_{\mathrm{k}+1}-\mu \mathrm{M}_{\mathrm{k}}\right)^{2}}{\sigma^{2} \mathrm{M}_{\mathrm{k}}}
$$

We shall not enter into the details, for which see [2], but end on a remark concerning what happens in a class such as $\mathbf{P}_{\mathrm{n}}$ defined here if one keeps $\boldsymbol{\mu}$ fixed but vary $\sigma$.

Let $p_{n}$ and $q_{n}$ be two elements of $P_{n}$. With the same expectation $\mu$ but respective variances $\sigma_{\mathrm{n}}^{2}$ and $\tau_{\mathrm{n}}^{2}$. Let $\mathrm{h}_{\mathrm{k}+1}^{2}$ be the conditional square Hellinger distance

$$
\mathrm{h}_{\mathrm{k}+1}^{2}=\frac{1}{2} \sum_{\mathrm{j}}\left|\left[\mathrm{P}\left(\mathrm{M}_{\mathrm{k}+1}=\mathrm{j} \mid \mathrm{M}_{\mathrm{k}}, \mathrm{P}_{\mathrm{n}}\right)\right]^{1 / 2}-\left[\mathrm{P}\left(\mathrm{M}_{\mathrm{k}+1}=\mathrm{j} \mid \mathrm{M}_{\mathrm{k}}, \mathrm{q}_{\mathrm{n}}\right)\right]^{1 / 2}\right|^{2}
$$

Up to terms of order $\mathrm{M}_{\mathrm{k}}^{-1 / 3}$ this is approximately $1-\left[1-\frac{\left(\sigma_{\mathrm{n}}-\tau_{\mathrm{n}}\right)^{2}}{\sigma_{\mathrm{n}}^{2}+\tau_{\mathrm{n}}^{2}}\right]^{1 / 2}$.
Let $P_{n}$ be the distribution of $M_{0}=1, M_{1}, \cdots, M_{n}$ under $p_{n}$. Let $Q_{n}$ be the corresponding distribution under $q_{n}$. It follows as in Lemma 3 that if $n\left(\sigma_{n}-\tau_{n}\right)^{2}$ remains bounded then the sequences $\left\{P_{n}\right\}\left\{Q_{n}\right\}$ will be contiguous.

The experiment $\left(\mathrm{P}_{\mathrm{n}}, \mathrm{Q}_{\mathrm{n}}\right)$ behaves asymptotically as if one observed a variable $\mathrm{T}_{\mathrm{n}}=\theta \chi_{\mathrm{n}}^{2}$ where $\chi_{\mathrm{n}}^{2}$ is a chi-square with n degrees of freedom and $\theta$ is allowed to take value $\sigma_{\mathrm{n}}^{2}$ or $\tau_{\mathrm{n}}^{2}$.

To apply Davies' Theorem 4.1 one would look at the bigger experiment in which the genealogies of the individuals in the $n^{\text {th }}$ generation are all observed. This would give other measures, say $\mathrm{P}_{\mathrm{n}}^{*}$ and $\mathrm{Q}_{\mathrm{n}}^{*}$. To insure contiguity for $\left\{\mathrm{P}_{\mathrm{n}}^{*}\right\}$, $\left\{\mathrm{Q}_{\mathrm{n}}{ }^{*}\right\}$ one needs to have that $\mu^{\mathrm{n}}\left(\sigma_{\mathrm{n}}-\tau_{\mathrm{n}}\right)^{2}$ remain bounded. Davies' Theorem 4.1 could then be applied, but the reduced experiment $\left(P_{n}, Q_{n}\right)$ would tend to the trivial uninformative experiment. Thus the note about estimating other parameters at the end of Section 5 of [1] should be given more weight than its short length seems to suggest.

We would also like to attract the reader's attention to the fact that the conditions (4) imposed on our classes $\mathbf{P}_{\mathbf{0}}$ and $\boldsymbol{P}_{\mathbf{1}}$ are not visible in the formulations of [1] or [6], but an analogous condition is implied.

The apparent lack of need for our condition (4) in the paper cited is due to the fact that the authors work with alternatives $q_{n}$ such that $\left\|q_{n}-p\right\| \rightarrow 0$ for a fixed $p$, independent of $n$. A similar remark applies to the third moment requirement. As we have seen above a condition such as (4) is needed to obtain uniform results on classes such as $\mathbf{P}_{\mathrm{n}}$. It is also needed for the uniform results of [2].

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