

## ASYMPTOTIC EXPANSIONS ASSOCIATED WITH POSTERIOR DISTRIBUTIONS<sup>1</sup>

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**1. Introduction and summary.** In this paper, an extension of the investigation of Johnson (1967b) is made by giving a larger class of posterior distributions which possess asymptotic expansions having a normal distribution as a leading term. Asymptotic expansions for the related normalizing transformation and percentiles are also presented.

Before asymptotic expansions were treated rigorously, LaPlace (1847) gave an expansion for certain posterior distributions. The method used in this paper is a variation of his technique. Bernstein (1934), page 406, and von Mises (1964), chapter VIII, Section C, also treat special cases of these expansions.

The conditions imposed are sufficient to make the maximum likelihood estimate strongly consistent and asymptotically normal. They also include higher order derivative assumptions on the log of the likelihood. As shown by Schwartz (1966), the posterior distribution may behave well even when the maximum likelihood estimate does not. However, we have not attempted to find the weakest assumptions under which the posterior distribution has an expansion. For general conditions under which the posterior distribution converges in variation to a normal distribution with probability one see LeCam (1953) and (1958) for the independent case and Kallianpur and Borwanker (1968) for Markov processes.

In Section 2, we show that with probability one, the centered and scaled posterior distribution possesses an asymptotic expansion in powers of  $n^{-\frac{1}{2}}$  having the standard normal as a leading term. The number of terms in the expansion obtained is two less than the number of continuous derivatives of the log likelihood. All terms beyond the first consist of a polynomial multiplied by the standard normal density. The coefficients of the polynomial depend on the prior density  $\rho$  and the likelihood. The moments of the posterior distribution are shown to possess an expansion in Section 3. The following two sections present the normalizing transformation and percentile expansions. These last three expansions also apply for the case considered by Johnson (1967b) as does the information on the form of the terms in the expansion of the posterior distribution. To simplify the already heavy notation, these results are first proved for independent identically distributed random variables. The extension of all these results to the case of certain stationary ergodic Markov processes is immediate; Section 6 presents the necessary modifications.

Throughout this paper,  $\Phi$  and  $\varphi$  will denote the standard normal cdf and pdf respectively. Also,  $n$  will be assumed to range over the positive integers; thus in

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some cases, the order of the error term in the expansion may be kept for smaller  $n$  if the bounding constant is modified.

**2. Expansion of the posterior distribution.** Let  $X_1, X_2, \dots$  be a sequence of real-valued random variables independently distributed as  $f(x, \theta_0)$  where  $f(x, \theta)$  is a family of densities with respect to a dominating measure  $\mu$ . It is assumed that  $\theta_0$  was chosen according to a law  $\rho(\cdot)$  which has a density with respect to Lebesgue measure.

We shall require certain assumptions to hold throughout Sections 2-5.

**ASSUMPTION 1.**  $\Theta$  is an open interval of the real line. Let  $\bar{\Theta}$  denote a closed interval containing  $\Theta$ .

**ASSUMPTION 2.** For any  $\theta \in \Theta$  and  $t \in \bar{\Theta}$ , provided  $t \neq \theta$ ,

$$\int_{-\infty}^{\infty} |f(x, \theta) - f(x, t)| d\mu > 0.$$

**ASSUMPTION 3.**  $f(x, \theta)$  is jointly measurable in  $(x, \theta)$ .

**ASSUMPTION 4.** For each  $x, f(x, \theta)$  admits partial derivatives of the first and second order with respect to  $\theta$  and these are continuous in  $\theta$  for  $\theta \in \bar{\Theta}$ .

**ASSUMPTION 5.** The measures  $\prod_{i=1}^n f(x_i, \theta)$  are mutually absolutely continuous for each  $n = 1, 2, \dots$ . Therefore, a null set will have probability zero for all  $\theta$ . Assumption 4 and Assumption 5 assure that  $\log f(x, \theta)$  is well defined and continuous in  $\theta$ .

**ASSUMPTION 6.** If  $\lim_{i \rightarrow \infty} |\theta_i| = \infty$ , then  $\lim_{i \rightarrow \infty} f(x, \theta_i) = 0$  for all  $x$  except for perhaps a null set not depending on the sequence  $\{\theta_i\}$ .

**ASSUMPTION 7.** For all  $\theta \in \bar{\Theta}$ ,

$$E_{\theta} |\log f(X, \theta)| < \infty$$

$$0 < I(\theta) \equiv -E_{\theta} \left[ \frac{\partial^2}{\partial \theta^2} \log f(X, \theta) \right].$$

**ASSUMPTION 8.** For each  $\theta_0 \in \Theta$ , there exist functions  $G_1(x)$  and  $G_2(x)$  satisfying

$$\left| \frac{\partial}{\partial \theta} \log f(x, \theta) \right| \leq G_1(x) \quad \left| \frac{\partial^2}{\partial \theta^2} \log f(x, \theta) \right| \leq G_2(x)$$

for  $\theta$  in a neighborhood of  $\theta_0$  and also  $E_{\theta_0}[G_1(X)] < \infty$  and  $E_{\theta_0}[G_2(X)] < \infty$ . The functions  $G_1$  and  $G_2$  may depend on  $\theta_0$ .

**ASSUMPTION 9.** Let

$$f(x, \theta, \rho) = \sup_{|\theta - \theta'| \leq \rho} f(x, \theta') \quad \rho > 0, \text{ and}$$

$$Q(x, r) = \sup_{|\theta| > r} f(x, \theta) \quad r > 0.$$

For every  $\theta \in \bar{\Theta}$  and  $\rho, r > 0, f(x, \theta, \rho)$  and  $Q(x, r)$  are measurable functions of  $x$ . Moreover, for sufficiently small  $\rho$  and sufficiently large  $r$ ,

$$\begin{aligned} E_{\theta_0}[\log f(x, \theta, \rho)]^+ &< \infty \\ E_{\theta_0}[\log Q(x, r)]^+ &< \infty \end{aligned} \quad \text{for each } \theta_0 \in \Theta.$$

These conditions are just one set of the many variations that imply that the maximum likelihood estimate  $\hat{\theta}$  is strongly consistent and that  $n^{1/2}(\hat{\theta} - \theta_0)$  is asymptotically normal with variance  $1/I(\theta_0)$ . The consistency conditions are essentially those of Wald while the other conditions are similar to those employed by Wolfowitz (1965) and Weiss-Wolfowitz (1966).

