## ON DEVIATIONS OF THE SAMPLE MEAN

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1. Introduction Let  $X_1$ ,  $X_2$ , ... be a sequence of independent and identically distributed random variables. Let a be a constant,  $-\infty < a < \infty$ , and for each  $n = 1, 2, \cdots$  let

$$p_n = P\left(\frac{X_1 + \cdots + X_n}{n} \ge a\right).$$

It is assumed throughout the paper that the distribution of  $X_1$  and the given constant a satisfy the conditions stated in the following paragraph. These conditions imply that  $p_n > 0$  for each n, and that  $p_n \to 0$  as  $n \to \infty$ . The object of the paper is to obtain an estimate of  $p_n$ , say  $q_n$ , which is precise in the sense that

$$(2) q_n/p_n = 1 + o(1) as n \to \infty.$$

Let t be a real variable, and let  $\varphi(t)$  denote the moment generating function (m.g.f.) of  $X_1$ , i.e.,  $\varphi(t) = E(e^{iX_1})$ ,  $0 < \varphi \le \infty$ . Define

$$\psi(t) = e^{-t} \varphi(t).$$

Let T denote the set of all values t for which  $\varphi(t) < \infty$ . We suppose that  $P(X_1 = a) \neq 1$ , that T is a non-degenerate interval, and that there exists a positive r in the interior of T such that  $\psi(\tau) = \inf_i |\psi(t)| = \rho$  (say). These conditions are satisfied if, for example,  $\varphi(t) < \infty$  for all t,  $E(X_1) = 0$ , a > 0, and  $P(X_1 > a) > 0$ . In any case, r and  $\rho$  are uniquely determined by

(4) 
$$\frac{\varphi'(\tau)}{\varphi(\tau)} = a \text{ and } \rho = \psi(\tau),$$

where  $\varphi' = d\varphi/dt$ , and we have  $0 < \rho < 1$ .

There are three separate cases to be considered.

Case 1: The distribution function (d.f.) of  $X_1$  is absolutely continuous, or, more generally, this d.f. satisfies Cramér's condition (C) [1, p. 81].

Case 2:  $X_1$  is a lattice variable, i.e., there exist constants  $x_0$  and d>0 such that  $X_1$  is confined to the set  $\{x_0+rd:r=0,\pm 1,\pm 2,\cdots\}$  with probability one.

Case 3: Neither Case 1 nor Case 2 obtains.

We can now state

THEOREM 1. There exists a sequence b. , b1 , ... of positive numbers ba such that

(5) 
$$p_n = \frac{\rho^n}{(2\pi n)^3} b_n [1 + o(1)], \quad \log b_n = O(1)$$

as  $n \to \infty$ . In Cases 1 and 3,  $b_n$  is independent of n. This last also holds in Cax: if  $P(X_1 = a) > 0$ .

The proof of Theorem 1, and of Theorem 2 below, is given in Sections 2: The present determination of b, is given by (4), (9) and (33) in Cases I and 3 and by (4), (8), (37), (38) and (46) in Case 2. The following refinements of Theorem 1 are available in Cases 1 and 2:

Theorem 2. (Cases 1 and 2). For each  $j=1, 2, \cdots$  there exists a bounded (possibly constant) sequence  $c_{i,1}, c_{j,2}, \cdots$  such that, for any given positive integral

(6) 
$$p_n = \frac{\rho^n}{(2\pi n)^4} b_n \left[ 1 + \frac{c_{1,n}}{n} + \frac{c_{2,n}}{n^2} + \cdots + \frac{c_{n,n}}{n^4} \right] \left[ 1 + O\left(\frac{1}{n^{4+1}}\right) \right]$$

as  $n \to \infty$ .

The sequences  $|c_{f,n}|$  are given explicitly for Cases 1 and 2 in Sections 3 and respectively. It would be interesting to know whether (6) holds in Case 3 as well, perhaps with the  $|c_{f,n}|$  determined according to the formula for Case 1.

Estimates in the form (5) or (6) were first obtained by Cramér [2, pp. 20-21] in the case when  $X_i$  has an absolutely continuous component (so that Case I obtains). Cramér showed that in the latter case (6) holds for every k (with k, and each  $c_{j,k}$  independent of n), and determined  $b_k$ . Our method of proof in the general case (cf. Sections 2-5) is essentially a variant or extension of Craméris method. Case 2 was treated recently by Blackwell and Hodges [3] by a different method. It is shown in [3] that (6) holds for k = 1 in Case 2, under the restriction on n and a that  $P(X_1 + \cdots + X_n = na) > 0$  for every admissible n, and the requisite  $b_k$  and  $c_{i,k}$  (which are then independent of n) are determined explicitly Some other references bearing on the problem under consideration are [4], [3] and [6].

