# A CENTRAL LIMIT THEOREM FOR FAMILIES OF STOCHASTIC PROCESSES INDEXED BY A SMALL AVERAGE STEP SIZE PARAMETER, AND SOME APPLICATIONS TO LEARNING MODELS

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## A CENTRAL LIMIT THEOREM FOR FAMILIES OF STOCHASTIC PROCESSES INDEXED BY A SMALL AVERAGE STEP SIZE PARAMETER, AND SOME APPLICATIONS TO LEARNING MODELS

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Let  $\theta > 0$  be a measure of the average step size of a stochastic process  $\{p_n^{(\ell)}\}_{n=1}^{\infty}$ . Conditions are given under which  $p_n^{(\ell)}$  is approximately normally distributed when n is large and  $\theta$  is small. This result is applied to a number of learning models where  $\theta$  is a learning rate parameter and  $p_n^{(\ell)}$  is the probability that the subject makes a certain response on the nth experimental trial. Both linear and stimulus sampling models are considered.

## 1. Introduction

Consider a learning experiment where response alternatives  $A_1$  and  $A_2$  are available to the subject on each of a sequence of trials, and, on any trial, response  $A_i$  is followed by outcome  $O_{ij}$ , presumed to be reinforcing to  $A_i$ , with probability  $\pi_{ij}$ . It is assumed that

$$\pi_{12} > 0$$
 and  $\pi_{21} > 0$ ,

which implies, in the learning models considered below, that a subject continues to choose both alternatives throughout the experiment. Let  $p_n$  be the probability that a subject will choose  $A_1$  on trial n, conditional upon past events (directly observable or hypothetical) that were crucial for learning. Mathematical learning theorists are often concerned with the asymptotic distribution of  $p_n$  as  $n \to \infty$ . Knowledge about the predicted distribution varies from model to model. There are interesting models for which not even the asymptotic mean of  $p_n$  is known, and, even in the few cases where formulas are available for some of the asymptotic moments of  $p_n$ , these usually give little insight into the shape of the asymptotic distribution.

Learning models typically involve one or more parameters whose magnitudes are directly related to the average step size of the corresponding process  $\{p_n\}$ . These may be loosely described as learning rate parameters. Quite often experimental data seem to call for small values of these parameters, so approximations to the asymptotic distribution of  $p_n$ , valid when learning rates are small, are of interest. An approximation of this sort, applicable to a wide variety of models, is obtained in this paper.

Let a family of models be indexed by a single parameter  $\theta$ , denote the stochastic process corresponding to  $\theta$  by  $\{p_n^{(\theta)}\}$ , and suppose that the cumulative distribution function of  $p_n^{(\theta)}$  converges as  $n\to\infty$  to a distribution  $F_\theta$ . The main result of this paper, the central limit theorem of Section 3 (henceforth abbreviated CLT), implies, roughly speaking, that if (a) certain means of  $p_{n+1}^{(\theta)} = p_n^{(\theta)}$  are of the order of magnitude of  $\theta$ , and (b) there is a point  $\rho$   $\epsilon$  (0, 1) toward which the process  $\{p_n^{(\theta)}\}$  drifts when  $\theta$  is small, then the normalized distribution

$$F_{\theta}(\theta^{1/2}x + \rho) = \lim_{n \to \infty} P_{\theta}((p_n - \rho)\theta^{-1/2} \leq x)$$

converges as  $\theta \to 0$  to a normal distribution with mean 0. (Throughout the paper  $\theta$  will appear as a subscript on P and E rather than as a superscript on  $p_n$  in probabilities and expectations involving the latter quantity.) A more precise statement of the theorem requires the functions