Denoting the maximum likelihood estimate of  $\theta$  by  $\hat{\theta}$ , we study the behavior of the posterior distribution of the centered and scaled variable  $\phi$  where

$$(2.1) \quad \phi = (\theta - \hat{\theta})b(\hat{\theta}) \quad \text{and}$$

$$(2.2) \quad b(\hat{\theta}) = \left[ -\frac{1}{n} \sum_{i=1}^n \frac{\partial^2}{\partial \theta^2} \log f(x_i, \theta) \Big|_{\theta=\hat{\theta}} \right]^{1/2}$$

when the prior distribution has pdf  $\rho(\theta)$  with respect to Lebesgue measure.

The notation  $b(\hat{\theta})$  suppresses the fact that  $b$  is the sum of  $n$  terms which depend on the observed sequence  $\mathbf{x} = (x_1, x_2, \dots, x_n, \dots)$ .

Denote the posterior cdf of  $n^{1/2} \phi$  by  $F_n$ . The observations  $x_1, x_2, \dots, x_n$  enter  $F_n$  both from the posterior distribution of  $\theta$  and the centering and scaling constants. The following theorem establishes the existence of an asymptotic expansion for  $F_n$ .

**THEOREM 2.1.** *Let  $\theta_0 \in \Theta$  be fixed. In addition to Assumptions 1 through 9, let  $K$  be an integer and assume that, for each  $x$ ,  $\log f(x, \theta)$  had  $K+3$  continuous partial derivatives with respect to  $\theta$ . Let there exist functions  $G_k(x)$  with  $E_{\theta_0}[G_k(X)] < \infty$  and  $|(\partial^k/\partial \theta^k) \log f(x, \theta)| \leq G_k(x)$  for  $\theta$  in a neighborhood of  $\theta_0$   $k = 3, 4, \dots, K+3$ .*

*If  $\rho(\theta_0) > 0$  and  $\rho(\cdot)$  has  $K+1$  continuous derivatives in a neighborhood of  $\theta_0$ , there exist functions  $\{\gamma_j(\xi, \mathbf{x})\}$  and constants  $D_1$  and  $N_x$  such that*

$$|F_n(\xi) - \Phi(\xi) - \sum_{j=1}^K \gamma_j(\xi, \mathbf{x}) n^{-j/2}| \leq D_1 n^{-\frac{1}{2}(K+1)} \quad \text{for all } n > N_x$$

*on an almost sure set  $(\pi f(x, \theta_0))$ . Here  $D_1$  depends on  $K$  and  $N_x$  depends on  $K$  and the observed sequence  $\mathbf{x}$ .*

If  $\rho(\cdot)$  is positive and sufficiently smooth and the derivatives of  $\log f$  are dominated, an expansion will exist for every  $\theta_0 \in \Theta$ .

Information is available on the form of the  $\gamma_j$ .

**PROPOSITION 2.1.** *Under the Assumptions of Theorem 2.1, each  $\gamma_j(\xi, \mathbf{x})$  is a polynomial in  $\xi$  having coefficients bounded in  $x$  multiplied by the standard normal density.*

**REMARK 1.** The proof below is also valid for the expansion appearing in Theorem 1.1 of Johnson (1967b). It extends Johnson's results since only the first

few terms were shown to be polynomials multiplied by  $\varphi(\xi)$ . Also  $\rho$  need only have  $K+1$  derivatives for an expansion of  $K+1$  terms.

REMARK 2. Specializing these results to the exponential family  $C(\phi) \exp(\phi R(x))$ , we see that Theorem 2.1 allows us to work with the usual parameter, for instance  $p$  in the binomial instead of  $\log [p/(1-p)]$ , as was necessary under the formulation in Johnson (1967b).

REMARK 3. Since the expansion is uniform in  $\xi$ , it may be used to find the rate at which the posterior distribution of  $\theta - \hat{\theta}$  approaches the degenerate distribution. In particular,  $P[-\varepsilon \leq \theta - \hat{\theta} \leq \varepsilon] = F_n(\varepsilon n^{-\frac{1}{2}} b^{-1}) - F_n(-\varepsilon n^{-\frac{1}{2}} b^{-1})$ .

A proof will be constructed using the lemmas below and will be patterned after the development in Johnson (1967a, 1967b).

2.1 Preliminaries. The posterior pdf of  $\phi$ , when viewed as a function of  $\phi$ , is proportional to

$$(2.3) \quad \rho(\hat{\theta} + \phi b^{-1}) \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1}) / f(x_i, \hat{\theta})]$$

where  $b = b(\hat{\theta})$  is given by (2.2). Since  $\theta_0$  must be an interior point, in the remainder of this section, it will be assumed that selected  $\theta$  intervals about  $\theta_0$  lie entirely within  $\Theta$ . Further, without loss of generality, the intervals will be taken small enough so that  $|(\partial^k / \partial \theta^k) \log f(x, \theta)| \leq G_k(x)$  holds for  $k = 1, 2, \dots, K+3$ . Before establishing certain properties of the function (2.3), we state a version of the uniform strong law which will be employed repeatedly. See LeCam (1953) or Rubin (1956) for the proof.

THEOREM 2.2. Let  $C$  be a compact set and let  $u(x, t)$  be a real-valued function measurable in  $x$  for each  $t \in C$  and continuous in  $t$  for each  $x$ . Let  $H(x)$  satisfy

$$|u(x, t)| \leq H(x) \quad \int H(x) dF(x) < \infty.$$

If the  $X_i$  are independent  $F$ ,

$$P[\lim_{n \rightarrow \infty} \sup_{t \in C} |\sum_{i=1}^n u(X_i, t) n^{-1} - \int u(x, t) dF(x)| = 0] = 1.$$

Let

$$(2.4) \quad a_{2n}(\theta) = \frac{1}{2} n^{-1} \sum_{i=1}^n \frac{\partial^2}{\partial \theta^2} \log f(x, \theta).$$

The next result is well known, but the proof is included for completeness.

