IMA Preprint Series #243, Univ. Minnesota, 1986

THE RATE OF CONVERGENCE IN A CENTRAL LIMIT THEOREM FOR DEPENDENT RANDOM VARIABLES WITH ARBITRARY INDEX SET

(Abbreviated title: RATE OF CONVERGENCE IN A CLT)

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Moty 1986 SUMMARY

The classical rate of convergence in the central limit theorem for independent random variables with finite third moments has been found to be too crude for dependent random variables in general. This paper has two objectives. The first is to find a generalization of the classical rate of convergence. The second is to apply Stein's method and in doing so develop general techniques in the application.

We obtain results for a class of dependent random variables with arbitrary index set which includes m-dependent random fields as special cases and is relevant to spatial statistics. Our results can easily be extended to mixing random fields and in particular certain Gibbs random fields.

Research partially sponsored by the United States Army under Contract No. DAAG29-80-C-0041.

AMS 1980 Subject Classifications: Primary 60F05, 60G60.

Key Words and Phrases: Central limit theorem, rate of convergence, normal approximation, Stein's method, exchangeability argument, conditional concentration inequality, dependent random variables with arbitrary index set, mixing random field, m-dependent random field, spatial statistic.

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1. Introduction. Let  $X_1, \dots, X_n$  be random variables with  $EX_i = 0$ ,  $Var(\sum\limits_{i=1}^n X_i) = B_n^2$  and  $E[X_i]^3 = \gamma_i < \infty$ . Let  $F_n$  be the distribution function of  $B_n^{-1} \sum\limits_{i=1}^n X_i$  and  $\Phi$  the standard normal distribution function. The Berry-i=1 Esseen theorem (Berry (1941) and Esseen (1945)) states that if  $X_1, \dots, X_n$  are independent then there exists an absolute constant C such that

$$\|\mathbf{F}_{\mathbf{n}} - \Phi\|_{\infty} < CB_{\mathbf{n}}^{-3} \sum_{i=1}^{n} \gamma_{i}$$

Here and throughout the rest of this paper  $\|\cdot\|_p$ ,  $1 , denotes the usual <math>L_p$  norm on the real line and C denotes an absolute constant. If in addition  $x_1, \dots, x_n$  are identically distributed, then the bound on  $\|F_n - \Phi\|_{\infty}$  if of the order  $n^{-1/2}$ . It is known that the rate of convergence  $n^{-1/2}$  is best possible.

There has been considerable interest in the literature in obtaining a bound for  $\|\mathbf{F}_n - \phi\|_{\infty}$  in the case when  $X_1, \dots, X_n$  are dependent. Many attempts have been made to obtain a bound having an order as close to  $n^{-1/2}$  as possible and in many cases existence of moments higher than the third was assumed. Some of the early works in this direction were by Philipp (1969), Petrov (1970), Egorov (1970), Stein (1972), Erickson (1973, 1974) and Leonenko (1975). More recent works include those of Maejima (1978), Neaderhouser (1978b), Shergin (1979), Tikhomirov (1980), Riauba (1980), Prakasa Rao (1981),

Schneider (1981), Bolthausen (1982a), Takahata (1981, 1983) and Guyon and Richardson (1984). Of these the better results for sequences of random variables are by Sherqin, Tikhomirov, Schneider and Bolthausen, and the better results for random fields are by Takahata (1983) and Guyon and Richardson. Except for Bolthausen and Takahata these authors assumed existence of the third and sometimes lower moments. Of the latter we will state only those results assuming the existence of the third moment. Under this assumption, Tikhomirov obtained a bound of order  $n^{-1/2}(\log n)^2$  for strongly mixing strictly stationary sequences of random variables and a bound of order  $n^{-1/2}(\log n)^{3/2}$  for completely regular ( $\varphi$ -mixing) strictly stationary sequences. In each case the mixing rate is assumed to be exponentially decreasing. The latter result was extended to nonstationary sequences by Schneider who also improved the order of the bound to  $n^{-1/2}\log n$ . Tikhomirov's result for strictly stationary m-dependent sequences was generalized by Shergin who proved that if  $X_1, \dots, X_n$  is an m-dependent sequence with  $EX_i = 0$ ,  $Var(\sum_{i=1}^n X_i) = B_n^2$  and  $E|X_i|^3 = Y_i < \infty$ , then

(1.1) 
$$\|F_n - \Phi\|_{\infty} \le C(m+1)^2 B_n^{-3} \sum_{i=1}^n \gamma_i .$$

By a result of Berk (1973), the order for m in (1.1) is best possible. Bolthausen was able to obtain a bound of order  $n^{-1/2}$  for strongly mixing stationary Harris recurrent Markov chain by assuming certain mixing rate and the existence of a moment higher than the third.

To state the results of Takahata (1983) and Guyon and Richardson (1984), we regard the indices 1,...,n of  $X_1,\ldots,X_n$  as being from  $\mathbb{Z}^d$  for notational convenience. Takahata proved that for strongly mixing random fields with exponentially decreasing mixing rate and bounded absolute (8+ $\delta$ )th moments for some  $\delta > 0$  such that

(1.2) 
$$\lim_{n\to\infty} \inf n^{-1} \operatorname{Var}(\sum_{i=1}^{n} X_{i}) > 0,$$

we have

(1.3) 
$$|| F_n - \phi ||_{\infty} = O(n^{-1/2} (\log n)^d).$$

Guyon and Richardson obtained a rate of  $n^{-1/2}(\log n)^{2d}$  for strongly mixing random fields with exponentially decreasing mixing rate assuming (1.2) and only the boundedness of the absolute third moments. But the mixing condition assumed by Guyon and Richardson is different from that assumed by Takahata who took into account the sizes of two disjoint subsets of  $\mathbb{Z}^d$  in addition to the distance between them. The mixing condition of Guyon and Richardson cannot be applied to Gibbs random fields (see Dobrushin (1968)). For m-dependent random fields, Takahata and Guyon and Richardson proved under condition (1.2) that  $\|F_n - \phi\|_{\infty} = O(n^{-1/2})$ .

While Takahata assumed the boundedness of the eighth moments for this result, Guyon and Richardson assumed only the boundedness of the absolute third moments. The authors did not indicate how the constants in (1.3) and (1.4) depend on d, the moments and the mixing rate.

If an m-dependent sequence is stationary with  $B_n \to \infty$ , then (1.1) yields a rate of  $n^{-1/2}$ . Similarly if an m-dependent random field is stationary and  $Cov(X_i, Y_i) > 0$  where  $Y_i$  is the sum of those  $X_j$  dependent of  $X_i$ , then condition (1.2) is satisfied and the rate of  $n^{-1/2}$  in (1.4) holds. However, without stationarity, even if the random variables are bounded and identically distributed, the bound is (1.1) may tend to infinity and condition (1.2) may not be satisfied. Here is an example due to Erickson (1974). Define a sequence of bounded and identically distributed random variables  $X_1, \ldots, X_n$  with  $EX_i = 0$  and  $EX_i^2 = 1$  as follows. Let  $B_k^2 = Var(\sum_{i=1}^k x_i)$  and let  $X_1, X_2$  be independent. For  $k \ge 2$ , define  $X_{k+1} = -X_k$  if  $B_k^2 > k^{1/2}$  and define  $X_{k+1}$  to be independent of  $X_1, \ldots, X_n$  if  $B_k^2 \le k^{1/2}$ . It is not difficult to see that  $X_1, \ldots, X_n$  is a 1-dependent

sequence,  $|B_n^2 - n^{1/2}| \le 2$  and  $\sum_{i=1}^n X_i$  is a sum of  $B_n^2 = n^{1/2}$  terms of independent and identically distributed random variables. By the Berry-Esseen theorem

where  $\gamma = E |x_i|^3$ . But the bound in (1.1) is asymptotically  $C\gamma n^{1/4}$  which  $+\infty$ , and condition (1.2) for (1.4) is not satisfied. In view of the example of Erickson it is natural to ask whether a better bound exists and if it does, what is it? One of the objectives of this paper is to answer this question.