$$V(p, \theta) = E_{\theta}[p_{n+1} - p_n \mid p_n = p]$$

and

$$M(p, \theta) = E_{\theta}[(p_{n+1} - p_n)^2 | p_n = p],$$

which are assumed to be independent of n, and

$$a_n(\theta) = E_{\theta}[|p_{n+1} - p_n|^3].$$

In accordance with (a),  $V(p, \theta)$ ,  $M^{1/2}(p, \theta)$ , and  $a_n^{1/3}(\theta)$  are assumed to be  $O(\theta)$ . The condition (b) means that  $V(p, \theta)$  has the same sign as  $(\rho - p)$  when  $\theta$  is sufficiently small, or, in view of (a), that  $(\partial/\partial\theta)V(p, 0)$  is positive for  $p < \rho$  and negative for  $p > \rho$ . Thus  $\rho$  can be simply described as the unique root of  $(\partial/\partial\theta)V(\rho, 0) = 0$  in (0, 1). Clearly  $(\partial^2/\partial \rho \ \partial\theta)V(\rho, 0) \le 0$ , and, by (a),  $(\partial^2/\partial\theta^2)M(\rho, 0) \ge 0$ . The expression

$$\sigma^2 = \frac{1}{4} \frac{(\partial^2/\partial \theta^2) M(\rho, 0)}{-(\partial^2/\partial p \ \partial \theta) V(\rho, 0)}$$

given by CLT for the variance of the limiting normal distribution suggests that  $(\partial^2/\partial p \ \partial \theta)V(\rho, 0)$  must be assumed strictly negative, and, if the asymptotic distribution is to be non-degenerate, that  $(\partial^2/\partial \theta^2)M(\rho, 0)$  must be assumed strictly positive. A completely precise statement of the theorem will be given in Section 3 along with its proof. The proof is similar to the proofs of the normal convergence theorems of Norman [1966], all of which are special cases of CLT.

Some applications of the theorem to specific learning models will be given in Section 2 below. These results partially overlap those in Norman [1966].

## 2. Some Corollaries of CLT

# 2.1 The Four-Operator Linear Model

The process  $\{p_n\}$  for the general four-operator or experimenter-subject controlled events linear model [Bush & Mosteller, 1955] satisfies the equation

$$p_{n+1} = (1 - \theta_{ij})p_n + \theta_{ij}\delta_{ij}$$

when  $A_i$  and  $O_{ij}$  occur on trial n. It is assumed that  $1 > \theta_{ij} \ge 0$  for all i and j, and, in addition, that  $\theta_{ij} > 0$  for  $i \ne j$ , so that the process  $\{p_n\}$  has no absorbing barriers. The distribution of  $p_n$  converges as  $n \to \infty$  to a distribution that does not depend on  $p_1$  [Norman, 1968, Theorem 3.1], but, in this generality, even the mean of the asymptotic distribution is not known.

A one parameter subfamily of this four parameter family of models is obtained by constraining the vector  $\Theta = (\theta_{11}, \theta_{12}, \theta_{21}, \theta_{22})$  to a one dimensional subspace of four dimensional space. Thus for any vector  $C = (c_{11}, c_{12}, c_{21}, c_{22})$  of constants with  $c_{ii} \ge 0$  and  $c_{ij} > 0$  for  $i \ne j$ , the family of parameter-vectors  $\Theta$  of the form  $\Theta = \theta C$ ,  $0 < \theta < 1/\max_{ij} c_{ij}$  is considered. The corresponding functions  $V(p, \theta)$  and  $M(p, \theta)$  can be written in the form

(1) 
$$V(p, \theta) = \theta W(p)$$
 and  $M(p, \theta) = \theta^2 N(p)$ ,

where  $W(p)=(\partial/\partial\theta)V(p,0)$  and  $N(p)=(\partial^2/\partial\theta^2)M(p,0)/2$  are quadratic and cubic polynomials, respectively. Since  $|p_{n+1}^{(\theta)}-p_n^{(\theta)}|\leq \theta \max_{i,j} c_{ij}$ ,  $a_n(\theta)\leq \theta^3\max_{i,j}^3 c_{ij}$ . Now  $W(0)=c_{21}\pi_{21}>0$  and  $W(1)=-c_{12}\pi_{12}<0$ , and, since it is quadratic, it follows that W(p) has a unique root  $\rho$  in (0,1) and  $W'(\rho)<0$ . Obviously  $N(\rho)>0$ . These facts and CLT have the following implication.