LEMMA 2.1. There exist a null set  $S_1^c$ ,  $\delta_1 > 0$  and for each  $x \in S_1$  an  $N_{1x}$  such that

$$|I(\theta_0) + 2a_{2n}(\theta)| \leq I(\theta_0)/2 \quad \text{all } |\theta - \theta_0| \leq \delta_1 \quad \text{when } n > N_{1x}.$$

PROOF. Since  $\Theta$  is open, Assumption 8 assures the existence of a  $\delta_1 > 0$  such that  $|E_{\theta_0} [\partial^2 / \partial \theta^2] \log f(X, \theta) + I(\theta_0)| \leq I(\theta_0)/3$  for  $|\theta - \theta_0| \leq \delta_1$ . An application of Theorem 2.2 to  $a_{2n}(\theta)$  gives the desired result.

Noting that  $b(\hat{\theta}) = [-2a_{2n}(\hat{\theta})]^\frac{1}{2}$ , the last result together with the strong consistency of  $\hat{\theta}$  provides an  $N_{2x} (> N_{1x})$  for each  $\mathbf{x} \in S_2$  such that  $|\hat{\theta} - \theta_0| \leq \frac{1}{2} \delta_1$

$$(2.5) \quad 0 < [I(\theta_0)/2]^\frac{1}{2} \leq b(\hat{\theta}) \leq [3I(\theta_0)/2]^\frac{1}{2} < \infty \quad \text{for } n > N_{2x}$$

where  $S_2^c$  is a null set containing  $\{\hat{\theta} \rightarrow \theta_0\}$ .

2.2. *Basic lemmas.* We now turn our attention to the product term in (2.3).

LEMMA 2.2. *There exists a  $\delta_2 (0 < \delta_2 < 1)$  with*

$$n^{-1} \log \left\{ \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1})/f(x_i, \theta)] \right\} \leq -\phi^2/6 \quad \text{for } |\phi| \leq \delta_2$$

when  $n > N_{2x}$  each  $\mathbf{x} \in S_{1 \cap} S_2$ .

PROOF. For each  $\mathbf{x} \in S_{1 \cap} S_2$ , choose  $\delta_2$  so that  $|\hat{\theta} + \phi b^{-1} - \theta_0| \leq \delta_1$  when  $|\phi| \leq \delta_2$  and  $n > N_{2x}$ . Here  $\delta_1$  appears in the statement of Lemma 2.1. An application of the extended mean value theorem together with the bound (2.5) yields the desired result.

The next approximation shows that the posterior mass outside of a  $\phi$ -interval may be neglected. A similar argument is given in LeCam (1953), Theorem 5b.

LEMMA 2.3. *There exists an  $\varepsilon > 0$ , a null set  $S_3^c$  with  $S_3 \subset S_{1 \cap} S_2$  and for each  $\mathbf{x} \in S_3$ , an  $N_{3x}$  such that*

$$\log \left\{ \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1})/f(x_i, \hat{\theta})] \right\} \leq -n\varepsilon \quad \text{for } |\phi| \geq \delta_2$$

whenever  $n > N_{3x}$ . This  $\delta_2$  appears in the statement of Lemma 2.2.

PROOF. Let  $\mathbf{x} \in S_2$  be fixed. According to (2.5),  $|\phi b^{-1}| \geq \delta_2 [3I(\theta_0)/2]^{-\frac{1}{2}}$  for  $|\phi| \geq \delta_2$  whenever  $n > N_{2x}$ . Let  $2\varepsilon_1 = \delta_2 [3I(\theta_0)/2]^{-\frac{1}{2}}$  and make  $n$  larger if necessary so that  $|\hat{\theta} - \theta_0| \leq \varepsilon_1$ . Then

$$(2.6) \quad \prod_{i=1}^n \{f(x_i, \hat{\theta} + \phi b^{-1})/f(x_i, \hat{\theta})\} \leq \sup_{|\theta - \theta_0| \geq \varepsilon_1} \left[ \prod_{i=1}^n f(x_i, \theta) \right] / \left[ \prod_{i=1}^n f(x_i, \theta_0) \right].$$

According to the argument given by Wald (1949), page 599-600, the right-hand side of (2.6) is bounded by an expression of the form

$$(2.7) \quad \left\{ \left[ \sum_{j=1}^J \prod_{i=1}^n f(x_i, \theta_j, \rho_{\theta,j}) + \prod_{i=1}^n Q(x_i, r_0) \right] / \prod_{i=1}^n f(x_i, \theta_0) \right\}$$

where  $\max_j E_{\theta_0} \log [f(X, \theta_j, \rho_{\theta,j})/f(X, \theta_0)]$  and  $E_{\theta_0} \log [Q(X, r_0)/f(X, \theta_0)]$  are negative. By Assumption 9, an application of the strong law of large numbers to the log of each term in (2.7) establishes the existence of an  $\varepsilon > 0$  and a set  $S_3$  included in  $S_{1 \cap} S_2$  such that

$$n^{-1} \log \left[ \sup_{|\theta - \theta_0| \geq \varepsilon_1} \left\{ \prod_{i=1}^n f(x_i, \theta) \right\} / \left\{ \prod_{i=1}^n f(x_i, \theta_0) \right\} \right] \leq -\varepsilon \quad n \geq N_{3x}$$

for some  $N_{3x}$  depending on  $\mathbf{x} \in S_3$ .

Let us introduce some further notation useful in the derivation. Set

$$(2.8) \quad a_{kn}(\theta) = n^{-1} \sum_{i=1}^n \partial^k / \partial \theta^k \log f(x_i, \theta) / k! \quad \text{for } k = 3, 4, \dots$$

where the dependence on  $\mathbf{x}$  is suppressed. In this notation, the log of the likelihood ratio under sufficient regularity conditions may be expressed as

$$(2.9) \quad \log \left[ \prod_{i=1}^n [f(x_i, \theta) / f(x_i, \hat{\theta})] \right] = n \sum_{i=2}^k a_{in}(\hat{\theta})(\theta - \hat{\theta})^i + n a_{k+1, n}(\theta_*) (\theta - \hat{\theta})^{k+1}$$

where  $\theta_*$  is an intermediate value. Denote the first  $k+1$  terms of the Taylor expansion of  $\rho(\cdot)$  about  $\hat{\theta}$  by  $\rho_k(\hat{\theta})$ . That is

$$(2.10) \quad \rho_k(\hat{\theta}) = \rho(\hat{\theta}) + \rho^{(1)}(\hat{\theta})(\theta - \hat{\theta}) + \cdots + \rho^{(k)}(\hat{\theta})(\theta - \hat{\theta})^k / k!$$

where  $\rho^{(j)}$  denotes the  $j$ th derivative of  $\rho$ .

In terms of the transformed parameter  $\phi$ , we establish the following lemma.