We consider random variables  $X_1, \dots, X_n$  with  $EX_i = 0$ ,  $Var(\sum_{i=1}^n X_i) = B^2$ and  $E | X_i |^3 = \gamma_i < \infty$  and satisfying the following dependence assumption. For each i  $\epsilon$  {1,...,n}, there exist nonempty  $A_i \subset B_i \subset C_i \subset \{1,...,n\}$  such that  $\{x_i\}$  is independent of  $x_{A_i^C}$ ,  $x_{A_i}$  is independent of  $x_{B_i^C}$  and  $x_{B_i}$ independent of  $X_{\underline{C}}$ . Here  $X_{\underline{A}}$  denotes  $\{X_{\underline{i}}: i \in A\}$ . We also assume that each  $C_i$  is relatively small compared to  $\{1,\ldots,n\}$ . One can see that the index set {1,...n} is actually arbitrary although the natural numbers are used. This class of dependent random variables has an interesting structure. If  $x_1, \dots, x_n$  is the n<sup>th</sup> row  $x_{n1}, \dots, x_{nn}$  of a triangular array with  $A_{ni}$  replacing  $A_i$  and so on, it has been proved in Chen (1978, theorem 4.1) that under very natural conditions,  $L(\sum_{i=1}^{n} X_{ni})$  converges weakly to an infinitely divisible distribution with finite canonical measure  $\nu$ and only if the signed measure  $\nu_{n}$  converges vaguely to  $\nu$  where  $Y_{ni} = \sum_{j \in A_{ni}} X_{nj}$  and  $v_n(A) = \sum_{i=1}^{n} Cov(X_{ni}, Y_{ni}I(Y_{ni} \in A))$  for  $A \in B(IR)$ . It is clear that this class of dependent random variables includes m-dependent random fields as special cases and many spatial statistics can be represented by a sum of such dependent random variables. (See for example Mantel (1967), David (1972), Sen (1976), Shapiro and Hubert (1979), and Cliff and Ord

(1981)). Furthermore the results proved in this paper can be extended in a fairly routine manner to mixing random fields and in particular certain Gibbs random fields. (See Dobrushin (1968), Holley and Stroock (1976) and Neaderhouser (1978a)). We have chosen not to include results for mixing random fields in this paper because we want to minimize the amount of tedious calculations in order to focus on the ideas and techniques of our proofs.

The main result in this paper is

(1.6) 
$$|F-\Phi| = C\lambda^{\frac{p-1}{p}} \text{ for } 1 \leq p \leq \infty ,$$

where F is the distribution function of  $B^{-1}$   $\sum_{i=1}^{n} X_i$ ,  $X_1, \ldots, X_n$  satisfy our dependence assumption, and  $\lambda$  and  $\beta$  are given by (2.6) and (2.7) respectively. We note that  $\lambda$  is a dimensionless quantity and  $\beta$  is a generalization of  $\sum_{i=1}^{n} \gamma_i$ . It is reasonable to assume that  $\lambda$  remains i=1 bounded as  $n + \infty$ . (See Remarks in Section 5.) It is shown in Section 5 that the bound in (1.6) coincides with that in (1.5) when it is applied to the example of Erickson. It is also shown that (1.6) yields

(1.7) 
$$\|F-\Phi\|_{p} < C\theta^{2}B^{-3} \sum_{i=1}^{n} \gamma_{i} \text{ for } 1 < p < \infty$$
,

where  $\theta$  is given by (2.5). When applied to m-dependent random fields on  $\mathbb{Z}^d$ ,  $\theta < (10m+1)^d$  (see (5.24)). For d=1 we obtain a bound of the same order as that in (1.1) with the best possible order for m.

Since  $I \cdot I_p^p < I \cdot I_{\infty}^{p-1} I \cdot I_1$ , it suffices to prove (1.6) for  $p = \infty$  and 1. An  $L_1$  bound for independent random variables was first proved by Esseen (1958) and that for m-dependent sequences by Erickson (1973, 1974) who obtained a bound of the order  $(m+1)^3 \sum_{i=1}^n \gamma_i$  in the second paper. The result of Erickson was generalized by Loh (1975) to completely regular sequences. Takahata (1983) considered the  $L_1$  problem under the same dependence assumptions as his  $L_{\infty}$  problem, and under the condition (1.2), obtained  $L_1$ 

bounds of the same orders as in (1.3) and (1.4) assuming the boundedness of the absolute (4+ $\delta$ )th moments ( $\delta$ >0) and the fourth moments respectively.

Most of the authors mentioned above used the characteristic function method and Bernstein's technique of dividing the sum into blocks and approximating the blocks by independent random variables. Stein (1972) invented his own method. Erickson (1974) and Loh (1975) adapted the method for L<sub>1</sub> bounds. Takahata's works (1981, 1983) are based on Stein's results. Tikhomirov (1980) and Guyon and Richardson (1984) combined techniques associated with Stein's method with the use of characteristic functions. Such an approach was also used by Chen (1978) and Bolthausen (1982b) but these authors were only concerned with limit theorems and not the rate of convergence.

Apart from obtaining new results, another objective of this paper is to apply Stein's method and in doing so develop general techniques which we hope are applicable to other problems. Stein's method may be briefly described as follows. Let  $W = B^{-1} \int\limits_{i=1}^{n} X_i$  where  $X_1, \dots, X_n$  are dependent random variables with  $EX_i = 0$  and  $Var(\sum_{i=1}^{n} X_i) = B^2$ . We derive an identity of the form

(1.8) 
$$E\{f'(W) - Wf(W)\} = ER_f(X_1, \dots, X_n)$$

for absolutely continuous f such that |f(w)| < C(1 + |w|) for some constant C, where the right hand side of (1.8) is an error term. By choosing f in (1.8) to be  $f_h$ , the unique bounded solution of the differential equation

(1.9) 
$$f'(w) - wf(w) = h(w) - Nh$$

where h is a bounded piecewise continuous function and Nh =  $(2\pi)^{-1/2} \int_{-\infty}^{\infty} h(x)e^{-x^2/2} dx$ , we obtain

(1.10) 
$$Eh(W) - Nh = ER_{f_h}(X_1, ..., X_n) .$$

If  $h = I_{(-\infty,z]}$ , then the left hand side of (1.10) is  $F(z) - \Phi(z)$ . We will be concerned only with  $h = I_{(-\infty,z]}$  and  $h = h_{z,\alpha}$  which is defined by (5.3). Bounds for  $\|F-\Phi\|_{\infty}$  and  $\|F-\Phi\|_{1}$  can then be obtained by appropriately bounding the right hand side of (1.10). There are two major steps in the application of Stein's method: the derivation of an identity of the form (1.8) and the bounding of the error term. The techniques involved in these two steps depend on the nature of dependence among  $X_1, \dots, X_n$  and may vary from problem to problem. Our aim is to develop these techniques in such a way that they may also be applicable to other problems.

For  $X_1,\dots,X_n$  satisfying our dependence assumption, we use an exchangeability argument to derive an identity of the form (1.8). Our exchangeability argument is related to Stein's idea of an exchangeability pair and an antisymmetric function but appears to be more intuitive and flexible in applications. (See Diaconis (1977) and Stein (1986).) By using (1.9) and the boundedness properties of  $f_h$ , a bound is immediately obtained for  $\|F-\Phi\|_1$  whereas the problem of bounding  $\|F-\Phi\|_{\infty}$  is reduced to that of bounding  $\|F^{\zeta}(a_{\zeta} \leq W \leq b_{\zeta})\|$  where  $\zeta$  is a random vector depending on a relatively small subset of  $\{X_1,\dots,X_n\}$ .

To bound  $p^{\zeta}(a_{\zeta} < W < b_{\zeta})$ , we use a conditional version of our identity and choose f in the identity to be dependent on  $\zeta$  and such that its derivative f' is an appropriate trapezoidal function. We call the resulting inequality for  $p^{\zeta}(a_{\zeta} < W < b_{\zeta})$  a conditional concentration inequality. The use of a concentration inequality in bounding  $\|F-\Phi\|_{\infty}$  was originally due to Stein in proving the Berry-Esseen theorem for independent and identically distributed random variables. A simplified version of Stein's proof is given in Ho and Chen (1978). However Stein's proof uses the symmetry inherited from the fact that the random variables are independent and identically distributed

and extension of the concentration inequality by using an identity of the form (1.8) even to independent but non-identically distributed random variables is difficult and remained unsolved for some time. In a discarded chapter of his monograph (1986), Stein has a proof of a concentration inequality involving an exhchangeable pair. The inequality is then applied to independent but non-identically distributed random variables. The proof of our conditional concentration inequality which also covers the independent but non-identically distributed case is different.

The foregoing discussion shows that bounding  $\|F-\Phi\|_{\infty}$  is more difficult than bounding  $\|F-\Phi\|_{1}$  using Stein's method. This fact was first observed by Erickson (1974).

We have only discussed the application of Stein's method to normal approximation. For an abstract formulation of Stein's method and other applications, see Stein (1986) or Chen (1979).