In the following Section 2 it is shown that  $p_n$  can be expressed as  $\rho^*I_n$ , where  $I_n$  is a certain integral;  $0 < I_n < 1$ , and  $I_n = O(n^{-1})$  as  $n \to \infty$ .  $I_n$  can be estimated by application of certain refinements [1], [7] of the central limit theorem. This estimation of  $I_n$  is carried out in Sections 3, 4 and 5 for Cases 1, 2 and 3 respectively. It may be added here that, as was pointed out in [2], direct application of the central limit theorem (or refinements thereof) to  $p_n$  defined by (1) does not, in general, yield approximations  $q_n$  which satisfy (2).

In Section 6 we describe certain numerical approximations to p, which are suggested by Theorems 1 and 2 and their proofs.

2. Lemmas. Let  $Y_1 = X_1 - a$ , and let F be the (left-continuous) distribution function (d.f.) of  $Y_1$ ,  $F(y) = P(Y_1 \stackrel{\checkmark}{<} y)$ . Let G be defined by  $G(x) = \int_{-\infty \sqrt{x}} e^{x} e^{x} f(y)$ . Since  $E(e^{x} Y_1) = \psi(\tau) = \rho$ , it is clear that G is a probability df. Let  $Z_1$  be a random variable distributed according to G.

LEMMA 1. The m.g.f. of  $Z_1$  exists in a neighborhood of the origin. We have

(7) 
$$E(Z_1) = 0, \quad 0 < \operatorname{Var}(Z_1) < \infty.$$

PROOF. Let  $\xi(t)$  denote the m.g.f. of  $Z_1$ . Then  $\xi(t) = \psi(\tau + t)/\rho$  for all t

by (3) and the definition of  $Z_1$ . Since  $\psi(t)<\infty$  in a neighborhood of  $t=\tau_1$  if follows that  $\xi(t)<\infty$  in a neighborhood of t=0. Consequently,  $E\mid Z_1\mid^r<\infty$  for  $t=1,2,3,\cdots$  and  $E(Z_1)=\{u^r\}/u^r\}$ . In particular,  $E(Z_1)=\{u^t\}/u^r\}/u^r$  is in the interior of T. It remains to show that  $\text{Var}(Z_1)>0$ . Suppose to the contrary that  $\text{Var}(Z_1)=0$ ; then  $P(Z_1=0)=1$ ; hence  $P(Y_1=0)=1$ , i.e.,  $P(X_1=a)=1$ , which is contrary to our assumptions. This completes the proof.

Let  $Var(Z_1)$  be denoted by  $\sigma^2$ . It follows from the preceding paragraph and (4) that

(8) 
$$\sigma^2 = \frac{\varphi''(\tau)}{\varphi(\tau)} - \alpha^2.$$

Define

(9) 
$$\alpha = \sigma r$$
,  $(0 < \alpha < \infty)$ .

Let  $Z_1$ ,  $Z_1$ ,  $\cdots$  be a sequence of independent and identically distributed random variables. For each n, let

$$U_n = \frac{Z_1 + \cdots + Z_n}{n!\sigma}$$

and

(11) 
$$H_n(z) = P(U_n < z), \quad (-\infty < z < \infty).$$

LEMMA 2.  $p_n = \rho^n I_n$ , where

(12) 
$$I_n = n^{\frac{1}{2}} \alpha \int_0^n e^{-n^{\frac{1}{2}} x} |H_n(x) - H_n(0)| dx.$$

Proof. Let  $Y_i = X_i - a$  for  $j = 1, 2, \dots, n$ . Then

$$p_n = P(Y_1 + \cdots + Y_n \ge 0) \qquad \text{by (1)}$$

$$= \int_{\substack{x_1 + \cdots + x_n \ge 0 \\ x_1 + \cdots + x_n \ge 0}} dF(y_1) \cdots dF(y_n)$$

$$= \rho^n \int_{x_1 + \cdots + x_n \ge 0} e^{-i(x_1 + \cdots + x_n)} dG(x_1) \cdots dG(x_n)$$

$$= \rho^n \int_{0 \le x < n} e^{-n \log x} dH_n(x) \qquad \text{by (9), (10), (11)}$$

$$= \rho^n \int_{0 \le x < n} e^{-n \log x} dH_n(x)$$

$$= \rho^n \int_{0 \le x < n} e^{-n \log x} dH_n(x)$$

It follows by integration by parts that  $I_n^{\bullet}$  defined in (13) is equal to  $I_n$ , and this completes the proof.

A theorem of Chernoff [4] states that  $p_n \leq p^n$  for every n, and that for any

given positive  $\rho_0 < \rho$ , we have  $p_n \ge \rho_0^n$  for all sufficiently large n. A simple proof of Chernoff's theorem can be given as follows. Since  $0 \le H_n(x) - H_n(0) \le 1$  for every n and  $x \ge 0$ , we have  $I_n \le 1$  and hence  $p_n \le \rho^n$  for every n, by Lemma 2. To establish the second part of the theorem, we note first that  $\lim_{n\to\infty} H_n(x) = \Phi(x)$  for every x, where