Corollary 1. For the four-operator linear model with  $\Theta = \theta C$  and  $c_{ij} > 0$  for  $i \neq j$ ,

$$\lim_{\theta \to 0} \lim_{n \to \infty} P_{\theta}((p_n - \rho)\theta^{-1/2} \leq x) = \Phi(x/\sigma)$$

where  $\Phi$  is the normal distribution with mean 0 and variance 1,  $\rho$  is the root in (0, 1) of  $W(\rho) = 0$ ,  $\sigma^2 = N(\rho)/-2W'(\rho)$ , and N and W are defined in (1).

It can be shown that  $\rho$  is the expected operator approximation to  $\lim_{n\to\infty} E[p_n]$  [Bush and Mosteller, 1955]. Thus Corollary 1 provides a justification for expected operator approximation when learning rates are small.

2.2 Families of Models That Predict Approximate Probability Matching for Small Learning Rates

Throughout this subsection and the next it is implicitly assumed that  $O_{1j} = O_{2j} = O_j$ , j = 1, 2, and that the outcomes  $O_1$  and  $O_2$  are complimentary in the sense that  $O_1$  has the same reinforcing effect on  $A_1$  as  $O_2$  has on  $A_2$ .

Norman and Yellott [1966] defined abstract families of learning models indexed by a learning rate parameter  $\epsilon$  by means of their axioms F1 — F4. Their paper should be consulted for statements of these axioms as well as the definitions of the italicized terms, the function w, and the number  $\mu$  that appear below. If such a family predicts approximate probability matching for small learning rates or, equivalently [Norman and Yellott, 1966, Theorem 4] approximate marginal constancy for small learning rates, and if, in addition, axiom F3 is strengthened to require that the function w be thrice continuously differentiable on  $[0, 1] \times [0, \mu]$ , then it can be shown that the hypotheses of CLT are satisfied, and the following conclusion is obtained.

Corollary 2. If 
$$0 < \pi_{12}$$
,  $\pi_{21} < 1$ , then

$$\lim_{\epsilon \to 0} \lim_{n \to \infty} P_{\epsilon}((p_n - \ell)\epsilon^{-1/2} \leq x) = \Phi(x/\sigma),$$

where  $\ell = \pi_{21}/(\pi_{12} + \pi_{21})$  and

$$\sigma^2 = \frac{1}{2} \frac{\ell(1-\ell)}{\pi_{12} + \pi_{21}} \frac{\partial w}{\partial \epsilon} (\ell, 0).$$

# 2.3 The Component Model with Fixed Sample Size

The proportion  $p_n^{(1/N)}$  of stimulus elements conditioned to  $A_1$  on trial n in the N element component model with fixed sample size s [Estes, 1959] satisfies the stochastic difference equation

$$p_{n+1} - p_n = \begin{cases} \frac{j}{N} \\ 0 \\ -\frac{k}{N} \end{cases} \text{ with prob.}$$

$$\binom{N}{s}^{-1} \begin{cases} \binom{(1-p_n)N}{j} \binom{p_nN}{s-j} \binom{j}{s} \pi_{21} + \frac{s-j}{s} \pi_{11} \\ \binom{p_nN}{s} \pi_{11} + \binom{(1-p_n)N}{s} \pi_{22} \\ \binom{p_nN}{k} \binom{(1-p_n)N}{s-k} \binom{k}{s} \pi_{12} + \frac{s-k}{s} \pi_{22} \end{cases}$$

where  $N \ge s \ge j$ ,  $k \ge 1$  and  $\binom{j}{i}$  is to be interpreted as 0 if i > j.