We have mentioned earlier that some of the above authors obtained their results assuming the existence of moments lower than the third. In particular, for m-dependent sequences Shergin (1979) obtained

(1.11) 
$$||\mathbf{F}_{n} - \Phi||_{\infty} \leq \inf_{\varepsilon > 0} C(m+1)^{2} [R_{2,n}(\varepsilon) + R_{3,n}(\varepsilon)]$$

and Erickson (1974) obtained

where  $R_{2,n}(\varepsilon) = B_n^{-2} \sum_{i=1}^n EX_i^2 I(|X_i| > \varepsilon B_n)$  and  $R_{3,n}(\varepsilon) = B_n^{-3} \sum_{i=1}^n E|X_i|^3 I(|X_i| < \varepsilon B_n)$ . It is possible to extend our techniques to obtain a bound for  $\|F-\Phi\|_p$  involving our  $\beta$  but in a way analogous to the bound in (1.11). However we sacrifice refinement for simplicity in this paper. Perhaps we will consider refinements and extension to mixing random fields in another paper.

2. A class of dependent random variables. In this section we state our dependence assumption and define notations which will be used throughout the paper. Let  $X_1, \ldots, X_n$  be random variables with  $EX_i = 0$ ,  $Var(\sum\limits_{i=1}^n X_i) = B^2$  and  $E[X_i]^3 = Y_i < \infty$ . We assume that for each  $i \in \{1, \ldots, n\}$ , there exist nonempty  $A_i \subset B_i \subset C_i \subset \{1, \ldots, n\}$  such that  $\{X_i\}$  is independent of  $X_{A_i}^{C}$ ,  $X_{A_i}^{C}$  is independent of  $X_{B_i}^{C}$  and  $X_{B_i}^{C}$  is independent of  $X_{A_i}^{C}$ . Here  $X_{A_i}^{C}$  denotes  $\{X_i: i \in A\}$ . It is clear that the index set of  $X_1, \ldots, X_n$  is arbitrary although we use the natural numbers for notational convenience.

For  $A \subset \{1, \ldots, n\}$ , we define

(2.1)  $N(A) \subset \{1,...,n\}$  to be such that for each  $i \in N(A)^C$ ,  $B_i \cap A = \emptyset$ . For each  $j \in \{1,...,n\}$ , we define

(2.2) 
$$N*(C^{j}) = \{i : j \in N(C_{i})\}$$
.

We assume that the  $C_i$  are relatively small compared to  $\{1,\ldots,n\}$  so that each  $N(C_i)$  is a proper subset of  $\{1,\ldots,n\}$ .

For each  $i \in \{1,...,n\}$ , let

$$\begin{cases} v_{i} = \sum_{j \in B_{i}^{C}} x_{j}, & y_{i} = \sum_{j \in A_{i}} x_{j}, \\ j \in B_{i} & j \in B_{i} \end{cases}, \quad T_{i} = \sum_{j \in C_{i}} x_{j}.$$

As usual A denotes the cardinality of the set A. We define

(2.4) 
$$\tilde{B}_{i}^{2} = Var(\sum_{j \in N(C_{i})^{C}} x_{j})$$

(2.5) 
$$\theta = \max_{1 \le i, j \le n} \{ |N(C_i)|, |N*(C^j)| \}$$

(2.6) 
$$\lambda = \max_{1 \le i \le n} \left( \frac{B}{\widetilde{B}_i} \right)^3 \vee 1$$

$$(2.7) \quad \beta = \sum_{i=1}^{n} \sum_{j \in N(C_{i}) \cup N^{*}(C^{i})} \{E|X_{j}X_{i}Y_{i}| + E|X_{j}X_{i}Z_{i}|\}$$

$$+ \theta \sum_{i=1}^{n} \{E|X_{i}^{2}Y_{i}| + EX_{i}^{2}E|Z_{i}|\}$$

$$+ \sum_{i=1}^{n} \{E|X_{i}Y_{i}Z_{i}| + E|X_{i}Y_{i}^{2}| + E|X_{i}Z_{i}^{2}|$$

$$+ E|X_{i}Y_{i}|E|Z_{i}| + E|X_{i}Y_{i}|E|T_{i}| + E|X_{i}Z_{i}|E|T_{i}|$$

$$+ E|X_{i}|E|Y_{i}^{2}| + E|X_{i}|E|Z_{i}^{2}$$

$$+ E|X_{i}|E|Y_{i}Z_{i}| + E|X_{i}|E|Z_{i}T_{i}| + E|X_{i}|E|T_{i}Y_{i}|\}.$$

Let  $\Lambda_0 \subset \Lambda \subset \{1,\ldots,n\}$  be such that  $X_{\Lambda_0}$  is independent of  $X_{\Lambda^c}$  and let  $\zeta$  be a random vector depending on  $X_{\Lambda_0}$ . Let  $X_1',\ldots,X_n'$  be an independent copy of  $X_1,\ldots,X_n$  and for i  $\epsilon$   $\{1,\ldots,n\}$ , let

$$(2.8) Y_{\underline{i}}' = \sum_{j \in A_{\underline{i}}} X_{\underline{j}}'.$$

We define

$$\widetilde{\mathbf{W}} = \sum_{\mathbf{i} \in \mathbf{N}(\Lambda)^{\mathbf{C}}} \mathbf{X}_{\mathbf{i}}$$

$$(2.10) \qquad \qquad \widetilde{B}^2 = Var(\widetilde{W})$$

(2.11) 
$$\Gamma(y,t) = I(y > t > 0) - I(y < t < 0)$$

(2.12) 
$$\hat{M}(t) = \sum_{i \in N(\Lambda)^{C}} X_{i} \Gamma(Y_{i} - Z_{i}, t)$$

(2.13) 
$$\hat{K}(t) = \sum_{i \in N(\Lambda)} c X_i' \Gamma(Y_i' - Z_i, t)$$

(2.14) 
$$K(t) = \hat{EK}(t)$$
.

Note that by the definitions of  $\Lambda_0$ ,  $\Lambda$  and  $N(\Lambda)$ , M(t) and  $\zeta$  are independent and by the dependence assumption  $X_i$  and  $Y_i - Z_i$  are independent for each i. Consequently

(2.15) 
$$E^{\zeta} \hat{M}(t) = EM(t) = 0 .$$

Likewise K(t) and  $\zeta$  are independent and so

(2.16) 
$$E^{\zeta} \hat{K}(t) = EK(t) = K(t)$$
.

By direct computations

(2.17) 
$$\int K(t) dt = \int_{i \in N(\Lambda)^{C}} EX_{i}'(Y_{i}' - Z_{i}) = \int_{i \in N(\Lambda)^{C}} EX_{i}'Y_{i}'$$
$$= \int_{i \in N(\Lambda)^{C}} EX_{i}Y_{i} = Var(\widetilde{W}) = \widetilde{B}^{2}.$$

If  $\Lambda_0 = \phi$ , then we take  $N(\Lambda) = \phi$ .

3. An exchangeability argument. Let  $X_1,\dots,X_n$  be dependent random variables with  $EX_i=0$  and  $Var(\sum\limits_{i=1}^n X_i)=1$  and let  $W=\sum\limits_{i=1}^n X_i$  (note i=1). In this section we describe a method for deriving identities of the form (1.8). For the time being we do not assume that  $X_1,\dots,X_n$  satisfy the dependence assumption of Section 2.

Suppose there exists a probability space  $(\Omega, B, P)$  on which the random variables  $X_i$ ,  $\xi_i$ ,  $X_i'$ ,  $\xi_i'$  and  $U_i$ ,  $i=1,\ldots,n$ , are defined such that for each i,

(i) 
$$L(x_i, \xi_i, x_i', \xi_i', U_i) = L(x_i', \xi_i', x_i, \xi_i, U_i)$$

Then we have

(3.1) 
$$\sum_{i=1}^{n} EX_{i}f(U_{i} + \xi_{i}) = \sum_{i=1}^{n} EX_{i}'f(U_{i} + \xi_{i}') ,$$

for absolutely continuous f such that  $|f(w)| \le C(1 + |w|)$  for some constant C. Suppose, in addition, for each i,

- (ii) U; is close to W in some sense,
- (iii)  $X_i$ ,  $\xi_i$ ,  $X_i$ ,  $\xi_i$  are small in some sense,
- (iv)  $(x_i',\xi_i')$  are weakly correlated to functions of  $(x_i,\xi_i,U_i,W)$ ,