(14) 
$$\Phi(x) = \int_{-\infty}^{x} (2\pi)^{-1} e^{-\frac{1}{2}t^2} dt \qquad (-\infty < x < \infty),$$

by (7), (10), (11) and the central limit theorem. Let  $\epsilon$  be a positive constant. Then

$$I_n \ge n^b \alpha \int_{\epsilon}^n e^{-n \beta n x} \{H_n(x) - H_n(0)\} dx$$

$$\ge [H_n(\epsilon) - H_n(0)] n^b \alpha \int_{\epsilon}^n e^{-n \beta n x} dx$$

$$= [H_n(\epsilon) - H_n(0)] e^{-n \beta n x}.$$

Hence  $\lim \inf_{n\to\infty} \{n^{-1} \log I_n\} \ge -\alpha \epsilon$ . Since  $I_n \le 1$  for every n, and since  $\epsilon$  is arbitrary, it follows that  $n^{-1} \log I_n = o(1)$ . Hence  $n^{-1} \log p_n = \log \rho + o(1)$ , by Lemma 2, and this is equivalent to the conclusion desired.

The preceding argument depends only on the central limit theorem. In the following sections we estimate  $I_n$  more accurately by substituting the expansions of  $H_n(x)$  due to Cramér [1] and Esseen [7] in the right side of (12). The remainder of this section is concerned with preparations for this application of the Cramér-Esseen expansions. Almost all the considerations of the following paragraphs are well known, and we include them here only for the sake of completeness.

Let  $\eta(w)$  denote the m.g.f. of  $Z_1/\sigma$ . According to Lemma 1,  $\eta < \infty$  in a neighborhood of w = 0. For  $j = 2, 3, \cdots$  let  $\lambda_j$  be defined by

(15) 
$$\lambda_1 = \frac{1}{2}; \quad \lambda_j = (j|\sigma^j)^{-1}(d^j/dt^j)|\log \varphi(t)|_{i=1}, \quad (j = 3, 4, \cdots).$$

It should be noted that  $j|\lambda_j$  is the jth cumulant of the distribution of  $Z_1/\sigma$ . The m.g.f. of  $U_n$ , with  $U_n$  defined by (10), is

$$[\eta(w/n^{\dagger})]^{n} = \exp \left[n \sum_{i=1}^{n} \lambda_{i} (w/n^{\dagger})^{t}\right].$$

Clearly,  $[n(w/n^{1})]^{n}$  exp  $(-w^{2}/2)$  is analytic in a domain independent of n, and can be expanded there as a power series in w. By regrouping the terms of this series according to powers of n we shall have

(16) 
$$[\eta(w/n^{i})]^{*}e^{-x^{i}n} = \sum_{i=0}^{\infty} n^{-i}i P_{i}(w)$$

where the  $P_j$  are polynomials.  $P_j$  is of degree 3j, and  $P_j$  is even or odd according as j is even or odd. The first few polynomials are

$$P_1(w) = \lambda_1 w^3,$$
(17) 
$$P_1(w) = \lambda_4 w^4 + \frac{1}{2} \lambda_3^2 w^6,$$

$$P_1(w) = \lambda_4 w^4 + \lambda_1 \lambda_4 w^7 + \frac{1}{2} \lambda_3^2 w^6,$$

$$P_4(w) = \lambda_4 w^4 + (\frac{1}{2} \lambda_4^2 + \lambda_1 \lambda_4) w^4 + \lambda_4^2 \lambda_4 w^{10} + \frac{1}{47} \lambda_4^4 w^{10}.$$

Write  $\Phi^{(0)}(x) = \Phi(x)$  and  $\Phi^{(1)}(x) = (d'/dx')\Phi(x)$  for  $r = 1, 2, \cdots$ , where  $\Phi$  is given by (14). Let  $P_I(-\Phi)$  denote the function of x obtained by replacing w' with  $(-1)'\Phi^{(1)}(x)$  in the polynomial  $P_I(w)$ . It is clear that each  $P_I(-\Phi)$  is absolutely continuous and of bounded variation in  $(-\infty, \infty)$ . It should also be noted that  $P_I(-\Phi)$  is square integrable with respect to Lebesgue measure.

In the following, for any function K(x) of bounded variation in  $(-\infty, \infty)$ , we denote the c.f. of K by  $\chi(t | K)$ , i.e.,

(18) 
$$\chi(t|K) = \int_{-\infty}^{\infty} e^{itx} dK(x)$$

 $P_0(w) = w^0 = 1$ 

for every real t. If K is absolutely continuous,  $\chi$  is, of course,  $(2\pi)^t$  times the Fourier transform of K'. The reader may refer to [8, Chapters I-III] for such elements of Fourier transform theory as are used in this page.