A simple computation shows that  $V(p, \psi) = \psi(\pi_{21} - (\pi_{21} + \pi_{12})p)s$ , and a more complicated one yields

$$M(p, \psi) = \psi^2 s \left\{ (\pi_{21} - \pi_{11})s - 3(\pi_{21} - \pi_{11})(s - 1)p(1 - p)/(1 - \psi) + \pi_{11}(1 - p) \left[ 1 + (s - 1) \frac{(1 - p - \psi)}{1 - \psi} \right] + \pi_{22} p \left[ 1 + (s - 1) \frac{(p - \psi)}{1 - \psi} \right] \right\},$$

where s is fixed and  $\psi=1/N$  is regarded as the learning rate parameter. The process  $\{p_n^{(1/N)}\}$  is a finite Markov chain with only one ergodic class, and this class is aperiodic, so the distribution of  $p_n^{(1/N)}$  converges as  $n\to\infty$  to a distribution  $F_{1/N}$  (which is independent of the distribution of  $p_1^{(1/N)}$ ). Obviously  $V(p,\psi)$ ,  $M^{1/2}(p,\psi)$  and  $a_n^{1/8}(\psi)$  are  $O(\psi)$ ;  $(\partial/\partial\psi)V(p,0)=s(\pi_{21}-(\pi_{21}+\pi_{12})p)$  has the same sign as  $\pi_{21}/(\pi_{21}+\pi_{12})=p$ , so  $\rho=\pi_{21}/(\pi_{21}+\pi_{12})$ ; and  $(\partial^2/\partial p \ \partial \psi)V(\rho,0)=-(\pi_{21}+\pi_{12})s<0$ . Finally, it can be shown that

$$(\partial^2/\partial\psi^2)M(\rho,0) = 2s\rho(1-\rho)(2(\pi_{21}+\pi_{12})+s-1) > 0,$$

so CLT is applicable and yields this result.

Corollary 3. If  $0 < \pi_{12}$ ,  $\pi_{21} < 1$ , then

$$\lim_{N\to\infty}\lim_{n\to\infty}P_{1/N}((p_n-\rho)N^{1/2}\leq x)=\Phi(x/\sigma),$$

where  $\rho = \pi_{21}/(\pi_{21} + \pi_{12})$  and

$$\sigma^2 = \rho(1-\rho)(1+(s-1)/2(\pi_{21}+\pi_{12})).$$

## 3. A General Central Limit Theorem

Let S be a bounded set of positive real numbers having 0 as a limit point,  $\delta = \sup S$ , and  $J = [0, \delta]$ . For every  $\theta \in S$ , let  $\{p_n^{(\delta)}\}_{n=1}^{\infty}$  be a stochastic process (not necessarily Markov) with state space contained in a closed bounded interval [a, b] = I. Suppose that there exist twice continuously differentiable real valued functions V and M on  $I \times J$  such that, for every  $\theta \in S$  and  $n \ge 1$ ,  $E_{\theta}[p_{n+1} - p_n \mid p_n] = V(p_n, \theta)$  and  $E_{\theta}[(p_{n+1} - p_n)^2 \mid p_n] = M(p_n, \theta)$  with probability 1. Suppose also that  $(\partial^3/\partial\theta^3)M(\cdot, \cdot)$  exists and is continuous on  $I \times J$ , and that  $(\partial^3/\partial p^2 \partial \theta)V(\cdot, 0)$  and  $(\partial^3/\partial p \partial \theta^2)M(\cdot, 0)$  exist and are continuous on I. Finally, assume that, for every  $\theta \in S$ , the distribution function of  $p_n^{(\theta)}$  converges to a distribution function  $F_{\theta}$  as  $n \to \infty$ .