(v) 
$$\sum_{i=1}^{n} Ex_{i}\xi_{i} \approx 1.$$

Then (3.1) can be rewritten as an identity of the form (1.8) with the error term small in some sense. We show heuristically how this can be done. We will derive explicit identities of the form (1.8) later when we apply this exchangeability argument to dependence random variables satisfying the dependence assumption of Section 2. By (ii) and (iii), the left hand side of (3.1)

$$\begin{array}{ccc}
 & n \\
 & \sum_{i=1}^{n} EX_{i}f(W) = EWf(W)
\end{array}$$

By (iv), the right hand side of (3.1)

$$= \sum_{i=1}^{n} EX_{i}'[f(U_{i} + \xi_{i}') - f(U_{i})]$$

$$= \int_{i=1}^{n} EX_{i}' \int \Gamma(\xi_{i}',t)f'(U_{i}+t)dt$$

where  $\Gamma$  is defined by (2.11). By (ii) and (iv), the last sum

$$= \int \left( \sum_{i=1}^{n} EX_{i}'\Gamma(\xi_{i}',t) \right) Ef'(W+t) dt$$

$$= \int \left( \sum_{i=1}^{n} Ex_{i} \Gamma(\xi_{i}, t) \right) Ef'(W+t) dt$$

$$= \int \tilde{K}(t) Ef'(W+t) dt$$

where  $\widetilde{K}(t) = \sum_{i=1}^{n} EX_{i}\Gamma(\xi_{i},t)$ . Now  $\widetilde{K}$  is the density of a signed measure and  $\sum_{i=1}^{n} \widetilde{K}(t) dt = \sum_{i=1}^{n} EX_{i}\xi_{i}$  which by (v) = 1. If the  $X_{i}$  and  $\xi_{i}$  are sufficiently small, then the mass of  $\widetilde{K}$  is concentrated near 0 and  $\int \widetilde{K}(t) Ef'(W+t) dt = Ef'(W)$ . Consequently (3.1) leads to

$$E\{f'(W) - Wf(W)\} \simeq 0$$

which is of the form (1.8) with small error term.

Sometimes it is more convenient to use

instead of (3.1) although in this paper we will use (3.1). It is easy to see that (i) also implies (3.2).

We now see how our exchangeability argument is related to Stein's use of an exchangeable pair and an antisymmetric function. Suppose W takes a finite number of values and there exists a probability space with W and W' defined on it such that L(W,W') = L(W',W). Stein (1986) used homology theory to prove that under certain connectedness condition if Eg(W) = 0 then there exists an antisymmetric function  $F: \mathbb{R}^2 + \mathbb{R}$  with  $E[F(W,W')] < \infty$  such that  $g(W) = E^W F(W,W')$ . Guided by this result he used the antisymmetric

function  $(w,w') \mapsto (w'-w)[f(w) + f(w')]$  and the easily verified fact  $E\{(W'-W)[f(W)+f(W')]\} = 0$  to derive an identity of the form (1.8). See Stein (1986, Lecture I). We show that (3.1) can be rewritten as

(3.3) 
$$EF_f(S,S') = 0$$

where L(S,S') = L(S',S), S and S' are approximately equal to W and  $F_f$  is an antisymmetric function similar to (w,w) + (w'-w)[f(w)+f(w')].

Let I be uniformly distributed on  $\{1,\ldots,n\}$  and independent of  $X_1$ ,  $\xi_1$ ,  $X_1$ ,  $\xi_1$ ,  $U_1$ ,  $U_1$ ,  $U_2$ ,  $U_3$ ,  $U_4$ ,  $U_5$ ,

$$(3.4) \qquad \qquad EU_{T}f(U_{T} + \xi_{T}) = EU_{T}f(U_{T} + \xi_{T}!)$$

where f is as in (3.1). Now (3.1) can be rewritten as

(3.5) 
$$EX_{T}f(U_{T} + \xi_{T}) = EX_{T}'f(U_{T} + \xi_{T}') .$$

Combining (3.4) and (3.5), we obtain

$$Esf(U_I + \xi_I) = Es'f(U_I + \xi_I')$$

which is

(3.6) 
$$\operatorname{ESE}^{S} f(U_{I} + \xi_{I}) = \operatorname{ES'E}^{S'} f(U_{I} + \xi_{I}') .$$

By (i) again, the regular conditional distribution of  $U_{\rm I}+\xi_{\rm I}$  given S=s is the same as that of  $U_{\rm I}+\xi_{\rm I}$ ' given S'=s. By letting  $\phi_{\rm f}(s)$  to be the regular conditional expectation  $E^{S=S}f(U_{\rm I}+\xi_{\rm I})$ , (3.6) becomes

$$ES\phi_{f}(S) = ES'\phi_{f}(S')$$

which is equivalent to

(3.7) 
$$E\{(S'-S)[\varphi_f(S) + \varphi_f(S')]\} = 0 .$$

The function  $(s,s') \leftrightarrow (s'-s)[\varphi_f(s) + \varphi_f(s')]$  is clearly antisymmetric.

Now we apply the exchangeability argument to  $X_1,\dots,X_n$  satisfying the dependence assumption of Section 2. We use the same definition and notations as in Section 2. We assume B=1 so that we have  $EX_1=0$ ,  $Var(\sum\limits_{i=1}^n X_i)=1$  and  $E|X_i|^3=\gamma_i$  <  $\infty$ . Let  $W=\sum\limits_{i=1}^n X_i$ . Since  $(X_1',\dots,X_n')$  is an

independent copy of  $(X_1, \dots, X_n)$  and for each  $i \in N(\Lambda)^C$ ,  $B_i \cap \Lambda = \emptyset$ , it follows that for each  $i \in N(\Lambda)^C$ ,

$$(X_i, Y_i)$$
,  $(X_i', Y_i')$  and  $(V_i, \zeta)$  are independent

and

$$L(x_{i}, Y_{i}, x_{i}', Y_{i}', V_{i}, \zeta) = L(x_{i}', Y_{i}', x_{i}, Y_{i}, V_{i}, \zeta)$$

Let  $(t,w) + f_t(w)$  be Borel measurable such that  $|f_{\zeta}(w)| < C(1+|w|)$  for some constant C. By the exchangeability argument

Here and throughout we drop the term "a.s." for simplicity. Now based on (3.8) we derive two slightly different identities of the form (1.8) which will be used separately for the next two sections. Write the left hand side of (3.8) as

$$(3.9) \qquad \sum_{\mathbf{i} \in \mathbf{N}(\Lambda)^{\mathbf{C}}} \mathbf{E}^{\zeta} \mathbf{X}_{\mathbf{i}} [f_{\zeta}(\mathbf{V}_{\mathbf{i}} + \mathbf{Y}_{\mathbf{i}}) - f_{\zeta}(\mathbf{W})] + \mathbf{E}^{\zeta} \widetilde{\mathbf{W}} f_{\zeta}(\mathbf{W})$$

$$= \sum_{\mathbf{i} \in \mathbf{N}(\Lambda)^{\mathbf{C}}} \mathbf{E}^{\zeta} \mathbf{X}_{\mathbf{i}} \int_{0}^{\mathbf{Y}_{\mathbf{i}} - \mathbf{Z}_{\mathbf{i}}} f_{\zeta}'(\mathbf{W} + \mathbf{t}) d\mathbf{t} + \mathbf{E}^{\zeta} \widetilde{\mathbf{W}} f_{\zeta}(\mathbf{W})$$

$$= \mathbf{E}^{\zeta} \int f_{\zeta}'(\mathbf{W} + \mathbf{t}) \widehat{\mathbf{M}}(\mathbf{t}) d\mathbf{t} + \mathbf{E}^{\zeta} \widetilde{\mathbf{W}} f_{\zeta}(\mathbf{W}) .$$

On the other hand, the right hand side of (3.8) equals

$$(3.10) \qquad \sum_{i \in N(\Lambda)} c \quad E^{\zeta} X_{i} \cdot [f_{\zeta}(V_{i} + Y_{i}') - f_{\zeta}(W)]$$

$$= \sum_{i \in N(\Lambda)} c \quad E^{\zeta} X_{i} \cdot \int_{0}^{Y_{i}' - Z_{i}} f_{\zeta}'(W+t) dt$$

$$= E^{\zeta} \int f_{\zeta}'(W+t) \hat{K}(t) dt$$

$$= E^{\zeta} \int f_{\zeta}'(W+t) [\hat{K}(t) - K(t)] dt$$

$$+ E^{\zeta} \int [f_{\zeta}'(W+t) - f_{\zeta}'(W)] K(t) dt + \tilde{B}^{2} E^{\zeta} f_{\zeta}'(W)$$

where we have used (2.17). Combining (3.8), (3.9) and (3.10), we obtain the first identity

$$(3.11) \qquad \widetilde{B}^{2} E^{\zeta} f_{\zeta}'(W) = E^{\zeta} \widetilde{W} f_{\zeta}(W) + E^{\zeta} \int f_{\zeta}'(W+t) \widehat{M}(t) dt$$

$$- E^{\zeta} \int f_{\zeta}'(W+t) [\widehat{K}(t) - K(t)] dt$$

$$- E^{\zeta} \int [f_{\zeta}'(W+t) - f_{\zeta}'(W)] K(t) dt .$$

For the second identity we go back to (3.8) and take  $\Lambda_0 = \phi$ . Then we have

