WEAK CONVERGENCE OF A TWO-SAMPLE EMPIRICAL PROCESS AND A NEW APPROACH TO CHERNOFF-SAVAGE THEOREMS¹

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- **0.** Summary. An empirical stochastic process for two-sample problems is defined and its weak convergence studied. The results are based upon an identity which relates the two-sample empirical process to the more usual one-sample empirical process. Based on this identity a relatively simple proof of a Chernoff-Savage theorem is obtained. The *c*-sample analogues of these results are also included.
- **1.** Introduction. Let X_1, \dots, X_m, \dots , and Y_1, \dots, Y_n, \dots be independent sequences of independent random variables (rv) with common continuous distribution functions (df) F and G respectively. Let N=m+n and $\lambda_N=m/N$ and suppose that $0<\lambda_*\leq \lambda_N\leq 1-\lambda_*<1$ for some $\lambda_*>0$. (Here, and throughout, when N is used as a subscript it denotes the pair (m,n).) Let $\Lambda=[\lambda_*,1-\lambda_*]$.

When dealing with a single sample X_1, \dots, X_m most distribution-free tests and statistics are based upon the empirical process $\{U_m(t): 0 \le t \le 1\}$ defined by

$$(1.1) U_m(t) = m^{\frac{1}{2}} [F_m F^{-1}(t) - t],$$

where F_m is the empirical df of the sample and where for any two df's, F, G we write FG^{-1} for the composed function $FG^{-1}(t) = F(G^{-1}(t))$. (Throughout this paper we take inverse functions to be left continuous; thus $F^{-1}(t) = \inf\{x\colon F(x) \geq t\}$.) These 1-sample empirical processes have been studied extensively. Much is known about them including their weak convergence to a separable tied-down Wiener process $\{U_0(t): 0 \leq t \leq 1\}$; that is, to a Gaussian process with $E[U_0(t)] = 0$ and $E[U_0(s)U_0(t)] = s(1-t), 0 \leq s \leq t \leq 1$.

Similarly for the second sample, let G_n denote the empirical df of Y_1 , \cdots , Y_n and let $\{V_n(t): 0 \le t \le 1\}$ be the corresponding empirical process

$$(1.2) V_n(t) = n^{\frac{1}{2}} [G_n G^{-1}(t) - t].$$

By assumption, the U_m - and V_n -processes are independent. Let $\{V_0(t): 0 \le t \le 1\}$ denote a second tied-down Wiener process independent of the U_0 -process.

Define

(1.3)
$$K_N = FH_N^{-1}, \quad K_\lambda = FH_\lambda^{-1}, \quad K = K_{\lambda_N} = FH^{-1},$$

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where

$$(1.4) H_N = \lambda_N F_m + (1 - \lambda_N) G_n, H_\lambda = \lambda F + (1 - \lambda) G, H = H_{\lambda_N}.$$

Most non-parametric tests and statistics for the 2-sample problem are based upon the relative ranks of one of the samples. Suppose we let $\{R_{Ni}: 1 \leq i \leq N\}$ denote the set of relative ranks of the X-sample; that is, R_{Ni} is the number of observations among X_1, \dots, X_m which do not exceed the *i*th order statistic of the combined sample $X_1, \dots, X_m, Y_1, \dots, Y_n$. In terms of previous notation one sees that $R_{Ni} = mF_mH_N^{-1}(i/N)$. It is thus natural to propose the 2-sample empirical process $\{L_N(t): 0 \leq t \leq 1\}$ defined by

(1.5)
$$L_N(t) = N^{\frac{1}{2}} [F_m H_N^{-1}(t) - FH^{-1}(t)].$$

We will now show how the L_N -process relates to the linear rank statistics studied by Chernoff and Savage (1958). Define

$$(1.6) T_N = m^{-1} \sum_{i=1}^{N} c_{Ni} R_{Ni} = \sum_{i=1}^{N} c_{Ni} F_m H_N^{-1}(i/N)$$

where $\{c_{Ni}: 1 \leq i \leq N\}$ is a given set of constants. If ν_N denotes the signed measure which puts measure c_{Ni} on the point i/N for $1 \leq i \leq N$, and puts zero measure elsewhere, then (1.6) may be written as

$$(1.7) T_N = \int_0^1 F_m H_N^{-1} d\nu_N.$$

The class of linear rank statistics described by (1.6) or (1.7) is the same as the one considered by Chernoff and Savage (1958). This can be seen by writing $Z_{Ni} = R_{Ni} - R_{Ni-1}$ for $1 < i \le N$ and $Z_{N1} = R_{N1}$. A summation by parts in (1.6) then yields

$$(1.8) T_N = m^{-1} \sum_{i=1}^N c_{Ni}^* Z_{Ni}$$

where $c_{Ni}^* = c_{Ni} + c_{Ni+1} + \cdots + c_{NN}$. Although the difference between (1.6) and (1.8) is only due to a summation by parts, the relative simplicity of the proofs of Theorems 4.1 and 5.1 given below indicates the usefulness and naturalness of the representation (1.6).

Write $\mu_N = \int_0^1 K \, d\nu_N$ and set

$$(1.9) T_N^* = N^{\frac{1}{2}} (T_N - \mu_N) = \int_0^1 L_N \, d\nu_N.$$

Let ν be another signed Lebesgue-Stieltjes measure on (0,1) for which $|\nu|([\epsilon, 1 - \epsilon]) < \infty$ for all $\epsilon > 0$ where $|\nu| = \nu^+ + \nu^-$ is the total variation of ν . In the statement of Theorems 4.1 and 5.1 conditions are given to insure that $\{\nu_N\}$ converges to ν in such a way as to permit one to substitute ν for ν_N in (1,9).

In Section 2, various results about 1-sample empirical processes are described. In Section 3, the basic identity relating the L_N -process to the 1-sample empirical processes of the two samples is presented. The results of Sections 2 and 3 are used in Sections 4 and 5 to establish in a relatively simple way the weak convergence of the L_N -process relative to various metrics. These results are then used to obtain central limit theorems for linear rank statistics. The extension of all results

to the c-sample problem is given in Section 6. Related work is contained in Pyke and Shorack ((1967a) and (1967b)). In the former the analogous limit theorems for random sample sizes are obtained, whereas in the latter questions of asymptotic relative efficiency are considered.

2. The 1-sample empirical process. For $m \geq 0$ let $\{W_m(t): 0 \leq t \leq 1\}$ denote stochastic processes on a probability space $(\Omega, \mathfrak{A}, P)$ whose sample functions are points in some metric space (\mathfrak{M}, δ) .