Central Limit Theorem. If

- (i)  $V(p, \theta) = O(\theta)$  and  $M(p, \theta) = O(\theta^2)$  for each  $p \in I$ ,
- (ii)  $a_n(\theta) = E_s[|p_{n+1} p_n|^3] = O(\theta^3)$  uniformly in n, and
- (iii)  $(\partial/\partial\theta)V(a, 0) > 0$ ,  $(\partial/\partial\theta)V(b, 0) < 0$ ,  $(\partial/\partial\theta)V(\cdot, 0)$  has a unique root  $\rho$  in (a, b), and

$$v = (\partial^2/\partial p \, \partial \theta) V(\rho, 0) < 0,$$

then  $m=(\partial^2/\partial\theta^2)M(\rho,0)\geq 0$ . Denote the normalized distribution  $F_{\theta}(\theta^{1/2}x+\rho)$  by  $G_{\theta}(x)$ . If m>0 then  $\lim_{\theta\to 0}G_{\theta}(x)=\Phi(x/\sigma)$  for all x, where  $\Phi$  is the standard normal distribution function and  $\sigma^2=m/-4v$ . If m=0 then  $\lim_{\theta\to 0}G_{\theta}(x)=\delta_0(x)$  for  $x\neq 0$ , where  $\delta_0(x)$  is the distribution function with all of its mass at 0.

Since, for every  $\theta \in S$  and  $n \ge 1$ ,  $|V(p_n^{(\theta)}, \theta)| \le M^{1/2}(p_n^{(\theta)}, \theta)$  with probability 1, one might hope that the first equality in (i) could be dispensed with.

No way of doing this has yet been found, and, at any rate, the development in Section 2 suggests that the verification of (i) and (ii) will be a triviality in most applications.

The following lemma is a basic component of the proof of CLT.

Lemma. Under hypotheses (i) and (iii) of CLT,  $\int_{-\infty}^{\infty} (p-\rho)^2 dF_{\delta}(p) = O(\theta)$  and  $m \ge 0$ . If m = 0 then

$$\int_{-\infty}^{\infty} (p-\rho)^2 dF_{\theta}(p) = O(\theta^{3/2}).$$

Proof of the Lemma. For  $\theta \in S$ ,

$$\begin{split} E_{\theta}[(p_{n+1} - \rho)^2] &= E_{\theta}[((p_n - \rho) + (p_{n+1} - p_n))^2] \\ &= E_{\theta}[(p_n - \rho)^2] + 2E_{\theta}[(p_n - \rho)(p_{n+1} - p_n)] + E_{\theta}[(p_{n+1} - p_n)^2] \\ &= E_{\theta}[(p_n - \rho)^2] + 2E_{\theta}[(p_n - \rho)V(p_{n+1} \theta)] + E_{\theta}[M(p_n \theta)]. \end{split}$$

Letting n approach  $\infty$  and using the Helly-Bray lemma [Loève, 1963, p. 180] we obtain

$$0 = 2 \int_a^b (p - \rho) V(p, \theta) dF_{\theta}(p) + \int_a^b M(p, \theta) dF_{\theta}(p).$$

It follows from (i) that  $V(p, 0) = M(p, 0) = (\partial/\partial\theta)M(p, 0) = 0$  for all  $p \in I$ . Hence, expanding  $V(p, \cdot)$  and  $M(p, \cdot)$  around 0, the former up to terms of second order and the latter up to terms of third order, and then expanding  $(\partial^2/\partial\theta^2)M(\cdot, 0)$  around  $\rho$  up to terms of first order we obtain

$$\begin{split} -2\int_a^b (p-\rho)\theta(\partial/\partial\theta)V(p,0)\;dF_\theta(p) \\ &=2\int_a^b (p-\rho)(\theta^2/2)(\partial^2/\partial\theta^2)V(p,\,\theta^*)\;dF_\theta(p)+(\theta^2/2)m \\ &+\int_a^b (\theta^2/2)(p-\rho)(\partial^3/\partial p\;\partial\theta^2)M(p^*,\,\theta)\;dF_\theta(p)+O(\theta^3) \end{split}$$