By using the fact that  $X_i$  and  $(Z_i - Y_i, V_i)$  are independent and  $EX_i = 0$  we write the left hand side of (3.12) as

$$(3.13) \sum_{i=1}^{n} EX_{i}[f(V_{i} + Y_{i}) - f(W)] + EWf(W)$$

$$= -\sum_{i=1}^{n} E \int f'(V_{i} + Y_{i} + t)X_{i}\Gamma(Z_{i} - Y_{i}, t)dt + EWf(W)$$

$$= -\sum_{i=1}^{n} E \int [f'(V_{i} + Y_{i} + t) - f'(V_{i} + t)]X_{i}\Gamma(Z_{i} - Y_{i}, t)dt + EWf(W)$$

Since  $(X_i', Y_i')$  and  $(V_i, W)$  are independent and  $EX_i' = EX_i = 0$ , the right hand side of (3.12) equals

(3.14) 
$$\sum_{i=1}^{n} EX_{i}'[f(V_{i} + Y_{i}') - f(V_{i})]$$

$$= \sum_{i=1}^{n} E \int f'(V_{i} + t)X_{i}'\Gamma(Y_{i}',t)dt$$

$$= \sum_{i=1}^{n} E \int [f'(V_{i}+t) - f'(W)]EX_{i}'\Gamma(Y_{i}',t)dt$$

$$+ Ef'(W) \sum_{i=1}^{n} \int EX_{i}'\Gamma(Y_{i}',t)dt$$

$$= -\sum_{i=1}^{n} \mathbb{E} \int [f'(W) - f'(V_i + t)] \mathbb{E} X_i \Gamma(Y_i, t) dt + \mathbb{E} f'(W)$$

where we have used the fact that  $\int_{i=1}^{n} \int EX_{i}'\Gamma(Y_{i}',t)dt = \int_{i=1}^{n} EX_{i}'Y_{i}'$   $= \int_{i=1}^{n} EX_{i}Y_{i} = Var(W) = 1. Combining (3.12), (3.13) and (3.14), we obtain the$ 

second identity

(3.15) 
$$E\{f'(W) - Wf(W)\}$$

$$= - \sum_{i=1}^{n} E \int [f'(V_{i} + Y_{i} + t) - f'(V_{i} + t)] X_{i} \Gamma(Z_{i} - Y_{i}, t) dt$$

$$+ \sum_{i=1}^{n} E \int [f'(W) - f'(V_{i} + t)] EX_{i} \Gamma(Y_{i}, t) dt .$$

4. A conditional concentration inequality. As was mentioned in Section 1, a conditional concentration inequality is required to bound  $\|F - \Phi\|_{\infty}$ . Since such a concentration inequality is of interest in itself we prove it in a separate section here.

Let  $X_1, \ldots, X_n$  be random variables with  $EX_i = 0$ ,  $Var\left(\sum_{i=1}^n X_i\right) = B^2$  and  $E\left|X_i\right|^3 = \gamma_i < \infty$  and satisfying the dependence assumption of Section 2. Let  $W = B^{-1} \sum_{i=1}^n X_i$  and let  $\Lambda_0$ ,  $\Lambda$  and  $\zeta$  be as defined in Section 2.

PROPOSITION 4.1. For Borel measurable functions  $a_{\zeta}$  and  $b_{\zeta}$  of  $\zeta$  such that  $a_{\zeta} \leq b_{\zeta}$  , we have

$$(4.1) \qquad \qquad P^{\zeta} \left( a_{\zeta} \leq W \leq b_{\zeta} \right) \leq (3/2) \stackrel{\gamma-1}{B} \left( b_{\zeta} - a_{\zeta} \right) + 15 \stackrel{\gamma-3}{B} \beta \text{ a.s.}$$
 where  $\beta$  is defined by (2.7) and  $\stackrel{\gamma}{B}$  by (2.10).

If  $\Lambda_0 = \phi$ , we take  $P^{\zeta} = P$ ,  $A_r = A$ ,  $A_r = A$  and  $A_r = B$ .

To prove the proposition we assume without loss of generality that  $B=1 \text{ and } b_{\zeta}-a_{\zeta} \text{ bounded and use the identity (3.11). Denote the terms on the right hand side of (3.11) by <math display="block">R_{\dot{1}}^{\zeta}(f_{\zeta}), \ i=1,2,3,4, \ \text{respectively.}$  Now choose  $f_{\zeta}$  in (3.11) to be such that  $f_{\zeta}(2^{-1}(a_{\zeta}+b_{\zeta}))=0, \ f_{\zeta}$  is continuous and given by

$$f_{\zeta}'(w) = \begin{cases} 1 & , a_{\zeta} \leq w \leq b_{\zeta}, \\ 0 & , w \leq a_{\zeta} - \alpha \text{ or } \geq b_{\zeta} + \alpha, \\ \\ \text{linear}, a_{\zeta} - \alpha \leq w \leq a_{\zeta} \text{ or } b_{\zeta} \leq w \leq b_{\zeta} + \alpha \end{cases}$$

where  $\alpha$  is a positive number to be determined later. Clearly  $\left|f_{\zeta}\left(w\right)\right|\leq 2^{-1}\left(b_{\zeta}-a_{\zeta}\right)+\alpha$  . With this  $f_{\zeta}$  ,

(4.2) 
$$R_{1}^{\zeta}(f_{\zeta}) \leq \left[2^{-1}(b_{\zeta}-a_{\zeta}) + \alpha\right]E^{\zeta}|\widetilde{w}|$$

$$= \left[2^{-1}(b_{\zeta}-a_{\zeta}) + \alpha\right]E|\widetilde{w}|$$

$$\leq \widetilde{B}\left[2^{-1}(b_{\zeta}-a_{\zeta}) + \alpha\right].$$

By the geometric mean-arithmetic mean (GM-AM) inequality,

$$(4.3) R_2^{\zeta}(f_{\zeta}) \le 2^{-1}\widetilde{B} E^{\zeta} \int [f_{\zeta}'(W+t)]^2 dt + 2^{-1}\widetilde{B}^{-1} E^{\zeta} \int [\widetilde{M}(t)]^2 dt$$

$$\le 2^{-1}\widetilde{B} E^{\zeta}(b_{\zeta} - a_{\zeta} + 2\alpha) + 2^{-1}\widetilde{B}^{-1} E \int [\widetilde{M}(t)]^2 dt$$

$$= \widetilde{B}[2^{-1}(b_{\gamma} - a_{\gamma}) + \alpha] + 2^{-1}\widetilde{B}^{-1} \int Var[\widetilde{M}(t)] dt .$$

Likewise

(4.4) 
$$R_3^{\zeta}(f_{\zeta}) \leq \tilde{B}[2^{-1}(b_{\zeta}-a_{\zeta}) + \alpha] + 2^{-1}\tilde{B}^{-1} \int Var[\tilde{K}(t)]dt$$
.

Now we abuse the notation a little and regard  $E^{\zeta}$  and  $p^{\zeta}$  as regular conditional expectation and probability given a particular value of  $\zeta$  until we state otherwise. Define

$$L_{\zeta}(\alpha) = \sup_{x \in \mathbb{IR}} P^{\zeta}(x \leq w \leq x+\alpha)$$
.