DEFINITION 2.1. We write $W_m \to_L W_0$ relative to (\mathfrak{M}, δ) if $\lim_{m\to\infty} E[\psi(W_m)] = E[\psi(W_0)]$ for all bounded real functionals ψ defined on \mathfrak{M} which are continuous in the δ -metric and are such that $\psi(W_m)$, $m \geq 0$, are measurable with respect to \mathfrak{M} . Such convergence is called *convergence in law* or *weak convergence*. (If the W_m process and W_0 -process are measurable with respect to the Borel sets of (\mathfrak{M}, δ) , so that their image laws on \mathfrak{M} are well defined, then the above definition is equivalent to the usual definition of weak convergence as given in Prokhorov (1956) for example.)

Suppose $\mathfrak{M}=D$, the set of all right continuous real valued functions on [0, 1] having only jump discontinuities. In this case two possible metrics are $\delta=\rho$, the uniform metric defined by

(2.1)
$$\rho(f, g) = \sup_{0 \le t \le 1} |f(t) - g(t)|,$$

and $\delta = d$, the Prokhorov metric on D as defined by Prokhorov (1956). Prokhorov showed that (D,d) is a complete separable metric space and that $U_m \to_L U_0$ relative to (D,d). Actually, since all jumps of the U_m -process equal $m^{-\frac{1}{2}}$, it is possible to show that $U_m \to_L U_0$ relative to the stronger uniform topology of the non-separable metric space (D,ρ) . We will obtain this result as Lemma 2.1. (It should be pointed out that Dudley (1966) gives a definition of weak convergence for non-separable spaces which is more general than Definition 2.1 above in that his use of upper and lower integrals enables him to place a less restrictive assumption of measurability upon the function ψ . Also, the statement " $U_m \to_L U_0$ on (D,ρ) in the sense of Prokhorov's definition" is false; see Chibisov (1965) for a statement of the measurability difficulties.)

Since $U_m \to_L U_0$ relative to (D, d) and (D, d) is a complete separable metric space it is possible, (see item 3.1.1 of Skorokhod (1956)), to construct processes $\{\tilde{U}_m(t): 0 \leq t \leq 1\}, m \geq 0$, with sample functions in D and having the same finite dimensional df's as $\{U_m(t): 0 \leq t \leq 1\}, m \geq 0$, but which in addition satisfy $d(\tilde{U}_m, \tilde{U}_0) \to_{a.s.} 0$. Let us make an independent construction for the V_n -processes so that

$$(2.2) d(\tilde{U}_m, \tilde{U}_0) \rightarrow_{\text{a.s.}} 0, d(\tilde{V}_n, \tilde{V}_0) \rightarrow_{\text{a.s.}} 0$$

where all processes are defined on a single probability space $(\tilde{\Omega}, \tilde{\mathfrak{A}}, \tilde{P})$. This is the probability space we shall work on in what follows. Note that if we set $\tilde{F}_m = m^{\frac{1}{2}} \tilde{U}_m(F) + F$, then \tilde{F}_m is a.s. a df having exactly m discontinuities each of magnitude m^{-1} . (We shall henceforth drop the symbol \sim from the notation.)

Based on the above construction, we shall prove a series of lemmas about the U_m -processes.

Warning. The results obtained below which involve convergence stronger than convergence in law may apply only to the specially constructed processes. Only the implied convergence in law should be assumed to hold for the original processes unless further checking is done.

LEMMA 2.1. In view of (2.2), $\rho(U_m, U_0) \rightarrow_{a.s.} 0$ as $m \rightarrow \infty$.

Proof. Since U_0 is continuous with probability one, it suffices to invoke Theorem 1 of Appendix 1 of Prokhorov (1956) which states that d-convergence to a continuous function is equivalent to ρ -convergence. \square

For many applications, such as those discussed in Sections 4 and 5 below, the weak convergence of the empirical processes in (D, ρ) is not sufficient. Rather, it is necessary to make use of the specific behavior of these processes near 0 and 1. For the U_0 -process, this behavior is best described by the law of the iterated logarithm. Set

$$q_0(t) = \{2t(1-t)\log[-\log t(1-t)]\}^{\frac{1}{2}}.$$

Then by the law of the iterated logarithm (cf. Lévy (1948)) it follows that

(2.3)
$$P[\limsup_{t\to a} |U_0(t)/q_0(t)| = 1] = 1, \qquad a = 0, 1.$$

To see this, let $\{W(s): s \ge 0\}$ be Brownian motion with mean zero, $E[W(s)]^2 = s$ and continuous paths a.s. By a transformation of Doob (1949), the finite dimensional df's of $\{U_0(t): 0 \le t \le 1\}$ agree with those of $\{(1-t)W(t/(1-t)) \ 0 \le t \le 1\}$. By applying the usual law of the iterated logarithm to the W-process, (both at s = 0 and at $s = +\infty$), and using this transformation, one obtains (2.3).

To study the behavior near 0 and 1 of the U_m -processes we use the following result. This result is equivalent to Lemma 7 of Govindarajulu et al. (1967), but the proof given here is simpler. For a similar result, see the result of Chibisov (1964) quoted as (4.5) in Chibisov (1965).

Lemma 2.2. Let q be any non-negative function which is non-decreasing on $[0, \theta]$ for $\theta < 1$. Then there exists a constant $c_{\theta} > 0$ such that

(2.4)
$$P[|U_m(t)| \le q(t), 0 \le t \le \theta] \ge 1 - c_\theta \int_0^\theta [q(t)]^{-2} dt$$

for all $m \geq 0$. Moreover, c_{θ} is non-increasing as $\theta \rightarrow 0$.

PROOF. For $m \ge 1$, let $\{\nu_m(t): 0 \le t \le \infty\}$ be a separable Poisson process with $E[\nu_m(1)] = m$. Then the left hand side of (2.4) equals $P(A \mid \nu_m(1) = m)$ where $A = [|\nu_m(t) - mt| \le m^{\frac{1}{2}}q(t), 0 \le t \le \theta]$.