LEMMA 2.4. *Let the integer  $K$  be given. For each  $x$ , let  $\log f(x, \theta)$  have  $K+3$  continuous partial derivatives in  $\theta$ . Let  $\theta_0 \in \Theta$  be fixed and assume that there exist functions  $G_k(x)$  satisfying  $|\partial^k / \partial \theta^k \log f(x, \theta)| \leq G_k(x)$ , for  $\theta$  in a neighborhood of  $\theta_0$ , and  $E_{\theta_0}[G_k(X)] < \infty$  each  $k = 2, 3, \dots, K+3$ . Also assume that  $\rho^{(K+1)}(\cdot)$  is continuous in a neighborhood of  $\theta_0$ . Then there exists a constant  $M_1$ , a null set  $S_K^c$  and for each  $\mathbf{x} \in S_K$ , an  $N_{4\mathbf{x}} (> N_{3\mathbf{x}})$  such that*

$$(2.11) \quad \int_{\delta_2}^{\delta_2} \left| \exp \left[ n \sum_{k=2}^{K+3} a_{kn}(\hat{\theta})(\phi/b)^k \right] \rho_K(\hat{\theta} + \phi b^{-1}) - \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1}) / f(x_i, \hat{\theta})] \rho(\hat{\theta} + \phi b^{-1}) \right| d\phi \leq M_1 n^{-\frac{1}{2}(K+2)}$$

for  $n > N_{4\mathbf{x}}$ .

Here  $a_{kn}$  is defined by (2.8) and  $\rho_K$  by (2.10) and  $\delta_2$  appears in the statement of Lemma 2.2.

PROOF. Together with the bound (2.5) for  $b(\hat{\theta})$  on the interval  $|\hat{\theta} - \theta_0| \leq \frac{1}{2} \delta_1$ , an application of Theorem 2.2 to each  $a_{kn}$ ,  $k = 2, \dots, K+3$ , establishes the existence of a null set  $S_K^c$ , a constant  $M_2$  and for each  $\mathbf{x} \in S_K$ , an  $N_{4\mathbf{x}}$  such that  $|a_{kn}(\hat{\theta})| \leq M_2$ ,  $k = 2, \dots, K+3$ , for  $n > N_{4\mathbf{x}}$  each  $\mathbf{x} \in S_K$ .

Add and subtract  $\rho_K(\hat{\theta} + \phi b^{-1}) \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1}) / f(x_i, \hat{\theta})]$  to the integrand in (2.11). It follows from Lemma 2.2 and the inequality  $|\log \{ \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1}) / f(x_i, \hat{\theta})] \} - n \sum_{k=2}^{K+3} a_{kn}(\hat{\theta})(\phi/b)^k| \leq n |\phi|^{K+3} M_2'$  for  $|\phi| \leq \delta_2$  some  $M_2'$  that the integrand is bounded by

$$e^{-n\phi^2/6} \{ |\rho^{(K+1)}(\theta_*)(\phi/b)^{K+1} / (K+1)!| + |\rho_K(\hat{\theta} + \phi b^{-1})| n |\phi|^{K+3} M_2'' \}$$

where  $\theta_*$  is an intermediate value and  $M_2''$  some constant. For each  $x \in S_K$ ,  $n$  is selected greater than  $N_{4\mathbf{x}}$ . Then the continuity assumptions imply that the integrand is bounded by  $M_1 \exp(-n\phi^2/6) [|\phi|^{K+1} + n|\phi|^{K+3}]$  for  $|\phi| \leq \delta_2$ . The result follows upon making the change of variable  $u = n^{\frac{1}{2}} \phi$ .

2.3. *Proof of asymptotic nature of the expansion.* Define a function  $\psi_{Kn}(z)$  for each  $K$  and  $n$  by

$$(2.12) \quad \psi_{Kn}(z) = \sum_{k=3}^{K+3} b^{-k} a_{kn}(\hat{\theta}) z^{k-3}$$

where the dependence on  $\mathbf{x}$  is suppressed. Recalling that  $b^2 = -2a_{2n}(\hat{\theta})$ , we write

$$(2.13) \quad \rho_K(\hat{\theta} + \phi b^{-1}) \exp \left[ n \sum_{k=2}^{K+3} a_{kn}(\hat{\theta})(\phi/b)^k \right] = e^{-n\phi^2/2} [\rho_K(\hat{\theta} + \phi b^{-1}) e^{n\phi^3} \psi_{Kn}(\phi)]$$

and note that for fixed  $n$  and sequence  $\mathbf{x}$ , the second factor on the right-hand side is a particular evaluation of the function

$$(2.14) \quad P_K(w, z, \mathbf{x}) = \rho_K(\hat{\theta} + zb^{-1}) e^{w\psi_{K_n}(z)} \quad |w| \leq 2, \quad |z| \leq 2\delta_2.$$

For each  $\mathbf{x} \in S_K$  and  $n > N_{4\mathbf{x}}$ , we have

$$(2.15) \quad P_K(w, z, \mathbf{x}) = \sum_{l,m} c_{lm}(\mathbf{x}) w^l z^m$$

where the series converges absolutely and uniformly in the region  $\{|w| \leq 3/2, |z| \leq 3\delta_2/2\}$  (see Markushevich 101–105). The coefficients  $c_{lm}(\mathbf{x})$  are given by

$$(2.16) \quad l!m! c_{lm}(\mathbf{x}) = \left. \frac{\partial^{l+m}}{\partial w^l \partial z^m} P_K(w, z, \mathbf{x}) \right|_{z=0, w=0} \quad l, m = 0, 1, \dots$$

and the usual approximation shows that

$$(2.17) \quad |c_{lm}(\mathbf{x})| \leq M_{lm} < \infty \quad l, m = 0, 1, \dots$$

where  $M_{lm}$  does not depend on  $n$  for  $n > N_{4\mathbf{x}}$ . These estimates enable us to find constants  $A_1$  and  $A_2$  such that

$$(2.18) \quad |P_K(w, z, \mathbf{x}) - \sum_{l+m \leq K} c_{lm}(\mathbf{x}) w^l z^m| \leq A_1 |w|^{K+1} + A_2 |z|^{K+1}$$

for  $|z| \leq \delta_2$  and  $|w| \leq 1$  when  $\mathbf{x} \in S_K$  and  $n > N_{4\mathbf{x}}$ .