LEMMA 3. For every j, t, and x

(19) 
$$\chi(t \mid P_i(-\Phi)) = P_i(it) e^{-it}$$

and

(20) 
$$P'_{j}(-\Phi) = (2\pi)^{-\frac{1}{2}} \int_{-\pi}^{\pi} e^{-itx} P_{j}(it) d\Phi(t).$$

PROOF. As is pointed out in [1, p. 49], we have

(21) 
$$\chi(t \mid \Phi^{(r)}) = (-it)^r e^{-it^{it}}$$

for  $r=0,1,\cdots$ . Suppose, for given j, that  $P_j(w)=\sum_{r=0}^N a_r w^r$ , where the  $a_r$  and N are constants (depending on j). Then  $P_j(-\Phi)=\sum_{r=0}^N a_r (-1)^r \Phi^{(j)}(x)$ ; hence the left side of (19) equals  $\sum_{i=0}^N a_i (-1)^r \chi(t \mid \Phi^{(i)})$ ; (19) now follows from (21). The relation (20) follows from (19) by the inversion formula for the Fourier transform, since  $dP_j(-\Phi)=P_j^r(-\Phi)\,dx$ , and  $d\Phi(t)=(2\pi)^{-1}e^{-t^2}\,dt$ .

A probability d.f. K(x) is said to satisfy condition (C) if

$$\lim \sup_{t \to \infty} |\chi(t \mid K)| < 1.$$

In the following lemma the  $F_i$  are arbitrary probability d.fs.

LEMMA 4. If  $F_1$  satisfies (C), and if  $F_1$  is absolutely continuous with respect to  $F_2$ , then  $F_1$  also satisfies (C).

PROOF. In this proof, for any probability d.f. K let  $K^*$  denote the symmetrized d.f. defined by  $K^*(x) = \int_{-\infty}^{\infty} K(x+y) dK(y)$ . We then have

$$\chi(t \mid K^{\bullet}) = \int_{-\infty}^{\infty} \cos(tx) dK^{\bullet} = |\chi(t \mid K)|^{2}$$

for all i.

Suppose, contrary to the lemma, that there exists a sequence  $\{i, j = 1, 2, \cdots\}$  such that  $|t_f| \to \infty$  and  $|\chi(t_f|F_f)| \to 1$  as  $j \to \infty$ . It then follows from the above paragraph with  $K = F_f$  that  $\int_{-\infty}^{\infty} \cos(tx) df_f^2 \to 1$ . Hence  $\cos(tx) \to 1$  in  $F_1^*$ -measure. Since  $F_1$ -measure dominates  $F_1$ -measure, it is easily seen that  $F_2^*$ -measure dominates  $F_1^*$ -measure. Consequently,  $\cos(tx) \to 1$  in  $F_1^*$ -measure. It now follows from the above paragraph with  $K = F_1$  that  $|\chi(t_f|F_1)|_{f=1}^2$  as  $j \to \infty$ , which is impossible. This completes the proof.

We conclude this section with a description of the functions  $S_1(x)$ ,  $S_1(x)$  which occur in the Euler-Maclaurin sum formulae, and which are required in the analysis of Case 2. It is convenient to define  $S_1$  as follows:

(22) 
$$S_1(x) = \frac{1}{2} - x$$
 for  $0 < x \le 1$ ;  $S_1(x+1) = S_1(x)$ .

For  $i \geq 2$ ,  $S_i$  may be defined as

(23) 
$$S_{j}(x) = \begin{cases} \frac{1}{2^{j-1}} \sum_{r=1}^{\infty} \frac{\cos(2\pi rx)}{(\pi r)^{j}} & (j \text{ even}) \\ \frac{1}{2^{j-1}} \sum_{r=1}^{\infty} \frac{\sin(2\pi rx)}{(\pi r)^{j}} & (j \text{ odd}). \end{cases}$$

Each  $S_j$  is a bounded and periodic function;  $S_j$  is absolutely continuous for  $j \ge 2$ ; and at each non-integral x we have

$$(24) S_1'(x) = -1, S_{j+1}'(x) = (-1)^j S_j(x) (j = 1, 2, \cdots).$$

3.  $I_n$  in Case 1. Suppose that the d.f. of  $X_1$  satisfies (C). Since  $Y_1 = X_1 - a$ , it is plain that  $P_1$  the d.f. of  $Y_1$ , also satisfies (C). It is easily seen that P and G (the d.f. of  $Z_1$ ) are absolutely continuous with respect to each other. It therefore follows from Lemma 4 with  $P_1 = P$  and  $P_2 = G$  that G also satisfies (C).

Let k be an arbitrary but fixed positive integer. It follows from the conclusion of the preceding paragraph by Cramér's theorem [1, p. 81] that  $H_*(z) = K_*(z) + R_*(z)$ , where

(25) 
$$K_n(z) = \sum_{j=0}^{k} n^{-jj} P_j(-\Phi)$$

and  $R_n(x)$  is of the order  $n^{-(k+1)/2}$  uniformly in x. It follows hence from (12) that

(26) 
$$I_n = n^{i} \alpha \int_0^{\infty} e^{-n^{i} a x} \left[ K_n(x) - K_n(0) \right] dx + O(n^{-i k - 1}).$$

We have

(27) 
$$\chi(t \mid K_n) = \sum_{i=1}^k n^{-ij} P_i(it) e^{-it^2}$$

by (19) and (25). Let 
$$f_n(x) = \exp(-n^{\frac{1}{2}}\alpha x)$$
 for  $x \ge 0$  and  $f_n(x) = 0$  other-

sise. Then  $\int_{-\infty}^{\infty} e^{ix} f_n(x) dx = 1/(n^i \alpha - it) = g_n(t)$  say. Consequently, by first using integration by parts and then Parseval's formula, it follows that