Since $\{\nu_m(t) - mt: t \ge 0\}$ is a martingale, $\{|\nu_m(t) - mt|^2: t \ge 0\}$ is a submartingale with $E |\nu_m(t) - mt|^2 = mt$. Therefore, by Theorem 5.1 of Birnbaum and Marshall (1961)

(2.5)
$$P(A^{c}) \leq \int_{0}^{\theta} [q(t)]^{-2} dt.$$

(Since the ν_m -process has independent increments, (2.5) may also be obtained

directly from the Hájek-Rényi version of Kolmogorov's inequality; cf. Gnedenko (1962).) Furthermore, since a Poisson process has stationary and independent increments, we may write

$$P(A^{c} | \nu_{m}(1) = m)$$

$$(2.6) = \sum_{k=0}^{m} P(A^{c} \cap [\nu_{m}(\theta) = k] \cap [\nu_{m}(1) - \nu_{m}(\theta) = m - k])$$

$$\cdot (P[\nu_{m}(1) = m])^{-1}$$

$$= \sum_{k=0}^{m} P(A^{c} \cap [\nu_{m}(\theta) = k]) P[\nu_{m}(1 - \theta) = m - k] (P[\nu_{m}(1) = m])^{-1}.$$

If Y is a Poisson rv with parameter λ , then $P[Y = k] \leq P[Y = [\lambda]]$ for every k where $[\lambda]$ is the greatest integer in λ . By applying this to $P[\nu_m(1-\theta)=m-k]$, it then follows by Stirling's approximation that there exists a constant $c_{\theta} = O((1-\theta)^{-\frac{1}{2}})$ such that $P[\nu_m(1-\theta)=m-k]/P[\nu_m(1)=m] \leq c_{\theta}$ for all k. Using this bound, (2.6) gives $P(A^c | \nu_m(1)=m) \leq c_{\theta}P(A^c)$ which in conjunction with (2.5) completes the proof for $m \geq 1$.

For m=0, simply apply the Birnbaum-Marshall inequality directly to the transformed W-process described following (2.3). Thus

$$P[|U_0(t)| \le q(t), 0 \le t \le \theta]$$

$$= P[(1-t)|W(t/(1-t))| \le q(t), 0 \le t \le \theta]$$

$$= P[|W(s)| \le (1+s)q(s/(1+s)), 0 \le s \le \theta^* = \theta/(1-\theta)]$$

$$\ge 1 - \int_0^{\theta^*} (1+s)^{-2} [q(s/(1+s))]^{-2} ds = 1 - \int_0^{\theta} [q(t)]^{-2} dt$$

by the Birnbaum-Marshall inequality since $E[W(s)]^2 = s$. (This result does not use (2.2).)

From Lemma 2.2 we see the possibility of obtaining weak convergence of the U_m -process relative to more general metrics on D. Whenever well defined set

(2.7)
$$\rho_q(f,g) = \rho(f/q,g/q)$$
 and $d_q(f,g) = d(f/q,g/q)$.

DEFINITION 2.2. Let \mathbf{Q}' denote the class of all non-negative functions defined on [0, 1] which for some $\epsilon > 0$ are bounded away from zero on $(\epsilon, 1 - \epsilon)$, are non-decreasing (non-increasing) on $[0, \epsilon]([1 - \epsilon, 1])$ and which have square-integrable reciprocals. Let $\mathbf{Q} = \{q \ \varepsilon \ D : q \ge q' \text{ for some } q' \ \varepsilon \ \mathbf{Q}'\}$.

THEOREM 2.1. For $q \in \mathbb{Q}$, $\rho_q(U_m, U_0) \to_p 0$ and $d_q(U_m, U_0) \to_p 0$ as $m \to \infty$. Proof. For any $0 < a \le \epsilon \le 1 - \epsilon \le b < 1$,

$$(2.8) \quad \rho_q(U_m, U_0) \leq \sup_{0 \leq t \leq a, b \leq t \leq 1} (|U_m(t)| + |U_0(t)|)/q(t) + \rho(U_m, U_0)/\inf_{a \leq t \leq b} q(t),$$

where ϵ is a value associated with q as in Definition 2.2. For each such a, b, the second term on the right converges to zero in probability by Lemma 2.1. Fix $\delta > 0$. Since q^{-2} is integrable, choose $a < \epsilon$ so that $\int_0^a [q(t)]^{-2} dt < c_a^{-1} \delta^3$. Then by Lemma 2.2

$$(2.9) \quad P[\sup_{0 \le t \le a} |U_m(t)/q(t)| \le \delta] \ge 1 - c_a \delta^{-2} \int_0^a [q(t)]^{-2} dt > 1 - \delta.$$

Since the reversed process defined by $U_m^-(t) = U_m((1-t)-)$ has the same finite dimensional df's as the U_m -process, one may in a similar manner choose $b \ge 1 - \epsilon$ for which $P[\sup_{b \le t \le 1} |U_m(t)/q(t)| \le \delta] > 1 - \delta$. Since this inequality and (2.9) also hold for m = 0, the first term on the right hand side of (2.7) does not exceed 4δ on an event whose probability exceeds $1 - 4\delta$. This completes the proof for ρ_q .

It is straightforward to show that there is a constant c for which $d \le c\rho$; Lemma 4.3 of Chibisov (1965) states that c=4 works. Thus, the result for ρ implies the result for d. (For continuous q we could have used Theorem 1 of Appendix 1 of Prokhorov (1956) which states that ρ -convergence to a continuous function is equivalent to d-convergence.)

This theorem implies that for $q \in \mathbb{Q}$, $U_m \to_L U_0$ relative to (D, ρ_q) and relative to (D, d_q) . Related results are given by Chibisov (1964), (1965) who proves the weak convergence of the empirical processes relative to d_q when q is a function of degree $\frac{1}{2}$ in neighborhoods of 0 and 1 and satisfies

$$\int_0^1 \exp\left\{-cq^2(u)/u(1-u)\right\} [u(1-u)]^{-1} du < \infty \text{ for all } c > 0.$$

(A function h is of degree r if h(cu)/h(u) converges to c^r as u approaches a specified limit.) In view of this weak convergence and the completeness and separability of (D, d_q) , one may apply Skorokhod's result, as was done prior to (2.2), to obtain versions of the empirical processes for which $d_q(U_m, U_0) \rightarrow_{a.s.} 0$ and $\rho_q(U_m, U_0) \rightarrow_{a.s.} 0$. It can be shown that the remaining results of this section are valid for these q functions of Chibisov and that therefore the results of the following sections remain valid for this larger class of q functions

The basic identity to be derived as Lemma 3.1 will express L_N as a random linear combination of $U_m(FH_N^{-1})$ and $V_n(GH_N^{-1})$. We therefore need an analogue of Theorem 2.1 for the process $U_m(FH_N^{-1}) = U_m(K_N)$; the analogue will appear as Theorem 2.2.