Denote the truncated series  $\sum_{l+m \leq K} c_{lm}(\mathbf{x}) w^l z^m$  by  $P_K'(w, z, \mathbf{x})$ . We are now ready to prove the main theorem.

**PROOF OF THEOREM 2.1.** Let  $K$  be an arbitrary but fixed integer,  $\mathbf{x} \in S_K$  and  $n > N_{4\mathbf{x}}$  where the latter quantities appear in the statement of Lemma 2.4. A simple approximation together with the lemmas above show that

$$(2.19) \quad \left| \int_{-\infty}^{\infty} \rho(\hat{\theta} + \phi b^{-1}) \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1}) / f(x_i, \hat{\theta})] d\phi - \int_{-\infty}^{\infty} e^{-n\phi^2/2} P_K'(n\phi^3, \phi, \mathbf{x}) d\phi \right| \leq B_1 n^{-\frac{1}{2}(K+2)} \quad \text{and}$$

$$(2.20) \quad \left| \int_{-\infty}^{\xi n^{-1/2}} \rho(\hat{\theta} + \phi b^{-1}) \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1}) / f(x_i, \hat{\theta})] d\phi - \int_{-\infty}^{\xi n^{-1/2}} e^{-n\phi^2/2} P_K'(n\phi^3, \phi, \mathbf{x}) d\phi \right| \leq B_2 n^{-\frac{1}{2}(K+2)}$$

for some  $B_1$  and  $B_2$ . The last expression is uniform in  $\xi$ .

Integrating the approximation and collecting terms, we obtain an  $N_{5\mathbf{x}}$  such that the expansions (2.21) and (2.22) hold for all  $\xi$  and all  $n > N_{5\mathbf{x}}$ .

$$(2.21) \quad \left| \int_{-\infty}^{\infty} \rho(\hat{\theta} + \phi b^{-1}) \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1}) / f(x_i, \hat{\theta})] d\phi - \sum_{j=0}^K \beta_j(\mathbf{x}) n^{-\frac{1}{2}(j+1)} \right| \leq B_3 n^{-\frac{1}{2}(K+2)} \quad \text{and}$$

$$(2.22) \quad \left| \int_{-\infty}^{\xi n^{-1/2}} \rho(\hat{\theta} + \phi b^{-1}) \prod_{i=1}^n [f(x_i, \hat{\theta} + \phi b^{-1}) / f(x_i, \hat{\theta})] d\phi - \sum_{j=0}^K \alpha_j(\xi, \mathbf{x}) n^{-\frac{1}{2}(j+1)} \right| \leq B_4 n^{-\frac{1}{2}(K+2)} \quad \text{all } \xi.$$

A change of variables gives

$$(2.23) \quad \alpha_j(\xi, \mathbf{x}) = \sum_{s=0}^j c_{s,j-s} \int_{-\infty}^{\xi} y^{2s+j} e^{-y^2/2} dy \quad \text{each } j = 0, 1, \dots, K$$

and  $\beta_j(\mathbf{x})$  corresponds to  $\alpha_j(\infty, \mathbf{x})$  and thus is zero when  $j$  is odd. Bound (2.17) implies that each  $|\alpha_j(\xi, \mathbf{x})|$  is bounded for all  $\xi$  and consequently  $|\beta_j(\mathbf{x})|$  is also bounded. Since  $\beta_0(\mathbf{x}) = (2\pi)^{\frac{1}{2}}\rho(\hat{\theta})$ , which is bounded away from zero, the quotient series corresponding to  $F_n(\xi)$  has leading term  $\Phi(\xi)$ . The remaining coefficients  $\{\gamma_j(\xi, \mathbf{x})\}$  satisfy

$$(2.24) \quad \alpha_j(\xi, \mathbf{x}) = \beta_0(\mathbf{x})\gamma_j(\xi, \mathbf{x}) + \sum_{s=1}^{j-1} \gamma_{j-s}(\xi, \mathbf{x})\beta_s(\mathbf{x}) + \beta_j(\mathbf{x})\Phi(\xi) \quad \text{for } j = 1, 2, \dots, K.$$

An induction argument using the lower bound on  $\beta_0$  and upper bounds on the other  $\beta$ 's and  $\alpha$ 's gives  $|\gamma_j(\xi, \mathbf{x})| \leq M_K < \infty$  all  $j \leq K$  uniformly in  $\xi$  for sufficiently large  $n$ . Consideration of the quotient series establishes the existence of a  $D_1$  and an  $N_x$  for each  $\mathbf{x}$  belonging to an almost sure set where the inequality in the theorem holds. This concludes the proof since the construction is valid for every integer  $K$ .

It should be noticed that the terms of the expansion may be obtained by formal division.

**PROOF OF PROPOSITION 2.1.** In (2.25), it is shown that  $\gamma_1$  is of the form asserted. From (2.24), it is sufficient to show that  $\alpha_j - \beta_j\Phi$  does not involve  $\Phi$ . If  $j$  is odd, this result follows from (2.23). When  $j$  is even, it suffices to show

$$\int_{-\infty}^{\xi} y^{2r+j} e^{-y^2/2} dy - \Phi(\xi)2^{r+(j+1)/2}\Gamma(r+j/2+\frac{1}{2})$$

does not involve  $\Phi$ . Integration by parts establishes the desired result.

**2.4. Calculation of terms.** Again as in the previous papers (Johnson (1967a, 1967b)), we have the same relationships between the  $\gamma$ 's and  $c_{im}$ 's except it must be remembered that both are functions of the observed sequence  $\mathbf{x}$ . Specifically, the leading term is  $\Phi$  and

$$(2.25) \quad \begin{aligned} \gamma_1(\xi, \mathbf{x}) &= -\varphi(\xi)c_{00}^{-1}[c_{10}(\xi^2+2)+c_{01}] \\ \gamma_2(\xi, \mathbf{x}) &= -\varphi(\xi)c_{00}^{-1}[c_{20}\xi^5+(5c_{20}+c_{11})\xi^3+(15c_{20}+3c_{11}+c_{02})\xi] \end{aligned}$$

where  $\varphi(\xi)$  denotes the first derivative of  $\Phi$ . Another term is given in the first reference above. The  $c_{im}$  given by (2.16) may also be expressed directly in terms of  $\rho$  and the likelihood together with their derivatives. In particular

$$(2.26) \quad \begin{aligned} c_{00} &= \rho(\hat{\theta}); & c_{01} &= b^{-1}\rho^{(1)}(\hat{\theta}); & c_{02} &= b^{-2}\rho^{(2)}(\hat{\theta}) \\ c_{10} &= b^{-3}a_{3n}(\hat{\theta})\rho(\hat{\theta}); & c_{11} &= b^{-4}a_{4n}(\hat{\theta})\rho(\hat{\theta})+b^{-4}a_{3n}(\hat{\theta})\rho^{(1)}(\hat{\theta}) \\ c_{20} &= 2^{-1}b^{-6}a_{3n}^2(\hat{\theta})\rho(\hat{\theta}) \end{aligned}$$

where  $b$  is defined by (2.2) and  $a_{kn}$  by (2.9) and  $\hat{\theta}$  is the maximum likelihood estimate.