where  $0 < \theta^* < \theta$  and  $p^*$  is between p and  $\rho$ . But, by (iii),  $(\partial/\partial\theta)V(p,0)/(p-\rho)$  < 0 when  $p \in I$ ,  $p \neq \rho$ . Since  $(\partial^2/\partial p \,\partial\theta)V(\rho, 0) < 0$  also, it follows that there is some  $\gamma > 0$  such that  $(\partial/\partial\theta)V(p, 0)/(p-\rho) < -\gamma$  or  $-(\partial/\partial\theta) \cdot V(p, 0)/(p-\rho) > \gamma$  for all  $p \in I$ , where the quotient is defined by continuity at  $p = \rho$ . Since

$$\begin{split} 2\int_a^b \left(p-\rho\right)^2 \left[-(\partial/\partial\theta) V(p,0)/(p-\rho)\right] dF_\theta(p) \\ &= O\left(\theta\int_a^b \left|p-\rho\right| dF_\theta(p)\right) + (\theta/2)m + O(\theta^2), \end{split}$$

and since  $\int_a^b |p-\rho| dF_{\theta}(p) = O(1)$ , we obtain  $\int_a^b (p-\rho)^2 dF_{\theta}(p) = O(\theta)$ . Therefore  $2\gamma \int_a^b (p-\rho)^2 dF_{\theta}(p) = O(\theta^{3/2}) + (\theta/2)m$ . From this it follows that  $m \ge 0$  and that, if m = 0, then  $\int_a^b (p-\rho)^2 dF_{\theta}(p) = O(\theta^{3/2})$ . Q.E.D.

Proof of CLT. If m=0 the lemma implies  $\int_{-\infty}^{\infty} x^2 dG_{\theta}(x) = O(\theta^{1/2})$ . From this it follows that  $G_{\theta}$  converges to  $\delta_{0}$ , as claimed. We assume henceforth that m>0 so that the quantity  $\sigma^{2}$  defined in the statement of the theorem is positive.

For  $\theta \in S$ ,

$$\begin{split} E_{\theta}\{\exp\ [it(p_{n+1}\ -\ \rho)\theta^{-1/2}]\} \\ &= E_{\theta}\{\exp\ [it(p_n\ -\ \rho)\theta^{-1/2}]\ \exp\ [it(p_{n+1}\ -\ p_n)\theta^{-1/2}]\} \\ &= E_{\theta}\{\exp\ [it(p_n\ -\ \rho)\theta^{-1/2}]E_{\theta}\{\exp\ [it(p_{n+1}\ -\ p_n)\theta^{-1/2}]\ |\ p_n\}\}. \end{split}$$

Now

$$egin{aligned} E_{\theta} \{ \exp \left[ it(p_{n+1} - p_n) heta^{-1/2} \right] \mid p_n \} &= 1 + it heta^{-1/2} V(p_n \mid \theta) \\ &- (t^2/2) heta^{-1} M(p_n \mid \theta) + \omega(|t|^3/3!) heta^{-3/2} E_{\theta} [|p_{n+1} - p_n|^3 \mid p_n] \end{aligned}$$

where  $|\omega| \leq 1$  [Loève, 1963, p. 199]. Therefore substituting this into the previous equation, using (ii), letting  $n \to \infty$ , using the Helly-Bray lemma, and cancelling the term on the left of the resulting equation and the first term on the right we obtain

$$\begin{split} 0 &= \int_a^b \exp \left[it(p-\rho)\theta^{-1/2}\right] it\theta^{-1/2} V(p,\,\theta) \; dF_\theta(p) \\ &- \int_a^b \exp \left[it(p-\rho)\theta^{-1/2}\right] (t^2/2)\theta^{-1} M(p,\,\theta) \; dF_\theta(p) \, + \, O(|t|^3 \; \theta^{3/2}). \end{split}$$