Lemma 2.3. As $N \to \infty$, $\rho(HH_N^{-1}, HH^{-1}) \to_{a.s.} 0$, uniformly in all continuous F and G and all $\lambda_N \in [0, 1]$.

Proof. Since $H_N-H=\lambda_N(F_m-F)+(1-\lambda_N)(G_n-G)$, the Glivenko-Cantelli theorem implies that $\rho(H_N\,,\,H)\to_{\mathrm{a.s.}}0$ uniformly in F and G. Also $\rho(HH_N^{-1},\,HH^{-1})\ \le\ \rho(HH_N^{-1},\,HH^{-1})\ +\ \rho(H_NH_N^{-1},\,HH^{-1})$ where $\rho(H_NH_N^{-1},\,HH^{-1})\ \le\ 1/N$. (This result does not use (2.2).)

LEMMA 2.4. If $q \in \mathbb{Q}$, then as $N \to \infty$,

$$\rho(U_m(K_N)/q(K_N), U_0(K)/q(K)) \rightarrow_p 0$$

uniformly in all continuous F and G and all $\lambda_N \in [0, 1]$.

Proof. Since $H = \lambda_N F + (1 - \lambda_N)G$, Lemma 2.3 implies that $\rho(K_N, K) \to_{\text{a.s.}} 0$ and $\rho(GH_N^{-1}, GH^{-1}) \to_{\text{a.s.}} 0$ uniformly; note that $K_N - K$ and $GH_N^{-1} - GH^{-1}$ are always of the same sign. Therefore, since U_0 is continuous a.s., $\rho(U_0(K_N)/q(K_N), U_0(K)/q(K)) \to_p 0$, the argument is like that of Theorem 2.1 where a.s. continuity of U_0 is needed to handle the term from the middle of the interval. Thus this lemma follows from Theorem 2.1 and the triangle inequality. \square

Lemma 2.5². For $\epsilon > 0$, there exists b > 0 such that

$$(2.10) P[K_N(t) \le 2b\lambda_*^{-1}t for all t \ge 1/N] \ge 1 - \epsilon.$$

PROOF. It is well known (cf. [9], Lemma 8) that for every $\epsilon > 0$, there exists b > 0 such that $P(A_m) > 1 - \epsilon$ where

$$(2.11) A_m = [F(t) \le bF_m(t) \text{ for all } t \text{ where } F_m(t) > 0].$$

Therefore, it follows that if $t \ge 1/N$ and $0 < F_m H_N^{-1}(t)$, then on A_m

$$K_N(t) = FH_N^{-1}(t) \le bF_mH_N^{-1}(t) \le b\lambda_*^{-1}H_NH_N^{-1}(t) \le 2b\lambda_*^{-1}t.$$

On the other hand if $t \ge 1/N$ but $0 = F_m H_N^{-1}(t) < 1/m$, then on A_m

$$K_N(t) = FH_N^{-1}(t) \le FF_m^{-1}(1/m) \le bm^{-1} \le 2b\lambda_*^{-1}t.$$

Thus on A_m we have $K_N(t) \leq 2b\lambda_*^{-1}t$ for all $t \geq 1/N$.

Let $U_m^*(K_N)$ equal $U_m(K_N)$ for $1/N \le t \le 1 - 1/N$ and equal 0 otherwise. Theorem 2.2. For $q \in \mathbb{Q}$

$$\rho_q(U_m^*(K_N), U_0(K)) \rightarrow_p 0$$

uniformly in all continuous F and G and all $\lambda_N \in \Lambda$.

Proof. As in the proof of Theorem 2.1 the problem reduces to a study of the supremum over intervals $[0, \delta]$ and $[1 - \delta, 1]$ for sufficiently small δ . For over the interior interval $(\delta, 1 - \delta)$, q(K) and, with high probability, $q(K_N)$ are bounded away from zero so that the supremum over $(\delta, 1 - \delta)$ converges in probability to zero by Lemma 2.4. Furthermore, it suffices to consider only the interval $[0, \delta]$ since the interval $[1 - \delta, 1]$ may be treated similarly by considering the reversed process $U_m^-(t) = U_m((1 - t) - t)$. Also without loss of generality assume q is non-decreasing.

For given ϵ , $\eta > 0$, choose b by Lemma 2.5 to satisfy (2.10). Then use Lemma 2.2 to choose $\alpha > 0$ to satisfy

$$P[|U_m(t)| \leq \eta q(\lambda_* t/2b), \quad 0 < t < \alpha] > 1 - \epsilon;$$

and then choose $\delta > 0$ so that

$$P[K_N(t) < \alpha, 0 < t \le \delta] > 1 - \epsilon$$

for all sufficiently large N. This choice of δ is possible by Lemma 2.3 since $K_N(t) \leq \lambda_*^{-1} H H_N^{-1}(\delta)$ for $t \leq \delta$. With these choices of b, α and δ one obtains $P[\sup_{0 \leq t \leq \delta} |U_m^*(K_N(t))|/q(t) \leq \eta]$

$$\geq P([|U_{\scriptscriptstyle m}(t)| \leq \eta q(\lambda_* t/2b), \quad 0 < t < \alpha] \cap [K_{\scriptscriptstyle N}(t) < \alpha, 0 < t \leq \delta]$$

$$\cap [K_{\scriptscriptstyle N}(t) \leq 2b\lambda_*^{-1}(t), t \geq 1/N])$$

for N sufficiently large. Since η and ϵ are arbitrary the proof is complete. \square

$$\int_0^1 [t(1-t)]^{-\frac{1}{2}} [q(t)]^{-1} dt < \infty.$$

² The authors are indebted to the referee for Lemma 2.5. The linear bounds on K_N provided by this result allow stronger statements to be made about the behavior of empirical processes in the tails than can be obtained by iterating Lemma 2.2. The main results of this paper are now able to be stated for $q \in Q$ rather than for the slightly more restrictive class of functions satisfying

3. The 2-sample empirical process. Recall that by Lemma 2.1

$$(3.1) \rho(U_m, U_0) \rightarrow_{\text{a.s.}} 0, \rho(V_n, V_0) \rightarrow_{\text{a.s.}} 0 \text{ as } m, n \rightarrow \infty.$$

Our study and application of the 2-sample empirical process $\{L_N(t): 0 \le t \le 1\}$ defined in (1.5) depends on the following identity which relates it to the 1-sample empirical U_{m^-} and V_{n^-} -processes.