(28) 
$$n^{\dagger}\alpha \int_{0}^{\infty} e^{-n^{\dagger}\alpha x} \left[K_{n}(x) - K_{n}(0)\right] dx = \int_{0}^{\infty} e^{-n^{\dagger}\alpha x} K'_{n}(x) dx \\ = \int_{-\infty}^{\infty} f_{n}(x) K'_{n}(x) dx - \frac{1}{2\pi} \int_{-\infty}^{\infty} \frac{g_{n}(t)}{g_{n}(t)} \chi(t \mid K_{n}) dt.$$

It follows from (26), (27) and (28) that

(29) 
$$\alpha(2\pi n)^{\frac{1}{2}} I_{\pi} = \int_{-\pi}^{\pi} \left(1 + \frac{it}{n!a}\right)^{-1} \left(\sum_{j=0}^{k} n^{-jj} P_{j}(it)\right) d\Phi(t) + O(n^{-k}).$$

Define

(30) 
$$\mu_{r,s} = \int_{-\infty}^{\infty} (it)^r P_s(it) d\Phi(t)$$
  $(r, s = 0, 1, 2, \cdots).$ 

Since P, is an even [odd] polynomial if s is even [odd], and since  $\int_{-\infty}^{\infty} t^{1/4} d\Phi(t) = 0$  for  $j = 0, 1, 2, \cdots$ , it follows that each  $\mu_{r,s}$  is a real constant, and that

31) 
$$\mu_{r,s} = 0 \quad \text{if } r + s \text{ is odd.}$$

Now,  $(1 + i\hbar^{-1}\alpha^{-1})^{-1} = \sum_{\delta \leq c, k} (-i\hbar^{-1}\alpha^{-1})^c + n^{-ik} t^k \omega_n(t)$ , where  $|\omega|$  is bounded in n and t. Since  $\Phi$  has finite moments of all orders, it therefore follows from (29), (30) and (31) that

(32) 
$$\alpha(2\pi n)^{\frac{1}{2}} I_n = \sum_{0 \le j \le \frac{1}{2}k} n^{-j} \left\{ \sum_{r+i=2j} \left( -\frac{1}{\alpha} \right)^r \mu_{r,i} \right\} + O(n^{-jk}).$$

Since  $p_n = p^n I_n$ , and since  $\mu_{0,0} = 1$ , it follows by replacing k with 2k + 2 in (32) that (6) holds for any given k, with

$$(33) b_n = \alpha^{-1}$$

and

(34) 
$$c_{j,n} = \sum_{\tau+i=1}^{r} \left(-\frac{1}{\alpha}\right)^r \mu_{\tau,s}$$
  $(j = 1, 2, \cdots)$ 

for every n. This establishes Theorem 2, and hence also Theorem 1, in Case 1. It follows from (17) and (30) that the coefficients  $\mu_{r,s}$  required to compute  $\epsilon_{s,a}$  according to (34) are

(35) 
$$\mu_{1,1} = 3\lambda_1$$

$$\mu_{0,1} = 3\lambda_4 - \frac{15}{2}\lambda_1^3$$

where the  $\lambda_i$  are given by (15). Similarly,  $c_{2,n}$  can be computed from

$$\mu_{4,0} = 3$$

$$\mu_{1,1} = -15\lambda_4$$

$$(36) \qquad \mu_{1,3} = -15\lambda_4 + 105\lambda_4^3$$

$$\mu_{1,4} = -15\lambda_4 + 105\lambda_4\lambda_4 - \frac{315}{2}\lambda_4^4$$

$$\mu_{4,4} = -15\lambda_4 + 105(\frac{1}{2}\lambda_4^4 + \lambda_2\lambda_4) - \frac{945}{2}\lambda_4^3\lambda_4 + \frac{10395}{24}\lambda_4^4$$

We conclude this section with a remark concerning the role of Cramér's theorem [1, p. 81] in the preceding argument. Suppose that  $H_n$  is absolutely continuous, and that  $H_n'$  is sequare integrable over  $(-\infty, \infty)$ . It then follows, by integrating (12) by parts and using Parseval's formula, that

(29°) 
$$\alpha(2\pi n)^{\frac{1}{2}} I_n = \int_{-\infty}^{\infty} \left(1 + \frac{it}{n^{\frac{1}{2}}\alpha}\right)^{-1} \left\{ \left[\eta\left(\frac{it}{n^{\frac{1}{2}}}\right)\right]^{\alpha} \delta^{\frac{1}{2}t^{\frac{2}{2}}} \right\} d\Phi(t)$$

where  $\eta$  is, as before, the m.g.f. of  $Z_1/\sigma$ . (The square integrability condition is imposed here for the validity of Parseval's formula, and can be replaced by others, e.g., that  $(1+i^2)^{-1} \mid \eta(ii) \mid be$  integrable). According to (16), the function in curly brackets on the right side of (29°) can be expressed as  $\sum_{i=1}^{n} n^{-i} P_i(ii)$ . By comparing (29) and (29°) it is seen that, from a technical point of view, the role of Cramér's theorem in the present special case is to guarantee that when  $\sum_{i=1}^{n} n^{-i} P_i$  is replaced by  $\sum_{i=1}^{n} n^{-i} P_i$  on the right side of (29°), the error introduced is indeed of the order  $n^{-13}$ . The same remark, but with (29°) replaced by a rather different formula for  $I_n$ , applies to the role of Esseen's theorem in the argument of the following section.