The prior density enters the expansion in the term of order  $n^{-\frac{1}{2}}$  as  $\rho^{(1)}(\hat{\theta})/\rho(\hat{\theta})$  and the term of order  $n^{-1}$  as  $\rho^{(2)}(\hat{\theta})/\rho(\hat{\theta})$  and if  $c_{11} \neq 0$  as  $\rho^{(1)}(\hat{\theta})/\rho(\hat{\theta})$ .

**3. Expansion of the moments.** We also study the moments of the posterior distribution of  $\phi$ . The moments of  $n^{\frac{1}{2}}\phi$  are obtained in the obvious manner. For



simplicity, only integer moments are considered. Denote the  $r$ th moment of the posterior distribution by  $E_p(\phi^r)$ .

**THEOREM 3.1.** *Under the assumptions of Theorem 2.1, if  $\int |\phi|^r \rho(\phi) d\phi < \infty$  where  $r < K$ , there exist functions  $\{\lambda_{r,j}(\mathbf{x})\}$ , a constant  $D_2$  and for all  $\mathbf{x}$  outside a null set, an  $N_x^2$  such that*

$$(3.1) \quad |E_p[\phi^r] - \sum_{j=r}^K \lambda_{r,j}(\mathbf{x}) n^{-j/2}| \leq D_2 n^{-\frac{1}{2}(K+1)} \quad n > N_x^2.$$

The odd terms are zero.

**PROOF.** The expansion is the quotient of two expansions of the form (2.21) where the numerator has  $\rho$  replaced by  $\phi^r \rho$ . The denominator has a non-zero leading term and the first few terms of the numerator series are zero. Also both series have odd terms equal to zero. Finding the reciprocal series for the denominator and multiplying, we obtain the desired result.

The terms may be obtained by formal division. In particular, the first non-zero term is

$$(3.2) \quad \lambda_{r,r}(\mathbf{x}) = 2^{r/2} \Gamma(\frac{1}{2}(r+1)) \Gamma^{-1}(\frac{1}{2})$$

if  $r$  is even and

$$(3.3) \quad \lambda_{r,r+1}(\mathbf{x}) = \Gamma^{-1}(\frac{1}{2}) 2^{\frac{1}{2}(r+1)} \{2(r+1) a_3(\hat{\theta}) \Gamma(\frac{1}{2}(r+4)) + \Gamma(\frac{1}{2}(r+2)) \rho^{(1)}(\hat{\theta}) / \rho(\hat{\theta})\}$$

when  $r$  is odd.

The moments of  $\theta$  about a fixed point may be obtained from those of  $\phi$  using relation (2.1). Two special cases which hold for sufficiently large  $n$  depending on  $\mathbf{x}$  belonging to a set of probability one,

$$(3.4) \quad |E_p(\theta) - \hat{\theta} - b^{-1}(6a_3(\hat{\theta}) + \rho'(\hat{\theta})/\rho(\hat{\theta}))n^{-1}| \leq D_2' n^{-2} \quad \text{and}$$

$$(3.5) \quad |E_p(\theta - \hat{\theta})^2 - b^{-2} n^{-1}| \leq D_2'' n^{-2},$$

may be compared with Gnedenko (1962), page 413, equations (5'), who further specializes to normal populations.

**4. Normalizing transformation.** Let  $\eta_n(\xi) = \Phi^{-1}(F_n(\xi))$  be the normalizing transformation. The notation does not reflect the dependence on  $\mathbf{x}$ . For fixed  $\mathbf{x}$ ,  $n$  and  $\xi$ ,  $\eta$  is defined by

$$(4.1) \quad \Phi(\eta) = F_n(\xi).$$

The following theorem shows that  $\eta_n$  has an asymptotic expansion.

**THEOREM 4.1.** *Under the assumptions of Theorem 2.1, there exist functions  $\{\omega_j(\xi) = \omega_j(\xi, \mathbf{x})\}$ , a constant  $D_3$  and for all  $\mathbf{x}$  outside a null set, an  $N_x^3$  such that (4.2) holds.*

$$(4.2) \quad |\eta_n(\xi) - \xi - \sum_{j=1}^K \omega_j(\xi) n^{-j/2}| \leq D_3 n^{-\frac{1}{2}(K+1)} \quad \text{for } n > N_x^3.$$

Here  $D_3$  is independent of  $\xi$  for  $\xi$  in a finite interval. The  $\omega_j(\xi)$  are polynomials in  $\xi$  having coefficients which are bounded for sufficiently large  $n$ .

PROOF. Let  $\mathbf{x} \in S_K$  where  $S_K$  is specified in Lemma 2.4. Denote by  $\Phi^{-1}(u)$  the branch of the analytic function which is the inverse of  $\Phi$  and is real valued for  $0 < u < 1$ . For any finite  $\xi$ -interval,  $F_n(\xi)$  is monotone and bounded away from 0 and 1 for all sufficiently large  $n$ . Substitution of the expansion for  $F_n$  into the series for  $\Phi^{-1}$  establishes the existence of an asymptotic expansion satisfying (4.2) with  $D_2$  independent of  $\xi$  over finite  $\xi$ -intervals. The  $\omega_j$  may be obtained by formal substitution. Alternatively, they may be obtained by inserting  $\eta = \sum_{j=0}^K \omega_j n^{-j/2}$  into

$$\Phi(\eta) \sim \Phi(\xi) + \sum_{j=1}^K \gamma_j(\xi, \mathbf{x}) n^{-j/2}$$

and identifying like powers of  $n^{-\frac{1}{2}}$ . Taking the limit as  $n \rightarrow \infty$ , we see that  $\omega_0 = \xi$ . The other coefficients satisfy

$$(4.3) \quad \Phi^{(1)}(\xi)\omega_j(\xi) = \gamma_j(\xi, \mathbf{x}) + L(\omega_0, \dots, \omega_{j-1}, \xi), \quad j \geq 1$$

where  $L$  is a polynomial in the  $\omega$ 's without constant terms. Each coefficient of the polynomial is a homogeneous of the first degree in the derivatives  $\Phi^{(l)}$   $1 \leq l \leq j$ . Canceling  $\Phi^{(1)}$ , each  $\omega_j$  is seen to be polynomial. The equations (4.3) show that the  $c_{lm}$  enter only in the numerator, except for  $c_{00}$ , so that the coefficients of the polynomial  $\omega_j$  are bounded for sufficiently large  $n$ . Although the argument is Wasow's, his equation (3.8) seems incorrect and the last portion of his proof should be modified as above.