LEMMA 3.1. With probability 1,

(3.2)
$$L_N(t) = (1 - \lambda_N) \{ \lambda_N^{-\frac{1}{2}} B_N(t) U_m(FH_N^{-1}(t)) - (1 - \lambda_N)^{-\frac{1}{2}} A_N(t) V_n(GH_N^{-1}(t)) + \delta_N(t) \}$$

for all $t \varepsilon (0, 1]$ where

(3.3)
$$\delta_N(t) = A_N(t) N^{\frac{1}{2}} [H_N H_N^{-1}(t) - t],$$

$$(3.4) A_N(t) = [K(u_t) - K(t)]/(u_t - t), u_t = HH_N^{-1}(t),$$

and B_N is defined by

(3.5)
$$\lambda_N A_N(t) + (1 - \lambda_N) B_N(t) = 1.$$

(In (3.2) L_N is defined by left continuity at any otherwise undefined points.) PROOF. The relationship (3.2) follows immediately from (1.3), (1.4) and (1.5) by writing

$$L_{N} = N^{\frac{1}{2}} [F_{m} H_{N}^{-1} - FH^{-1}] = N^{\frac{1}{2}} [m^{-\frac{1}{2}} U_{m} (FH_{N}^{-1}) + FH_{N}^{-1} - FH^{-1}]$$

$$= N^{\frac{1}{2}} [m^{-\frac{1}{2}} U_{m} (FH_{N}^{-1}) + A_{N} (HH_{N}^{-1} - HH^{-1})]$$

provided $HH_N^{-1} \neq HH^{-1}$. Since with probability one this proviso is satisfied except at a finite number of points, and since

$$\begin{split} N^{\frac{1}{2}}[HH_{N}^{-1} - HH^{-1}] \\ &= N^{\frac{1}{2}}[HH_{N}^{-1} - H_{N}H_{N}^{-1} + H_{N}H_{N}^{-1} - HH^{-1}] \\ &= -\lambda_{N}^{\frac{1}{2}}U_{m}(FH_{N}^{-1}) - (1 - \lambda_{N})^{\frac{1}{2}}V_{n}(GH_{N}^{-1}) + N^{\frac{1}{2}}[H_{N}H_{N}^{-1} - HH^{-1}] \end{split}$$

the desired result follows. If $HH_N^{-1} = HH^{-1}$, it follows from the first line of the proof that the proper expression for L_N is $L_N = \lambda_N^{-\frac{1}{2}} U_m(FH^{-1})$.

In the representation (3.2), one has for all $t \in [0, 1]$ that

$$(3.6) |A_N(t)| \le \lambda_N^{-1}, |B_N(t)| \le (1 - \lambda_N)^{-1}, |\delta_N(t)| \le \lambda_N^{-1} N^{-\frac{1}{2}}.$$

The first two inequalities follow from (3.5) since $A_N(t)$ and $B_N(t)$ are of the same sign, whereas the third inequality derives from (3.3) and the fact that $0 \le H_N H_N^{-1}(t) - t \le N^{-1}$.

From the representation (3.2) of the L_N -process, it is easy to see what must be the natural limiting process. We define first of all the \bar{L}_N -process by

$$(3.7) \quad \bar{L}_{N}(t) = (1 - \lambda_{N}) \{ \lambda_{N}^{-\frac{1}{2}} B_{N}(t) U_{0}(FH^{-1}(t)) - (1 - \lambda_{N})^{-\frac{1}{2}} A_{N}(t) V_{0}(GH^{-1}(t)) \}$$

for $0 \le t \le 1$. In view of Lemma 2.4 (for $q \equiv 1$) and its analogue for the V_n -process, the following limiting result is immediate.

LEMMA 3.2. As
$$N \to \infty$$
, $\rho(L_N, \bar{L}_N) \to_p 0$.

To obtain the natural limiting process from (3.7) it remains to study the limiting behavior of B_N and A_N . By their definition, A_N and B_N are difference quotients of FH^{-1} and GH^{-1} respectively. In view of (1.4), $\lambda_N FH^{-1}(t)+(1-\lambda_N)GH^{-1}(t)=t$, so that both FH^{-1} and GH^{-1} are absolutely continuous. Let a_N and b_N denote the derivatives of FH^{-1} and GH^{-1} respectively, which therefore exist a.e. on [0,1] with respect to Lebesgue measure. In some results, we will let λ_N converge to a fixed λ_0 . Let a_0 , b_0 denote the derivatives of FH_0^{-1} , GH_0^{-1} where $H_0 = \lambda_0 F + (1-\lambda_0)G$. Now, wherever defined, set

$$(3.8) \quad L_0(t) = (1 - \lambda_0)^{\frac{1}{2}} \{ \lambda_0^{-\frac{1}{2}} b_0(t) U_0(FH_0^{-1}(t)) - (1 - \lambda_0)^{-\frac{1}{2}} a_0(t) V_0(GH_0^{-1}(t)) \}$$
 and

$$(3.9) \quad L_{0N}(t) = (1 - \lambda_N)^{\frac{1}{2}} \{\lambda_N^{-\frac{1}{2}} b_N(t) U_0(FH^{-1}(t)) - (1 - \lambda_N)^{-\frac{1}{2}} a_N(t) V_0(GH^{-1}(t)) \}.$$

These expressions respectively determine the natural limiting processes for L_N when λ_N does or does not converge. In the following sections the convergence of L_N to these limits will be considered.

4. Convergence in the ρ_q and d_q metrics. We now consider the convergence of the L_N -process to the limiting L_0 -process when $\lambda_N \to \lambda_0$. Let $q \in \mathbb{Q}$, let δ_N^* be defined to equal δ_N on [1/N, 1 - 1/N] and 0 elsewhere, and let L_N' equal L_N on [1/N, 1] and 0 elsewhere. By (3.6), $\rho_q(\delta_N^*, 0) = o(1)$. Also

$$\sup_{1-1/N \le t < 1} |L_N(t)/q(t)| = o(1)$$

since $|L_N(t)| = N^{\frac{1}{2}}|1 - FH^{-1}(t)| \le N^{\frac{1}{2}} \lambda_*^{-1}(1-t)$ in this interval. Thus, in order to prove that

$$\rho_q(L_N', L_0) \to_p 0$$

it will, by (3.2) and the triangle inequality, suffice to prove that

$$\rho(B_N, 0)[\rho_q(U_m^*(K_N), U_0(K)) + \rho_q(U_0(K), U_0(K_0))]$$

$$+ \rho(B_N, b_0)\rho_{\sigma}(U_0(K_0), 0),$$

and the analogous quantity in A_N , V_{n_q} and G, converge in probability to 0. Apply (3.6) to $\rho(B_N, 0)$, Theorem 2.2 to $\rho_q(U_m^*(K_N), U_0(K))$, the proof of Lemma 2.4 to $\rho_q(U_0(K), U_0(K_0))$ and Lemma 2.2 to $\rho_q(U_0(K_0), 0)$ to obtain that (4.1) holds provided

$$\rho(B_N, b_0) \to_p 0.$$

We will now make a differentiability assumption on the functions $\{K_{\lambda} : \lambda \in \Lambda\}$ that will imply (4.2). This assumption is satisfied in many important cases, as Corollary 4.1 indicates.