The change of variables  $x = (p - \rho)\theta^{-1/2}$  yields

$$0 = \int_{-\infty}^{\infty} \exp(itx)it\theta^{-1/2}V(\theta^{1/2}x + \rho, \theta) dG_{\theta}(x)$$
$$-\int_{-\infty}^{\infty} \exp(itx)(t^{2}/2)\theta^{-1}M(\theta^{1/2}x + \rho, \theta) dG_{\theta}(x) + O(|t|^{3}\theta^{3/2}).$$

By (i) and the boundedness of the relevant third partial derivatives we have the Taylor expansions

$$V(p, \theta) = \theta(p - \rho)(\partial^2/\partial p \, \partial \theta)V(\rho, 0) + O(\theta(p - \rho)^2) + O(\theta^2)$$

and

$$M(p, \theta) = (\theta^{2}/2)(\partial^{2}/\partial\theta^{2})M(\rho, 0) + O(\theta^{2}|p-\rho|) + O(\theta^{3})$$

where  $O(\theta^2)$  and  $O(\theta^3)$  are uniform in p. Therefore

$$V(\theta^{1/2}x + \rho, \theta) = \theta^{3/2}xv + O(\theta^2x^2) + O(\theta^2)$$

and

$$M(\theta^{1/2}x + \rho, \theta) = (\theta^2/2)m + O(\theta^{5/2}|x|) + O(\theta^3).$$

But the lemma implies that  $\int_{-\infty}^{\infty} x^2 dG_{\theta}(x) = O(1)$  and  $\int_{-\infty}^{\infty} |x| dG_{\theta}(x) = O(1)$ . Thus we can write

$$0 = \int_{-\infty}^{\infty} \exp(itx)itxv \, dG_{\theta}(x)$$

$$-\int_{-\infty}^{\infty} \exp{(itx)(t^2/4)m} \ dG_{\theta}(x) + O((|t| + |t|^3)\theta^{1/2}).$$

Therefore

$$(d/dt)g_{\theta}(t) + \sigma^{2}tg_{\theta}(t) = O((1 + t^{2})\theta^{1/2})$$

where  $g_{\theta}$  is the characteristic function of  $G_{\theta}$ . The solution of this differential equation for which  $g_{\theta}(0)=1$  is

$$g_{\theta}(t) = \exp(-\sigma^2 t^2/2) \left(1 + \int_0^t \exp(\sigma^2 s^2/2) O((1 + s^2) \theta^{1/2}) ds\right)$$

Thus  $g_{\theta}(t) \to \exp(-\sigma^2 t^2/2)$  as  $\theta \to 0$ . It follows that  $G_{\theta}(x) \to \Phi(x/\sigma)$  for all x as  $\theta \to 0$ .

It is easy to verify that all of the normal convergence theorems discussed in Norman [1966] follow from CLT. From the present vantage point, however, it can be seen that the hypotheses of some of these theorems are unnecessarily strong, while others are completely superfluous. In the latter category are assumptions (45) and (49) in Theorem 4. In the former class are the important assumptions (10) of Theorem 1, (56) and (57) of Theorem 4, and (70) of Theorem 5, which imply that

(2) 
$$(\partial^2/\partial p \,\partial \theta)V(p,0) < 0$$
 for all  $p \in I$ .

This condition, in conjunction with  $(\partial/\partial\theta)V(a,0)>0$  and  $(\partial/\partial\theta)V(b,0)<0$ , implies hypothesis (iii) of CLT above. The inequality (2) need not be satisfied under the hypotheses of Corollaries 1 and 2. Consider, for instance, the case treated by Corollary 1 under the additional assumption that  $c_{1i}=c_{2i}=c_i$  (i.e.  $\theta_{1i}=\theta_{2i}=\theta_i$ ), j=1,2. The condition (2) can be shown to hold for all  $\pi_{12}>0$  and  $\pi_{21}>0$  if  $c_1=c_2$ , but to fail if  $c_1\neq c_2$  and  $\pi_{12}$  and  $\pi_{21}$  are sufficiently small.

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