Then by writing

$$\begin{split} R_{\mathbf{4}}^{\zeta}(\mathbf{f}_{\zeta}) &= -\mathbf{E}^{\zeta} \int_{0}^{\infty} \int_{0}^{t} \mathbf{f}_{\zeta}^{\prime\prime}(\mathbf{W} + \mathbf{s}) \, \mathrm{d}\mathbf{s} K(\mathbf{t}) \, \mathrm{d}\mathbf{t} + \mathbf{E}^{\zeta} \int_{-\infty}^{0} \int_{\mathbf{t}}^{0} \mathbf{f}_{\zeta}^{\prime\prime\prime}(\mathbf{W} + \mathbf{s}) \, \mathrm{d}\mathbf{s} K(\mathbf{t}) \, \mathrm{d}\mathbf{t} \\ &= -\alpha^{-1} \int_{0}^{\infty} \int_{0}^{t} \left[ \mathbf{p}^{\zeta}(\mathbf{a}_{\zeta} - \boldsymbol{\alpha} < \mathbf{W} < \mathbf{a}_{\zeta}) - \mathbf{p}^{\zeta}(\mathbf{b}_{\zeta} < \mathbf{W} < \mathbf{b}_{\zeta} + \boldsymbol{\alpha}) \right] \, \mathrm{d}\mathbf{s} K(\mathbf{t}) \, \mathrm{d}\mathbf{t} \\ &+ \alpha^{-1} \int_{-\infty}^{0} \int_{\mathbf{t}}^{0} \left[ \mathbf{p}^{\zeta}(\mathbf{a}_{\zeta} - \boldsymbol{\alpha} < \mathbf{W} < \mathbf{a}_{\zeta}) - \mathbf{p}^{\zeta}(\mathbf{b}_{\zeta} < \mathbf{W} < \mathbf{b}_{\zeta} + \boldsymbol{\alpha}) \, \mathrm{d}\mathbf{s} K(\mathbf{t}) \, \mathrm{d}\mathbf{t} \right] \end{split}$$

we have

$$(4.5) \qquad R_{\mathbf{4}}^{\zeta}(f_{\zeta}) < \alpha^{-1} \int_{0}^{\infty} \int_{0}^{t} L_{\zeta}(\alpha) ds \left| K(t) \right| dt + \alpha^{-1} \int_{-\infty}^{0} \int_{t}^{0} L_{\zeta}(\alpha) ds \left| K(t) \right| dt$$

$$= \alpha^{-1} L_{\zeta}(\alpha) \int_{-\infty}^{\infty} \left| tK(t) \right| dt .$$

Since  $p^{\zeta}(a_{\zeta} \leq W \leq b_{\zeta}) \leq E^{\zeta}f_{\zeta}(W)$ , it follows from (3.11), (4.2), (4.3), (4.4) and (4.5) that

$$(4.6) p^{\zeta}(a_{\zeta} \leq w \leq b_{\zeta}) \leq (3/2) \widetilde{B}^{-1}(b_{\zeta} - a_{\zeta} + 2\alpha)$$

$$+ 2^{-1} \widetilde{B}^{-3} \int Var[\widehat{M}(t)] dt + 2^{-1} \widetilde{B}^{-3} \int Var[\widehat{K}(t)] dt$$

$$+ \alpha^{-1} L_{\zeta}(\alpha) \widetilde{B}^{-2} \int |tK(t)| dt .$$

Let  $\rho = \tilde{B}^{-3} \int \text{Var}[\tilde{M}(t)]dt + \tilde{B}^{-3} \int \text{Var}[\tilde{K}(t)]dt$ . Substituting  $a_{\zeta} = x$  and  $b_{\zeta} = x + \alpha$  in (4.6), we obtain

(4.7) 
$$L_{\zeta}(\alpha) < 9\tilde{B}^{-1}\alpha/2 + \rho/2 + \alpha^{-1}L_{\zeta}(\alpha)\tilde{B}^{-2} \int |tK(t)|dt$$
.

Now let  $\alpha = 2B^{-2} \int |tK(t)| dt = \alpha*$  say. Then (4.7) yields

(4.8) 
$$L_{\zeta}(\alpha^*) \leq 9B^{-1}\alpha^*/2 + \rho/2 + L_{\zeta}(\alpha^*)/2$$

From this we get

(4.9) 
$$L_{\zeta}(\alpha^*) \leq 9B^{-1}\alpha^* + \rho$$

$$= 18B^{-3} \int |tK(t)|dt + \rho .$$

Finally combining (4.6) and (4.9) and letting  $\alpha = \alpha^*$  we get

$$(4.10) P^{\zeta}(a_{\zeta} \leq W \leq b_{\zeta}) \leq (3/2) \stackrel{\sim}{B}^{-1}(b_{\zeta} - a_{\zeta}) + 15 \stackrel{\sim}{B}^{-3} \int |tK(t)| dt$$

$$+ \stackrel{\sim}{B}^{-3} \int Var[M(t)] dt + \stackrel{\sim}{B}^{-3} \int Var[K(t)] dt .$$

Now we can go back and regard the left hand side of (4.10) as ordinary conditional probability and the right hand side as a random variable. What remains to do is to bound the last three terms on the right hand side of (4.10) which we denote by  $R_1$ ,  $R_2$  and  $R_3$  respectively. For  $R_1$  we have

(4.11) 
$$\int |tK(t)|dt < \int |t| \sum_{i \in N(\Lambda)^C} E|x_i'\Gamma(Y_i' - Z_i, t)|dt$$

$$\left\langle \sum_{i=1}^{n} E[x_{i}'] \right\rangle |t|[I(Y_{i}'-Z_{i}) + I(Y_{i}'-Z_{i}) + I(Y_{i}'-Z_{i})]dt$$

$$= 2^{-1} \int_{i=1}^{n} E|X_{i}'(Y_{i}'-Z_{i})^{2}|$$

$$< \int_{i=1}^{n} \{E|X_{i}'Y_{i}'^{2}| + E|X_{i}'Z_{i}^{2}|\}$$

$$= \int_{i=1}^{n} \{E|X_{i}Y_{i}^{2}| + E|X_{i}|EZ_{i}^{2}\}.$$

For R2, direct computations yield

$$(4.12) \int Var[M(t)]dt$$

$$= \sum_{i \in N(\Lambda)^{C}} \sum_{j \in N(C_{i})} E\{X_{i}X_{j}[(Y_{i}-Z_{i})\wedge(Y_{j}-Z_{j})I((Y_{i}-Z_{i})\wedge(Y_{j}-Z_{j}) > 0) + (Y_{i}-Z_{i})\vee(Y_{j}-Z_{j})I((Y_{i}-Z_{i})\vee(Y_{j}-Z_{j}) < 0)]\}$$

$$< \sum_{i=1} \sum_{j \in N(C_{i})} E\{|X_{i}X_{j}|[|Y_{i}-Z_{i}|\wedge|Y_{j}-Z_{j}|]\}$$

 $< 2^{-1} \sum_{i=1}^{n} \sum_{j \in N(C_i)} E\{|x_i x_j|[|y_i| + |z_i| + |y_j| + z_j|]\}$ 

For R3, we have

(4.13) 
$$\int Var[\hat{K}(t)]dt$$

$$= \int \sum_{i \in N(\Lambda)^{C}} \sum_{j \in N(C_{i})} Cov(X_{i}'\Gamma(Y_{i}'-Z_{i},t),X_{j}'\Gamma(Y_{j}'-Z_{j},t))dt$$

$$< \int \sum_{i=1}^{n} \sum_{j \in N(C_{i})} \{ Var[X_{i}'\Gamma(Y_{i}'-Z_{i},t)] Var[X_{j}'\Gamma(Y_{j}'-Z_{j},t)] \}^{1/2}dt$$

$$< 2^{-1} \int \sum_{i=1}^{n} \sum_{j \in N(C_{i})} \{ Var[X_{i}'\Gamma(Y_{i}'-Z_{i},t)] + Var[X_{j}'\Gamma(Y_{j}'-Z_{j},t)] \} dt$$

$$< 2^{-1} \sum_{i=1}^{n} \sum_{j \in N(C_{i})} \{ E[X_{i}'^{2}\Gamma(Y_{i}'-Z_{i},t)] + E[X_{j}'^{2}\Gamma(Y_{j}'-Z_{j},t)] \} dt$$

$$< 2^{-1} \sum_{i=1}^{n} \sum_{j \in N(C_{i})} \{ E[X_{i}'^{2}\Gamma(Y_{i}'-Z_{i})] + E[X_{j}'^{2}(Y_{j}'-Z_{j})] \} dt$$

$$< 2^{-1} \sum_{i=1}^{n} \sum_{j \in N(C_{i})} \{ E[X_{i}'^{2}(Y_{i}'-Z_{i})] + E[X_{j}'^{2}(Y_{j}'-Z_{j})] \}$$

$$<2^{-1}\sum_{i=1}^{n}\sum_{j\in N(C_{i})}\{E|x_{i}^{2}Y_{i}| + Ex_{i}^{2}E|z_{i}| + E|x_{j}^{2}Y_{j}| + Ex_{j}^{2}E|z_{j}|\}$$

By (4.10), (4.11), (4.12) and (4.13), the proposition is proved.

By using the Hölder and the GM-AM inequalities,  $\beta \leqslant 2100^2 \sum_{i=1}^n \gamma_i$  and Proposition 4.1 yields

$$(4.14) P^{\zeta}(a_{\zeta} \le W \le b_{\zeta}) \le (3/2) \tilde{B}^{-1}(b_{\zeta} - a_{\zeta}) + 2100^{2} \tilde{B}^{-3} \int_{i=1}^{n} \gamma_{i} a.s.$$

Since  $\beta$  contains more terms than are contributed by (4.11), (4.12) and (4.13), the second absolute constant in (4.14) can be reduced by using (4.11), (4.12) and (4.13) directly. The result is the following corollary.