**4.**  $I_n$  in Case 2. Suppose that  $X_1$  is a lattice variable. Let d be the maximum span of  $X_1$ , i.e., d > 0 is the g.c.d. of the differences between consecutive possible values of  $X_1$ . Let  $x_0$  be the number such that  $a \le x_0 < a + d$ , and such that the possible values of  $X_1$  are included in the set  $\{x_0 + rd : r = 0, \pm 1, \pm 2, \cdots\}$ . Let

$$\beta = d/\sigma, \quad \gamma = rd, \quad \kappa = (x_0 - a)/d$$

It should be noted that  $0 \le x < 1$ . For each n, let

(38) 
$$\theta_{\Lambda} = \pi_{K} - [\pi_{\Lambda}], \qquad 0 \leq \theta_{\Lambda} < 1,$$

where [x] denotes the greatest integer contained in x.

Let k be an arbitrary but fixed positive integer. It follows from Esseen's theorem for the lattice case [7, p. 61] that  $H_n(x) = K_n(x) + L_n(x) + R_n(x)$ , where  $K_n(x)$  is given by (25),  $R_n$  is of the order  $n^{-(k+1)/3}$  uniformly in x, and L.

is defined as follows. For any  $j=1,2,\cdots$  let  $h_j=1$  if  $j\equiv 1$  or 2 (mod 4) and  $h_j=-1$  if  $j\equiv 0$  or 3 (mod 4). Then

(39) 
$$L_n(x) = \sum_{j=1}^k n^{-\frac{1}{2}j} h_j \beta^j S_j(n^{\frac{1}{2}}\beta^{-1}x - \theta_n) K_n^{(j)}(x) = \sum_{j=1}^k M_{j,n}(x) \text{ say,}$$

where  $K_n^{(j)}$  is the jth derivative of  $K_n$ . It follows hence from (12) that

$$I_n = n^i \alpha \int_0^\infty e^{-n^i \alpha x} \{K_n(x) - K_n(0)\} dx$$

$$+ \sum_{i=1}^k n^i \alpha \int_0^\infty e^{-n^i \alpha x} \{M_{j,n}(x) - M_{j,n}(0)\} dx + O(n^{-k-1}).$$

The first term on the right side of (40) is (cf., (28)) equal to  $\int_0^\infty e^{-n^4 cx} K_n^{(1)}(x) dx$ . We observe next that, for  $i \ge 2$ .

$$n^{1} \alpha \int_{0}^{\infty} e^{-n^{1}\alpha x} [M_{j,n}(x) - M_{j,n}(0)] dx = \int_{0}^{\infty} e^{-n^{1}\alpha x} M'_{j,n}(x) dx$$

$$= n^{-ij} h_{j} \beta^{j} \int_{0}^{\infty} e^{-n^{1}\alpha x} [S_{j}(y_{n}) K_{n}^{(j+1)}(x)$$

$$+ (-1)^{j-1} n^{i} \beta^{-1} S_{j-1}(y_{n}) K_{n}^{(j)}(x)] dx$$

$$= n^{-ij} h_{j} \beta^{j} \int_{0}^{\infty} e^{-n^{1}\alpha x} S_{j}(y_{n}) K_{n}^{(j+1)}(x) dx$$

$$- n^{-i(j-1)} h_{j-1} \beta^{j-1} \int_{0}^{\infty} e^{-n^{1}\alpha x} S_{j-1}(y_{n}) K_{n}^{(j)}(x) dx$$

$$= N_{j,n} - N_{j-1,n} (say).$$

In (41), we have put  $n^{\dagger}\beta^{-1}x - \theta_n = y_n$ , and used integration by parts, (24), and the identity  $(-1)^{j}h_j = h_{j-1}$ . In order to evaluate the contribution of  $M_{1,n}$  to the right side of (40), suppose for the moment that  $0 < \theta_n < 1$ , and let

(42) 
$$\zeta_0 = 0, \quad \zeta_r = (r - 1 + \theta_n)\beta/n! \quad (r = 1, 2, \cdots).$$

Let A, denote the open interval  $(f, , f_{r+1})$ . Then  $S_1(y_n)$  is linear in x over each A, (cf. (22)), and its derivative there equals  $-n^n \beta^{-1}$ . By writing  $\int_0^\infty = \sum_{r=0}^\infty \int_{A_r}^\infty$  and applying integration by parts to  $f_{A_r}$ , it follows without difficulty that

(43) 
$$n^{1}\alpha \int_{0}^{\pi} e^{-n^{1}\alpha x} M_{1,n}(x) dx = -\int_{0}^{\pi} e^{-n^{1}\alpha x} K_{n}^{(1)}(x) dx + M_{1,n} + M_{1,n}(0) + \beta n^{-1} \sum_{i}^{n} e^{-i(\nu-i+\Phi_{n})} K_{n}^{(1)}(f_{i}),$$

where  $\gamma = \alpha \beta = rd$  (cf. (37)). Now,  $S_1(x)$  is a left-continuous function of x. It follows hence that, for given n, the left and right sides of (43) are right-

continuous in  $\theta_n$ . Since (43) holds for each  $\theta_n$  in (0, 1), we conclude that (43) is valid for  $\theta_n = 0$  also.