The terms  $\omega_j(\xi)$  of the expansion are expressible directly in terms of the  $\gamma_j(\xi, \mathbf{x})$  and by (2.26), they may be written in terms of the  $c_{lm}(\mathbf{x})$ . For the second and third terms, we have

$$(4.4) \quad \omega_1 = \gamma_1/\varphi \quad \text{and} \quad \omega_2 = (\gamma_2/\varphi) + \xi(\gamma_1^2/2\varphi^2).$$

In the next section, this expansion is inverted so as to express  $\eta$  in terms of  $\xi$ .

**5. Expansion of the percentiles.** Here we consider the solution of the equation  $\Phi(\eta) = F_n(\xi_n(\eta))$ .

**THEOREM 5.1.** *Under the assumptions of Theorem 2.1, there exist functions  $\{\tau_j(\eta) = \tau_j(\eta, \mathbf{x})\}$ , a constant  $D_4$  independent of  $\eta$  for  $\eta$  in a finite interval, and for each  $\mathbf{x}$  outside a null set, an  $N_x^4$  such that*

$$(5.1) \quad \left| \xi_n(\eta) - \eta - \sum_{j=1}^K \tau_j(\eta) n^{-j/2} \right| \leq D_4 n^{-\frac{1}{2}(K+1)} \quad n > N_x^4$$

where the  $\tau_j(\eta)$  are polynomials whose coefficients depend on the  $c_{lm}(\mathbf{x})$  and are bounded for sufficiently large  $n$ .

PROOF. Let  $\mathbf{x} \in S_K$  where  $S_K$  is specified in Lemma 2.4.  $F_n$  has a positive derivative over any finite  $\xi$ -interval when  $n$  is sufficiently large. Therefore  $\xi_n(\eta)$  is unique and  $\xi_n(\eta) \rightarrow \eta$ .

Formally set,

$$(5.2) \quad \xi = \sum_{j=0}^K \tau_j(\eta) n^{-j/2} + \tau_{K+1}^* n^{-\frac{1}{2}(K+1)}.$$

Following Wasow, insert this expression into

$$(5.3) \quad \eta = \sum_{j=0}^K \omega_j(\xi)n^{-j/2} + \omega_{K+1}^* n^{-\frac{1}{2}(K+1)}$$

and equate the coefficients of  $n^{-j/2}, j = 1, \dots, K+1$ , to zero. The  $\tau_j(\xi)$  turn out to be polynomials with bounded coefficients since the same is true of the first term  $\eta$  and the  $\omega_j$ 's.

The condition that the  $c_{lm}$  are bounded implies that the equation for determining  $\tau_{K+1}^*$  has a solution which remains bounded as  $n \rightarrow \infty$ . With this value, the right-hand side of (5.2) is a solution of (5.3) and hence of (4.1). Since the right-hand side of (5.2) is real, it must equal  $\xi_n(\eta)$ . This proves the theorem.

The coefficients may be obtained formally. We have

$$(5.4) \quad \tau_1 = c_{00}^{-1}(c_{10}\eta^2 + 2c_{10} + c_{01}) \quad \text{and}$$

$$(5.5) \quad \tau_2 = (5(c_{20}c_{00}^{-1} + c_{11}c_{00}^{-1} - c_{01}c_{10}c_{00}^{-2})\eta^3 + (2c_{10}^2c_{00}^{-2} - c_{01}^2c_{00}^{-2}/2 + 15c_{20}c_{00}^{-1} + 3c_{11}c_{00}^{-1} + c_{02}c_{00}^{-1})\eta$$

where the  $c_{lm}$  depend on  $x$  and are given by (2.26).

As is usual for such expansions,  $\eta + \sum_{j=1}^K \tau_j(\eta)n^{-j/2}$  is an approximate  $\alpha$ th percentile for  $F_n$  when  $\eta$  is the  $\alpha$ th percentile of  $\Phi$ . This follows directly if  $F_n$  is expanded about  $\xi_n(\eta)$  since  $F_n^{(1)}$  is bounded in a fixed neighborhood of  $\eta$  for all sufficiently large  $n$ .

**THEOREM 5.2.** *Under the assumptions of Theorem 2.1, there exist a constant  $D_6$  independent of  $\eta$  for  $\eta$  in a finite interval and for each  $x$  outside a null set, an  $N_x^6$  such that*

$$|F_n(\eta + \sum_{j=1}^K \tau_j(\eta)n^{-j/2}) - \alpha| \leq D_6 n^{-\frac{1}{2}(K+1)} \quad n > N_x^6.$$

**6. Extension of Markov processes.** The results in the previous sections extend immediately to include strictly stationary Markov processes. We shall only state sufficient conditions for the existence of a  $K+1$  term expansion. Kallianpur and Borwanker (1968) prove convergence in variation and our conditions essentially include theirs.