COROLLARY 4.1. Let  $\zeta$ ,  $a_{\zeta}$  and  $b_{\zeta}$  be as in Proposition 4.1. Then

(4.15) 
$$P^{\zeta}(a_{\zeta} \leq W \leq b_{\zeta}) \leq (3/2) \tilde{B}^{-1}(b_{\zeta} - a_{\zeta}) + 34\theta^{2} \tilde{B}^{-3} \int_{i=1}^{n} \gamma_{i} \quad a.s.$$

If  $X_1,\dots,X_n$  are independent, then  $\theta=1$ . By taking  $\Lambda_0=\Lambda=N(\Lambda)=\varphi$ , we have the following classical result.

COROLLARY 4.2. For a < b,

(4.16) 
$$P(a \le W \le b) \le (3/2) B^{-1}(b-a) + 34B^{-3} \sum_{i=1}^{n} \gamma_{i}.$$

Of course the absolute constants in (4.16) can be substantially reduced if we prove the inequality for independent random variables directly.

5. The main theorem. Let  $X_1, \dots, X_n$  be random variables with  $EX_i = 0$ ,  $Var(\sum_{i=1}^n X_i) = B^2$  and  $E|X_i|^3 = \gamma_i < \infty$  and satisfying the dependence assumption of Section 2. Let F be the distribution function of W =  $B^{-1} \sum_{i=1}^n X_i$  and  $\Phi$  the standard normal distribution function.

THEOREM 5.1. There exists an absolute constant C such that for  $1 \le p \le \infty$ ,

(5.1) 
$$||\mathbf{F} - \Phi||_{\mathbf{p}} < C\lambda^{\frac{\mathbf{p} - 1}{\mathbf{p}}} \mathbf{B}^{-3}\beta$$

where  $\lambda$  is given by (2.6) and  $\beta$  by (2.7).

To prove Theorem 5.1 we assume without loss of generality that B=1 and use the identity (3.15). Denote the error terms on the right hand side of (3.15) by  $R_1(f)$  and  $R_2(f)$  respectively. Now choose f in (3.15) to be  $f_{z,\alpha}$ , the unique bounded solution of the differential equation

(5.2) 
$$f'(w) - wf(w) = h_{z,\alpha}(w) - Nh_{z,\alpha}$$

where  $\alpha > 0$  is to be determined later and for  $\alpha > 0$ ,

(5.3) 
$$h_{z,\alpha}(w) = \begin{cases} 1 & , & w < z , \\ 1 + \frac{z-w}{\alpha} & , & z < w < z+\alpha , \\ 0 & , & w > z+\alpha , \end{cases}$$

and  $Nh_{z,\alpha} = (2\pi)^{-1/2} \int_{-\infty}^{\infty} h_{z,\alpha}(x) e^{-x^2/2} dx$ . For  $\alpha = 0$ ,  $h_{z,0} = I_{(-\infty,z]}$  and  $Nh_{z,0} = \Phi(z)$ . Then (3.15) yields

(5.4) 
$$\operatorname{Eh}_{z,\alpha}(W) - \operatorname{Nh}_{z,\alpha} = R_1(f_{z,\alpha}) + R_2(f_{z,\alpha}) .$$

The solution of (5.2) is given by

$$f_{z,\alpha}(w) = e^{w^2/2} \int_{-\infty}^{w} [h_{z,\alpha}(x) - Nh_{z,\alpha}] e^{-x^2/2} dx$$
  
=  $-e^{w^2/2} \int_{w}^{\infty} [h_{z,\alpha}(x) - Nh_{z,\alpha}] e^{-x^2/2} dx$ .

We need the boundedness properties of  $f_{z,\alpha}$  which we state in the following lemmas.

LEMMA 5.1. For all w, z and  $\alpha$ ,

(5.5) 
$$0 < f_{z,\alpha}(w) < 1$$
,

$$\left| wf_{z,\alpha}(w) \right| < 1 ,$$

$$|f_{\mathbf{Z},\alpha}^{\dagger}(\mathbf{w})| < 1 .$$

proof. For  $\alpha=0$ , see Stein (1972) or Stein (1986). For  $\alpha>0$ , use the fact that  $h_{z,\alpha}=\lim_{n\to\infty} n^{-1}\sum_{i=0}^{n-1}h_{z+\alpha n}^{-1}$ .

LEMMA 5.2. For all w, z and 
$$\alpha > 0$$
, 
$$|f'_{z,0}(w+s) - f'_{z,0}(w+t)| < [2\wedge(|s|+|t|)][|w|+1]$$
 
$$+ I(z-s < w < z-t)I(s > t)$$

(5.9) 
$$|f_{z,\alpha}^{\dagger}(w+s) - f_{z,\alpha}^{\dagger}(w+t)| \leq [2\wedge(|s|+|t|)][|w|+1]$$

$$+ \alpha^{-1}|\int_{t}^{s} I(z \leq w+u \leq z+\alpha) du| .$$

PROOF. Use (5.2), (5.5) and (5.7), keeping in mind that the left hand side of (5.8) or (5.9) is at most 2.

LEMMA 5.3. For all w,

(5.10) 
$$\int f_{z,0}(w)dz = 1 ,$$

+  $I(z-t\leq w\leq z-s)I(s\leq t)$ .

PROOF. See Ho and Chen (1978, Lemma 2.1).

LEMMA 5.4. For all w,

PROOF. Use (5.2), (5.10) and (5.11), keeping in mind that the left hand side of (5.12) is at most 2.

We also need a lemma which relates  $\|F-\Phi\|_{\infty}$  to  $\sup_{z \in \mathbb{R}} |Eh_{z,\alpha}(W) - Nh_{z,\alpha}|$ . LEMMA 5.5. For  $\alpha > 0$ ,

(5.13) 
$$\|F-\Phi\|_{\infty} \leq \sup_{\mathbf{z},\alpha} |\operatorname{Eh}_{\mathbf{z},\alpha}(W) - \operatorname{Nh}_{\mathbf{z},\alpha}| + (2\pi)^{-1/2} \alpha .$$

PROOF. See Bolthausen (1984, Section 2) or use the fact that  $\operatorname{Eh}_{z-\alpha}$ , $\alpha$  (W)  $\leq$  F(z)  $\leq$  Eh<sub>z,\alpha</sub>(W) and that  $\Phi(z+\alpha) - \Phi(z) \leq (2\pi)^{-1/2}\alpha$ .

Now we bound the error terms in (5.4). Since  $\|\cdot\|_p^p < \|\cdot\|_{p-1}^{p-1}\|\cdot\|_{1}$  it suffices to prove (5.1) for  $p=\infty$  and 1. To facilitate calculations, we let C denote an absolute constant with possibly different magnitudes at different places. By Lemma 5.2,

$$(5.14)$$
  $|R_1(f_{z,\alpha})|$ 

$$\begin{cases} \sum_{i=1}^{n} E \int [2 \wedge (|Y_{i}| + t| + |t|)] (E^{i} |V_{i}| + 1] |X_{i} \Gamma(Z_{i} - Y_{i}, t)| dt \end{cases}$$

+ 
$$\alpha^{-1}$$
  $\sum_{i=1}^{n}$  E  $\int \{I(Y_i>0) \int_{t}^{Y_i+t} \sum_{p=1}^{q} (Z_i+z \leq W+s \leq Z_i+z+\alpha) ds \}$ 