Since  $S_k$  and  $K_n^{(k+1)}$  are bounded functions, it is plain from the definition of  $N_{f,n}$  (cf. (41)) that  $N_{k,n}$  is of the order  $n^{-k-1}$ . It therefore follows from (40), (41) and (43) that

(44) 
$$I_{n} = \beta n^{-1} \sum_{r=1}^{n} e^{-\gamma(r-1+\theta_{n})} K_{n}^{(1)}(\zeta_{r}) + O(n^{-|k-1|}).$$

Now, according to (20) and (25),

(45) 
$$K_n^{(1)}(\xi_r) = (2\pi)^{-\frac{1}{2}} \int_{-\pi}^{\pi} e^{-itt_r} \left( \sum_{i=0}^{k} n^{-\frac{1}{2}i} P_j(it) \right) d\Phi(t)$$

for every r. Let us write

(46) 
$$z = e^{-\tau}, \quad b_n = [\beta/(1-z)]z^{\theta_n}$$

It follows from (44) and (45) that

$$b_n^{-1} (2\pi n)^1 f_n = \int_{-\infty}^{\infty} \frac{(1-z) \exp[-it\beta\theta_n/n^4]}{(1-z \exp[-it\beta/n^4])} \cdot \left(\sum_{i=1}^k n^{-kj} P_j(it)\right) d\Phi(t) + O(n^{-kj}).$$

For any  $\theta$  and any  $j = 0, 1, 2, \cdots$  let

$$f_j(\theta) = \frac{1}{j!} \left\{ \frac{d^j}{dw^j} \left( \frac{(1-z)e^{-\theta w}}{(1-ze^{-w})} \right) \right\}_{w=0}.$$

It then follows easily from (31) and (47) that

(49) 
$$b_n^{-1}(2\pi n)^{\frac{1}{2}}I_n = \sum_{0 \le j \le k/2} n^{-j} \{ \sum_{r+s=2j} \beta' \ell_r(\theta_n) \mu_{r,s} \} + O(n^{-jk}).$$

By replacing k with 2k + 2 in (49) we see that (6) holds for any given k, with  $b_n$  given by (46), and

(50) 
$$c_{f,n} = \sum_{r,k=n,j} \beta' \, f_r(\theta_n) \, \mu_{r,k} \, .$$

This establishes Theorem 2 in Case 2, and hence also the first part of Theorem 1. To complete the proof of Theorem 1 in Case 2, we see from (37), (38) that  $P(X_1 + \cdots + X_n = na) > 0$  implies  $\theta_n = 0$ . Consequently if  $P(X_1 = a) > 0$  then  $\theta_n = 0$  for every n, and hence  $b_n = \beta/(1 - z)$  for every n.

It may be worthwhile to note that in the present case  $b_n$  can be expressed as  $\alpha^{-1}[\gamma e^{n(1-\delta_n)}/(c^n-1)]$ , which shows that, in general,  $b_n$  oscillates about the value  $\alpha^{-1}$  (cf. (33)) as  $n \to \infty$  through the sequence 1, 2, ...

An alternative formula for the coefficients t, required in (50) is

(51) 
$$t_{j}(\theta) = (-1)^{j} \sum_{r \neq z = j} \frac{\theta'}{r! \delta!} \left\{ (1-z) \left( z \frac{d}{dz} \right)^{s} (1-z)^{-1} \right\}.$$

From (51) it is easily seen that, with u = s/(1-s),

$$\begin{aligned} & f_1 = 1 \\ & f_2 = -(\theta + u) \\ & (52) \quad f_3 = \frac{1}{2}[(\theta + u)^3 + u(1 + u)] \\ & f_4 = -\frac{1}{2}[(\theta + u)^4 + 3u(1 + u)\theta + u(1 + u)(1 + 5u)] \\ & f_4 = \frac{1}{2}[(\theta + u)^4 + 8u(1 + u)\theta^3 + 4u(1 + u)(1 + 5u)\theta \\ & + 23u^4 + 36u^3 + 14u^2 + u]. \end{aligned}$$

The coefficients  $c_{1,n}$  and  $c_{1,n}$  can be computed from (35), (36), (50) and (52). The formulae for  $b_n$  and  $c_{1,n}$  with  $\theta = 0$  agree with the results of [3].