Assumption 4.1. The functions K_{λ} have derivatives a_{λ} for all $t \in (0, 1)$, and for some λ' , $a_{\lambda'}$ is continuous on (0, 1) and has one-sided limits at 0 and 1.

It is an easy consequence of (1.4) and Assumption 4.1 that the functions GH_{λ}^{-1} have derivatives b_{λ} satisfying $\lambda a_{\lambda} + (1 - \lambda)b_{\lambda} = 1$ and that $b_{\lambda'}$ is continuous on (0, 1) and has one-sided limits at 0 and 1.

Definition 4.1. A family of functions $\{h_{\lambda} : \lambda \in \Lambda\}$ is said to be uniformly equicontinuous (uec) if for all $\epsilon > 0$ there exists $\delta_{\epsilon} > 0$ such that $|h_{\lambda}(u) - h_{\lambda}(v)| < \epsilon$ for all $\lambda \in \Lambda$ and for all u, v in the domain of the functions with $|u - v| < \delta_{\epsilon}$.

LEMMA 4.1. Under Assumption 4.1 $\{a_{\lambda} : \lambda \in \Lambda\}$ and $\{b_{\lambda} : \lambda \in \Lambda\}$ are two uec families of functions on [0, 1] and $a_{\lambda}(b_{\lambda})$ converges uniformly to $a_0(b_0)$ as $\lambda \to \lambda_0$. Proof. Differentiate $K_{\lambda} = K_{\lambda'}(H_{\lambda'}H_{\lambda}^{-1})$ to obtain

$$(4.3) a_{\lambda} = a_{\lambda'}(H_{\lambda'}H_{\lambda}^{-1})[(\lambda' - \lambda)a_{\lambda} + (1 - \lambda')](1 - \lambda)^{-1}.$$

If we set $\alpha_{\lambda} = \alpha_{\lambda'}(H_{\lambda'}H_{\lambda}^{-1})$ and let ϕ_{λ} denote the remaining factors in (4.3) so that $a_{\lambda} = \alpha_{\lambda} \phi_{\lambda}$, then

$$[a_{\lambda}(u) - a_{\lambda}(v)][1 - \alpha_{\lambda}(v)(\lambda' - \lambda)(1 - \lambda)^{-1}] = [\alpha_{\lambda}(u) - \alpha_{\lambda}(v)]\phi_{\lambda}(u).$$

But for $\lambda \varepsilon \Lambda$, then ϕ_{λ} 's are uniformly bounded since $|a_{\lambda}| \leq 1$ and $\alpha_{\lambda}(\lambda' - \lambda)(1 - \lambda)^{-1}$ is uniformly bounded away from 1. Thus for $\lambda \varepsilon \Lambda$, $|a_{\lambda}(u) - a_{\lambda}(v)| \leq C |\alpha_{\lambda}(u) - \alpha_{\lambda}(v)|$ for some constant C. To complete the proof that $\{a_{\lambda}: \lambda \in \Lambda\}$ is uec, apply to the definition of a_{λ} the fact that $\{H_{\lambda'}H_{\lambda}^{-1}: \lambda \in \Lambda\}$ is uec; since for $u \le v$ we have $0 \le FH_{\lambda}^{-1}(v) - FH_{\lambda}^{-1}(u) \le \lambda^{-1}(v - u) \le \lambda^{-1}(v - u)$ $\lambda_*^{-1}(v-u)$.

Lemma 4.2. Under Assumption 4.1 $\rho(A_N, a_0) \rightarrow_{a.s.} 0$ and $\rho(B_N, b_0) \rightarrow_{a.s.} 0$ as $N o \infty$.

Proof. By the definition of A_N in (3.4), the mean value theorem and Lemma 4.1 we have that a_N is continuous and $A_N(t) = a_N(t_N)$ for some t_N between t and $HH_N^{-1}(t)$. Consequently, $\rho(A_N, a_N) = \sup_{0 \le t \le 1} |a_N(t_N) - a_N(t)|$. Since, by Lemma 2.3, $t_N \to t$ uniformly on [0, 1] as $N \to \infty$, it follows from the uec of the a_N 's that $\rho(A_N, a_N) \to_{a.s.} 0$. Also, since $\lambda_N \to \lambda_0$, $\rho(a_N, a_0) \to 0$ by Lemma 4.1. An application of the triangle inequality completes the proof.

Due to the left continuity of the L_{N} -process we introduce D^{-} for the set of all left continuous functions on [0, 1].

Theorem 4.1. (a) Suppose Assumption 4.1 holds, $\lambda_N \to \lambda_0$ and $q \in \mathbb{Q}$. Then $\rho_q(L_N', L_0) \to_p 0$ so that $L_N' \to_L L_0$ relative to (D^-, ρ_q) . The same statement holds

- (b) If in addition the measures $\{\nu_N : N' \geq 1\}$ and ν of Section 1 satisfy
- (i) $\int_{1/N}^{1} L_N d(\nu_N \nu) \rightarrow_p 0$, and, (ii) $\int_{0}^{1} q d|\nu| < \infty$,

then $T_N^* \xrightarrow{}_p \int_0^1 L_0 d\nu$, a $N(0, \sigma_0^2)$ rv where

$$(4.4) \quad \sigma_0^2 = 2(1-\lambda_0)^2 \{\lambda_0^{-1} \int_0^1 \int_0^v b_0(u) \ b_0(v) F H_0^{-1}(u) [1-F H_0^{-1}(v)] \ d\nu(u) \ d\nu(v) + (1-\lambda_0)^{-1} \int_0^1 \int_0^v a_0(u) a_0(v) G H_0^{-1}(u) [1-G H_0^{-1}(v)] \ d\nu(u) \ d\nu(v) \}.$$

Proof. (a) Lemma 4.2 shows that (4.1) holds under Assumption 4.1. Con-

vergence in law follows easily by definition. Results for d_q follow by an analogous argument, since $d_q \leq 4\rho_q$ as remarked in the proof of Theorem 2.1.