+ 
$$I(Y_i < 0) \int_{Y_i+t}^{t} P^{i}(Z_i + z < W + s < Z_i + z + \alpha) ds | X_i \Gamma(Z_i - Y_i, t)| dt$$

$$= R_{11} + R_{12}$$
 say,

where  $\zeta_i = (X_i, Y_i, Z_i)$ . Now

$$E^{\zeta_{i}}|V_{i}| \leq E^{\zeta_{i}}|W-T_{i}| + E^{\zeta_{i}}|T_{i}-Z_{i}|$$

$$= E|W-T_{i}| + E^{\zeta_{i}}|T_{i}-Z_{i}|$$

$$\leq E|W| + E|T_{i}| + E^{\zeta_{i}}|T_{i}-Z_{i}|$$

$$\leq 1 + E|T_{i}| + E^{\zeta_{i}}|T_{i}-Z_{i}|.$$

Using this, we have

$$|R_{11}| < C \sum_{i=1}^{n} E \int [|Y_{i}| + 2|t|] |X_{i}\Gamma(Z_{i}-Y_{i},t)| dt$$

$$+ C \sum_{i=1}^{n} E |T_{i}| E \int |X_{i}\Gamma(Z_{i}-Y_{i},t)| dt$$

$$+ C \sum_{i=1}^{n} E \int |T_{i}-Z_{i}| |X_{i}\Gamma(Z_{i}-Y_{i},t)| dt$$

$$< C \sum_{i=1}^{n} \{E|X_{i}Y_{i}(Z_{i}-Y_{i})| + E|X_{i}(Z_{i}-Y_{i})^{2}|$$

$$+ C \sum_{i=1}^{n} E|X_{i}(Z_{i}-Y_{i})| E|T_{i}|$$

$$+ C \sum_{i=1}^{n} E|X_{i}(Z_{i}-Y_{i})|T_{i}-Z_{i}|$$

By Proposition 4.1,

$$|R_{12}| < \lambda \alpha^{-1} \sum_{i=1}^{n} E \int |Y_{i}| (\alpha + \beta) |X_{i}\Gamma(Z_{i} - Y_{i}, t)| dt$$

$$= \lambda \alpha^{-1} (\alpha + \beta) \sum_{i=1}^{n} E |X_{i}Y_{i}(Z_{i} - Y_{i})|$$

$$< C\lambda \alpha^{-1} (\alpha + \beta) \beta .$$

Combining (5.14), (5.15) and (5.16), we obtain

(5.17) 
$$|R_1(f_{z,\alpha})| < c\lambda[1 + \alpha^{-1}(\alpha+\beta)]\beta$$
.

By similar arguments,

(5.18) 
$$|R_2(f_{z,\alpha})| < C\lambda[1 + \alpha^{-1}(\alpha+\beta)]\beta$$

Now choose  $\alpha = \beta$ . Then by (5.4), (5.17), (5.18) and Lemma 5.5, we obtain

This proves (5.1) for  $p = \infty$ .

By Lemma 5.4,

(5.20) 
$$\int |R_{1}(f_{z,0})| dz < \sum_{i=1}^{n} E \int [1_{\Lambda}(|Y_{i}+t|+|t|)] [E^{i}|V_{i}|+2] |X_{i}\Gamma(Z_{i}-Y_{i},t)| dt$$

which is the same as  $R_{11}$  except for the constants and therefore <  $C\beta$ .

Likewise

Combining (5.4), (5.20) and (5.21), we obtain

(5.22) 
$$\|F-\Phi\|_1 \le C\beta$$

This proves (5.1) for p = 1 and therefore theorem 5.1.

REMARKS. (1) It is reasonable to assume that  $\lambda$  remains bounded as  $n+\infty$ . In fact  $\lambda+1$  as  $n+\infty$  if  $\max_{1\leq i\leq n} B^{-2} Var(\sum_{i} X_{j})+0$  as  $n+\infty$  (here  $\lambda$  should be  $\lambda_{n}$ ,  $N(C_{i})$  should be  $N(C_{ni})$  and  $B^{2}$  should be  $B_{n}^{2}$ ), which is an infinitesimality condition and is natural to assume in order for the central limit theorem to hold.

(2) In bounding  $\|F-\Phi\|_{\infty}$ , we could use  $h_{z,0}$  instead of  $h_{z,\alpha}$  in (5.4) and use (5.8) instead of (5.9) of Lemma 5.2. But in doing so we would end up having the quantity  $\max\{1, B^{-2} \sum_{i=1}^{n} E|X_{i}Y_{i}|, B^{-2} \sum_{i=1}^{n} E|X_{i}Z_{i}|\}$  as a factor in the bound.

COROLLARY 5.1. There exists an absolute constant C such that for

1 < p < ∞,

(5.23) 
$$||F-\Phi||_{p} < c\theta^{2}B^{-3} \sum_{i=1}^{n} \gamma_{i}$$

where  $\theta$  is given by (2.5).

PROOF. Without loss of generality we assume B=1. Let  $\rho=\theta^2\sum_{i=1}^n\gamma_i$ . First, by the Hölder and the GM-AM inequalities we obtain  $\frac{p-1}{p}$  (CA)  $\frac{p}{p}$  from (5.1). It suffices to prove (5.23) for  $p=\infty$ . Now choose C>8 so that for  $\rho>\frac{1}{8}$ ,  $\|F-\Phi\|_{\infty}< C\rho$  is trivially true. For  $\rho<\frac{1}{8}$ , we use  $B_i>1-\left[\operatorname{Var}(\sum_i X_j)\right]^{1/2}$  which by the Hölder and the  $\operatorname{jcN}(C_i)$  GM-AM inequalities  $>1-\left[\theta^2\sum_{j\in N(C_i)} Y_i\right]^{1/3}>1-\rho^{1/3}>\frac{1}{2}$ , so that  $\lambda<8$ . This proves the corollary.

If  $x_1, \dots, x_n$  are independent, then  $\theta = 1$  and Corollary 5.1 yields the classical result.

Recall that the distance between two points  $i=(i_1,\ldots,i_d)$  and  $j=(j_1,\ldots,j_d)$  in  $\mathbb{Z}^d$  is defined by  $|i-j|=\max_{1\leq l\leq d}|i_l-j_l|$  and the distance between two subsets A and B of  $\mathbb{Z}^d$  is defined by  $\rho(A,B)=\inf\{|i-j|:i\in A,\ j\in B\}$ . A set of random variables  $\{X_i:i\in \mathbb{Z}^d\}$  is said to be an m-dependent random field if  $X_A$  and  $X_B$  are independent whenever  $\rho(A,B)>m$  where A and B are finite sets.

Suppose  $\{X_i:i\in\mathbb{Z}^d\}$  is an m-dependent random field with  $EX_i=0$  and  $E|X_i|^3=\gamma_i<\infty$ . If  $\{X_1,\ldots,X_n\}=\{X_i:i\in\Lambda\}$  where  $\Lambda\subset\mathbb{Z}^d$ , then in applying Theorem 5.1 or Corollary 5.1, we have  $A_i=\{j:|j-i|< m\}\cap\Lambda$ ,  $B_i=\{j:|j-i|<2m\}\cap\Lambda$ ,  $C_i=\{j:|j-i|<3m\}\cap\Lambda$  and  $N(C_i)=\{j:|j-i|<5m\}\cap\Lambda=N^*(C^i)$ . Therefore  $0<(10m+1)^d.$ 

Now suppose further that  $\{x_i:i\epsilon\mathbb{Z}^d\}$  is stationary and  $\Lambda=\{i\epsilon\mathbb{Z}^d:|i|\leq N\}$ .

Assume  $\text{EX}_{\mathbf{i}}(\sum_{|j-\mathbf{i}| \leq m} X_j) = \tau^2 > 0$ . Then  $B^2 \sim \tau^2 |\Lambda|$ . Let  $Y_{\mathbf{i}} = Y$  for all  $\mathbf{i}$ . The following corollary follows from Corollary 5.1.

COROLLARY 5.2. There exists an absolute constant  $\,$  C such that for sufficiently large  $\,$  N  $\,$ 

$$\left|\left|F-\Phi\right|\right|_{p} \leq C\left(10m+1\right)^{2d} \gamma \tau^{-1} \left|\Lambda\right|^{-1/2}, \quad 1 \leq p \leq \infty \ .$$

Finally we show that Theorem 5.1 gives a rate of  $n^{-1/4}$  if  $X_1, \ldots, X_n$  is the 1-dependence sequence example of Erickson (1974). Recall that the rate of  $n^{-1/4}$  is the correct one. It is not difficult to see that in this example if  $X_i$  is independent of all the other random variables then  $A_i = B_i = C_i = N(C_i) = \{i\}$  and  $X_i = Y_i = Z_i = T_i$ . On the other hand if  $X_i = -X_{i-1}$  or  $-X_{i+1}$ , then  $A_i = B_i = C_i = N(C_i) = \{i, i-1\}$  or  $\{i, i+1\}$  accordingly. In this case  $Y_i = Z_i = T_i = 0$  identically. Consequently

 $\beta \leq C Y$  (number of X which is independent of all the other random variables)

$$= CYB_n^2$$

where  $B_n^2 = Var \left(\sum_{i=1}^n X_i\right)$  and  $\gamma = E \left|X_i\right|^3$  for all i. Similarly,  $\lambda \le C$  for sufficiently large. Hence for  $1 \le p \le \infty$ ,

$$\|F - \Phi\|_{p} \le CB_{n}^{-3} \Upsilon B_{n}^{2} = C\Upsilon B_{n}^{-1}$$

$$\le C\Upsilon n^{-1/4}$$

for sufficiently large n.

Acknowledgement. The author would like to thank Andrew Barbour for simplifying some of the steps in the proof of the concentration inequality in Section 4.

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