5.  $I_n$  in Case 3. If  $X_1$  is not a lattice variable, then neither is  $Z_1$ . It follows hence from a theorem of Easeen [7, p. 49] that  $H_n(x) = \Phi(x) + n^{-1}f(x) + n^{-1}r_n(x)$ , where  $f(x) = (\text{const.}) (1 - x^2) \exp{(-\frac{1}{2}x^2)}$ , and  $r_n(x) = 0$  uniformly in x as  $n \to \infty$ . The contribution of  $n^{-1}f$  to  $I_n$  is  $n^{-1}f_n^{\infty}e^{-(n)^{\log x}}f'(x)$  dx, which is easily seen to be of the order  $n^{-1}$ . It follows that

(53) 
$$I_{n} = n^{1}\alpha \int_{0}^{\infty} e^{-n^{1}\alpha x} \left[\Phi(x) - \Phi(0)\right] dx + o(n^{-1})$$

$$\approx \int_{0}^{\infty} e^{-n^{1}\alpha x} \Phi'(x) dx + o(n^{-1})$$

$$\approx e^{\ln n^{2}} \left[1 - \Phi(n^{1}\alpha)\right] + o(n^{-1})$$

$$= (2\pi n)^{-1}\alpha^{-1} + o(n^{-1}).$$

In (53), we have used integration by parts, a linear change of variable, and the leading term of the asymptotic formula [9, p. 179]

$$(54) \quad 1 - \Phi(x) \, = \, (2\pi)^{-1} e^{-1x^2} \{ x^{-1} - x^{-1} + 3x^{-1} + O(x^{-1}) \} \text{ as } x \to \infty \, .$$

It follows from (53) that (5) holds, with  $b_n = \alpha^{-1}$  for every n. This completes the proof of Theorem 1.

Since  $\Phi(x) + n^{-1}f(x) = K_n(x)$ , where  $K_n$  is defined by (25) with k = 1, the conclusion of the preceding paragraph is also available from the argument of Section 3. We have used a direct calculation instead because this calculation suggests the form of the numerical approximations described in the following section.

8. Concluding remarks. Suppose, in a given case, and for given n and a, that it is required to compute the numerical value of p, defined by (1). In this section we consider approximations of the form

(55) 
$$q_n = \rho^n e^{i v_n^2} ((1 - \Phi(v_n)),$$

where  $\rho$  and  $\Phi$  are defined by (4) and (14), and  $\nu_n$  is a suitably chosen number.

We shall describe four choices of  $v_n$ , called  $v_n^{\bullet}$ ,  $v_n^{(0)}$ ,  $v_n^{(1)}$ , and  $v_n^{(1)}$ . The resulting values of  $q_n$  are denoted by  $q_n^{\bullet}$ ,  $q_n^{(0)}$ , etc.

First consider

$$v_n^* = n^{\dagger} \alpha$$

where  $\alpha$  is given by (4), (8) and (9). This choice of  $v_n$  amounts (cf. (53)) to approximating  $I_n$  by replacing  $H_n$  with  $\Phi$  on the right side of (12). It therefore follows from the Esseen-Berry theorem that we always have

(57) 
$$|p_{*} - q_{*}^{\bullet}| \leq 2C \frac{\rho^{\bullet}}{n^{1}} \frac{E |Z_{1}|^{1}}{\sigma^{1}}$$

where C is a universal constant. Wallace [10, p. 637] states that  $C \le 2.05$ . Next. consider

$$v_n^{(0)} = n^{\frac{1}{2}}/b_n$$

where  $b_n$  is defined by (33) in Cases 1 and 3, and by (46) in Case 2. (Of course,  $q_n^{(0)} = q_n$  in Cases 1 and 3). Then  $q_n^{(0)}$  satisfies (2), and the o(1) term in (2) is known to be of the order  $n^{-1}$  in Cases 1 and 2. Finally, let  $c_{f,n}$  be defined according to Section 4 in Cases 1 and 3, and according to Section 5 in Case 2. Define

(59) 
$$v_n^{(1)} = v_n^{(0)} \left[1 - (b_n^2 + c_{1,n})/n\right]$$

if the expression within the square brackets is positive and  $v_n^{(1)} = 0$  otherwise; and

(60) 
$$v_n^{(1)} = v_n^{(1)} \left[1 + (b_n^4 + c_{1,n}^2 - b_n^2 c_{1,n} - c_{2,n})/n^4\right]$$

if the expression in square brackets is positive and  $v_n^{(1)}=0$  otherwise. Then  $q_n^{(1)}$  also satisfies (2), and  $o(1)=O(n^{-j-1})$  in Cases 1 and 2 (j=1,2). The stated theoretical properties of the approximations  $q_n^{(1)}$  are easy consequences of (5), (6), (54), and (58).

Although (unlike  $q_n^*$ ) the approximations  $q_n^{(l)}$  are derived from asymptotic expansions corresponding to the case when  $n \to \infty$  and a is held fixed, the usefulness of these approximations may be wider than is suggested by the derivation. Some evidence to this effect is provided by the fact that if  $X_1$  is normally distributed then  $p_n = q_n^{(l)} = q_n^{(l)} = q_n^{(l)}$  for every admissible a and every n.

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