(b) Now

$$|T_N^* - \int_0^1 L_0 d\nu| = |\int_0^1 L_N' d(\nu_N - \nu) + \int_0^1 (L_N' - L_0) d\nu|$$

$$\leq |\int_0^1 L_N' d(\nu_N - \nu)| + \rho_q(L_N', L_0) \int_0^1 q d|\nu| \to_{\nu} 0.$$

The fact that the limiting rv has a normal distribution relies on the existence and finiteness of its variance, which in turn follows from (ii) since

$$E\left[\int_{0}^{1} L_{0} d\nu\right]^{2} \leq c \int_{0}^{1} \int_{0}^{\nu} u(1-\nu) d|\nu|(u) d|\nu|(v)$$

$$\leq c \int_{0}^{1} \int_{0}^{\nu} \left\{ \left[u(1-u)v(1-\nu)\right]^{\frac{1}{2}} \left[q(u)q(v)\right]^{-1} \right\} q(u)q(v) d|\nu|(u) d|\nu|(v)$$

for some constant c.

Since the L_N -process takes jumps only at points i/N and since the L_0 -process has continuous paths a.s., the mappings L_N ($N \ge 0$) from (Ω, \mathfrak{A}) to (D^-, ρ_q) are measurable; see Billingsley (1967), p. 308. Thus the weak convergence in Theorem 4.1 is valid in the sense of Prokhorov's usual definition.

It is necessary to write L_N' and not L_N in Theorem 4.1 since if the smallest observation is an X, then $L_N(0+) = N^{\frac{1}{2}}m^{-1}$.

REMARK. If λ_N does not converge, then in Theorem 4.1 we may write $\rho_q(L_N', L_{0N}) \to_p 0$ and $T_N^* - \int_0^1 L_{0N} d\nu \to_p 0$ where $\int_0^1 L_{0N} d\nu$ is $N(0, \sigma_{0N}^2)$ with σ_{0N}^2 given by the obvious modification of (4.4). In this form, the theorem relates most closely to the main theorem of Chernoff and Savage (1958); σ_{0N}^2 is given as equation (4.3) in this reference.

COROLLARY 4.1. The conclusion of Theorem 4.1 is valid if in place of Assumption 4.1 we have either (i) F = G or (ii) there are numbers $-\infty \le \alpha_1$, α_2 , γ_1 , $\gamma_2 \le +\infty$ for which $F(\alpha_1) = G(\alpha_1) = 0$, $F(\alpha_2) = G(\alpha_2) = 1$, the boundary of the set of common zeros of the continuous functions f = F' and g = G' is finite and $\gamma_i = \lim_{x \to \alpha_i} [g(x)/f(x)]$ exists for i = 1, 2.

Proof. (i). This follows since Assumption 4.1 is satisfied when F = G. (ii). Suppose first that the set of common zeros is empty. Since $a_{\lambda}(t) = (d/dt)FH_{\lambda}^{-1}(t) = (\lambda + (1-\lambda)[gH_{\lambda}^{-1}(t)/fH_{\lambda}^{-1}(t)])^{-1}$, the existence of γ_1 and γ_2 implies in this case that Assumption 4.1 holds.

In the general case, consider the finite set of points t for which $H_{\lambda_0}^{-1}(t)$ is in the boundary of the set of common zeros; these points are interior to (0, 1). Hence, we may throw away intervals of arbitrarily small $|\nu|$ -measure that contain these points. On the complement of these intervals Assumption 4.1, with Λ replaced by a sufficiently small interval about λ_0 , holds. Now apply the proof of Lemma 4.2 to each interval in the complement. \square

5. Convergence in other metric spaces. A study of linear rank statistics leads naturally to a study of integral functionals on stochastic processes; see (1.9). It thus seems natural to consider integral metrics on these processes. Indeed, this will allow us to obtain results under weaker hypotheses than previously

assumed. This follows since the B_N 's may converge in an integral norm when they fail to converge in a supremum norm.

Whenever well defined let

(5.1)
$$||f||_{\nu} = \int_0^1 |f(t)| \, d \, |\nu| \, (t).$$

We shall drop the subscript and write $\|\cdot\|$ in case ν is Lebesgue measure on

LEMMA 5.1. If
$$\int_0^1 q \, d|\mathbf{p}| < \infty$$
 for some $q \in \mathbf{Q}$ then
$$||U_m^*(K_N) - U_0(K)||_{\mathbf{p}} \to_{\mathbf{p}} 0$$

uniformly in all continuous F and G and all $\lambda_N \in \Lambda$.

Proof. This is an immediate consequence of Theorem 2.2 since $||f||_{\nu} \leq$ $\rho_{q}(f, 0) \int_{0}^{1} q \, d|\nu|. \quad \square$

As in Section 4, one may show that $\|\delta_N^*\|_{\nu} = o(1)$, and $|\int_{1-1/N}^1 L_N d\nu| = o(1)$. Thus by analogy with equations (4.1) and (4.2) we have that in order to establish

(5.2)
$$||L_{N}' - L_{0}||_{r} \rightarrow_{p} 0 \qquad (||L_{N}' - L_{0N}||_{r} \rightarrow_{p} 0)$$

it suffices to find conditions under which

(5.3)
$$||(B_N - b_0)q||_{\nu} \to_{p} 0 \qquad (||(B_N - b_N)q||_{\nu} \to_{p} 0)$$

for some $q \in \mathbf{Q}$.

Let $C_0 \subset (0,1)$ be the set on which $K_0 = FH_0^{-1}$ is differentiable. By the remark at the end of Section 3, the Lebesgue measure of C_0 is 1.

Lemma 5.2. If $\lambda_N = \lambda_0 + O(N^{-\frac{1}{2}})$ and μ is a Lebesgue-Stieltjes measure on (0,1) for which $\mu(C_0^c)=0$, then $A_N\to a_0$ and $B_N\to b_0$ in $P\times \mu$ -measure as $N\to\infty$.