For each  $\theta \in \Theta$ , let  $\{X_n, n \geq 0\}$  be a strictly stationary ergodic Markov process defined on  $(\mathcal{X}, \mathcal{A}) = \prod_{i=1}^{\infty} (R, \mathcal{B})$  into  $(R, \mathcal{B})$  where  $(R, \mathcal{B})$  is the Borel real line and the probability measure is  $P_\theta$ . Each  $P_\theta$  is induced by a stationary initial distribution  $p_\theta(\cdot)$  and the stationary transition probabilities  $p_\theta(\cdot \cdot)$  defined on  $\mathcal{B}$  and  $\mathcal{B} \times R$  respectively. Let  $P_{n,\theta}$  denote the restriction of  $P_\theta$  to the  $\sigma$ -field generated by  $\{X_0, \dots, X_n\}$ . It will be assumed that the probability measures  $\{P_{n,\theta}, \theta \in \Theta\}$  are mutually absolutely continuous. In order to use the transition densities as derivatives, we make the additional but unnecessary assumption that both  $p_\theta(\cdot)$  and  $p_\theta(y, \cdot)$  are absolutely continuous with respect to a measure  $\mu$  having densities  $f(z, \theta)$  and  $f(y, z, \theta)$  for  $\theta \in \Theta$ . The notation  $E_\theta[\cdot]$  denotes expectation with respect to  $P_{1,\theta}$ . The original assumptions together with the additional smoothness conditions become

ASSUMPTION 1'.  $\Theta$  is an open interval of the real line.

ASSUMPTION 2'. For any  $\theta \in \Theta$  and  $t \in \bar{\Theta}$ , provided  $t \neq \theta$

$$\int_{-\infty}^{\infty} |f(y, z, \theta) - f(y, z, t)| d\mu > 0 \quad \text{for all } y.$$

ASSUMPTION 3'.  $f(y, z, \theta)$  and  $f(z, \theta)$  are jointly measurable in their arguments.

ASSUMPTION 4'. For each  $(y, z)$ , both  $f(z, \theta)$  and  $f(y, z, \theta)$  admit partial derivatives of the first and second order with respect to  $\theta$  and these are continuous for all  $\theta \in \bar{\Theta}$ .

ASSUMPTION 5'.  $\{P_{n, \theta}, \theta \in \bar{\Theta}\}$  for each  $n \geq 0$  are mutually absolutely continuous.

ASSUMPTION 6'. If  $\lim_{i \rightarrow \infty} |\theta_i| = \infty$ , then  $\lim_{i \rightarrow \infty} f(y, z, \theta_i) = 0$  for all  $(y, z)$  except perhaps for a  $P_{1, \theta_0}$ -null set which does not depend on the sequence  $\{\theta_i\}$ .

ASSUMPTION 7'. For all  $\theta \in \bar{\Theta}$ ,  $E_{\theta} |\log f(X_0, X_1, \theta)| < \infty$  and  $0 < I(\theta) \equiv -E_{\theta} [\partial^2 / \partial \theta^2 \log f(X_0, X_1, \theta)]$ .

ASSUMPTION 8'. For each  $(y, z)$ , there exist  $K+3$  partial derivatives of  $\log f(y, z, \theta)$  and  $\log f(z, \theta)$ , continuous for all  $\theta \in \bar{\Theta}$ . For each  $\theta_0 \in \Theta$ , there exist functions  $G_j, j = 1, 2, \dots, K+3$  satisfying  $|\partial^j \theta / \partial \theta^j \log f(y, z, \theta)| \leq G_j(y, z)$ , for  $\theta$  in a neighborhood of  $\theta_0$ , and  $E_{\theta_0} [G_j] < \infty$  for  $j = 1, 2, \dots, K+3$ .

ASSUMPTION 9'. Let

$$f(y, z, \theta, \rho) = \sup_{|\theta - \theta'| \leq \rho} f(y, z, \theta') \quad \rho > 0, \text{ and}$$

$$Q(x, r) = \sup_{|\theta| > r} f(y, z, \theta) \quad r > 0.$$

For every  $\theta \in \bar{\Theta}$  and  $\rho, r > 0$ ,  $f(y, z, \theta, \rho)$  and  $Q(y, z, r)$  are measurable functions of  $(y, z)$ . Moreover for sufficiently small  $\rho$  and sufficiently large  $r$ ,

$$E_{\theta_0} [\log f(y, z, \theta, \rho)]^+ < \infty \quad \text{and}$$

$$E_{\theta_0} [\log Q(y, z, r)]^+ < \infty \quad \text{for } \theta_0 \in \Theta.$$

ASSUMPTION 10'. The prior density  $\rho(\cdot)$  has  $K+1$  continuous derivatives in a neighborhood of  $\theta_0$  and  $\rho(\theta_0) > 0$ . Here

$$(6.1) \quad \phi = (\theta - \hat{\theta})b(\hat{\theta}) \quad \text{where}$$

$$(6.2) \quad b(\hat{\theta}) = \left\{ -n^{-1} \left[ \frac{\partial^2}{\partial \theta^2} \log f(x_0, \theta) + \sum_{i=1}^n \frac{\partial^2}{\partial \theta^2} \log f(x_i, \theta) \right]_{\theta=\hat{\theta}} \right\}^{\frac{1}{2}}.$$

The  $a_{kn}$  are now defined in terms of  $n^{-1} \log [f(x_0, \theta) \prod_{i=1}^n f(x_{i-1}, x_i, \theta)] / [f(x_0, \hat{\theta}) \prod_{i=1}^n f(x_{i-1}, x_i, \hat{\theta})]$  instead of (2.9).

Let  $F_n(\cdot)$  denote the posterior cdf of  $\phi$ .

**THEOREM 6.1.** *The assumptions 1' through 10' imply the conclusion of Theorem 2.1 where the null set is taken w.r.t.  $f(x_0, \theta_0) \prod_{i=1}^{\infty} f(x_i, x_{i-1}, \theta_0)$ .*

Further, the results of Proposition 2.1, Theorem 3.1, Theorem 4.1, Theorem 5.1

and Theorem 5.2 hold under Assumptions 1' through 10'. The expressions for the terms in the expansions remain the same.

The proofs of Theorem 4.1, Theorem 5.1 and Theorem 5.2 do not require any changes. In the other cases, the initial distribution is handled separately with Assumption 8'. The proof of Theorem 2.1 is modified by employing the uniform strong law for stationary ergodic processes instead of Theorem 2.2. See Kallianpur and Borwanker (1968), Lemma 2.1, for a statement of this result. However the reasoning is the same as previously and there is no need to give new expressions for the terms of any expansion.

This result may be compared with Kallianpur and Borwanker (1968) who show that, almost surely  $P_{\theta_0}$ , the posterior distribution converges in variation to a normal distribution. This implies that the percentiles also converge. However they do not consider expansions of these quantities.

Another interesting comparison may be made with the work of Welch and Peers (1963). Dealing with the likelihood theory of Barnard, they consider weighted likelihoods which are mathematically identical to posterior distributions. Our results support their formal manipulations.

The author is presently investigating the multivariate extensions.

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