PROOF. By (3.5), it suffices to prove the result about A_N . From (1.3) and (3.4) one obtains that $K = FH^{-1} = FH_0^{-1}(H_0H^{-1}) = K_0(H_0H^{-1})$ and that $A_N = A_{N1}A_{N2}$ where

$$A_{N1} = [K_0(H_0H_N^{-1}) - K_0(H_0H^{-1})][H_0H_N^{-1} - H_0H^{-1}]^{-1},$$

$$A_{N2} = [H_0H_N^{-1} - H_0H^{-1}][HH_N^{-1} - HH^{-1}]^{-1}.$$

By (1.4),
$$H_0 = \lambda_0 F + (1 - \lambda_0)G = H + (\lambda_0 - \lambda_N)(F - G)$$
, so that

$$H_0H_N^{-1}(t) - H_0H^{-1}(t) = HH_N^{-1}(t) - t + (\lambda_0 - \lambda_N)D_N(t)$$

where $D_N = FH_N^{-1} - FH^{-1} - GH_N^{-1} + GH^{-1}$. Therefore,

(5.4)
$$A_{N2} = 1 + (\lambda_0 - \lambda_N) D_N (d_N - N^{-\frac{1}{2}} W_N)^{-1}$$

where $0 \le d_N(t) = H_N H_N^{-1}(t) - t \le 1/N$ and $W_N = N^{\frac{1}{2}} (H_N H_N^{-1} - H H_N^{-1})$. By definitions (1.1) and (1.2),

$$W_N = \lambda_N^{\frac{1}{2}} U_m(FH_N^{-1}) + (1 - \lambda_N)^{\frac{1}{2}} V_n(GH_N^{-1}).$$

It therefore follows by arguments similar to those used to prove Lemma 2.3 and 2.4 that $\rho(W_N, W_0) \to_p 0$ where $W_0 = \lambda_0^{\frac{1}{2}} U_0(FH_0^{-1}) + (1 - \lambda_0)^{\frac{1}{2}} V_0(GH_0^{-1})$. By Lemma 2.3, $\rho(D_N, 0) \rightarrow_{a.s.} 0$. Let $\{N_k : k \geq 1\}$ be a subsequence for which $\rho(W_{N_k}, W_0) \to_{\text{a.s.}} 0$. For each $\omega \in \Omega$ let $C_{0,\omega} = \{t \in [0, 1]: W_0(t) \neq 0\}$. By applying to (5.4) the convergences just mentioned and the hypothesis $\lambda_N - \lambda_0 = O(N^{-\frac{1}{2}})$, it follows that every subsequence has a subsequence, $\{N_k\}$ say, for which

(5.5)
$$P\{w \in \Omega: A_{N_k 2}(t) \to 1 \text{ for all } t \in C_{0,\omega}\} = 1.$$

To study the limit of A_{N1} , observe first that for any function p

$$\frac{p(x) - p(y)}{x - y} = \frac{p(x) - p(x_0)}{x - x_0} \frac{x - x_0}{x - y} + \frac{p(x_0) - p(y)}{x_0 - y} \frac{x_0 - y}{x - z},$$

Thus if p is differentiable at x_0 with derivative p', then $[p(x) - p(y)]/(x - y) \to p'(x_0)$ if $x, y \to x_0$ in such a way that $(x_0 - y)/(x - y)$ remains bounded. As $N \to \infty$ over a subsequence $\{N_k\}$ for which $\rho(W_{N_k}, W_0) \to_{a.s.} 0$, one obtains that for each $t \in C_{0,\omega}$ and almost all ω ,

$$(H_0H_0^{-1} - H_0H^{-1})(H_0H_N^{-1} - H_0H^{-1})^{-1}$$

$$= (\lambda_N - \lambda_0)(FH^{-1} - GH^{-1})[d_N - N^{-\frac{1}{2}}W_N + (\lambda_0 - \lambda_N)D_N]^{-1}$$

remains bounded. Therefore

$$(5.6) P\{\omega \varepsilon \Omega \colon A_{N_k^{-1}}(t) \to a_0(t) \text{ for all } t \varepsilon C_0 \cap C_{0,\omega}\} = 1.$$

Since $P[W_0(t) = 0] = 0$ for each $t \in (0, 1)$ and since the set of pairs (t, ω) for which $t \in C_{0,\omega}$ is measurable because of the almost sure continuity of the W_0 -process, it follows from Fubini's theorem that $\mu(C_{0,\omega}^c) = 0$ for almost all ω . The proof is completed by combining (5.5) and (5.6). \square

THEOREM 5.1. (a) Suppose K_0 is differentiable a.e. $|\overline{\nu}|$, $\lambda_N = \lambda_0 + O(N^{-\frac{1}{2}})$. and $\int_0^1 q \, d|\nu| < \infty$ for some $q \in \mathbb{Q}$. Then $||L_N' - L_0||_{\nu} \to_{\mathcal{P}} 0$ so that $L_N' \to_{\mathcal{L}} L_0$ relative to $(D^-, ||\cdot||_{\nu})$.

(b) Suppose in addition that the measures $\{\nu_N : N \geq 1\}$ of Section 1 satisfy (i) $\int_{1/N}^1 L_N d(\nu_N - \nu) \to_p 0$.

Then $T_N \to_p \int_0^1 L_0 d\nu$ which is a $N(0, \sigma_0^2)$ rv with σ_0^2 given by (4.4).

PROOF. (a) By Lemma 5.2 with $\mu = |\nu|$, every subsequence of B_N contains a further subsequence converging to b_0 a.e. $P \times |\nu|$. Thus by the dominated convergence theorem we obtain $\|(B_N - b_0)q\|_{|\nu|} \to_{\text{a.s.}} 0$. Convergence in law follows easily by definition.

(b) Now

$$(5.7) |T_N^* - \int_0^1 L_0 d\nu| \leq |\int_0^1 L_N' d(\nu_N - \nu)| + |\int_0^1 (L_N' - L_0) d\nu|$$

which equals $o_p(1)$ by (i) and since we have established (5.3) for $|\nu|$.

Note that Theorem 4.1(a) has a stronger conclusion than Theorem 5.1(a) since the ρ_q -topology on D^- contains the $\|\cdot\|_{\nu}$ -topology on D^- when $\int_0^1 q \, d \, |\nu| < \infty$. However, the functional $\int_0^1 \cdot d\nu$ is $\|\cdot\|_{\nu}$ -continuous a.s. on D^- with respect to the measure of the L_0 -process, so that it was possible to obtain part (b) of Theorem 5.1 in analogy with Theorem 4.1(b).

We next discuss a relationship between Theorem 5.1(b) and Theorem 1 of Chernoff and Savage (1958). Define a function J_N on [0, 1] by letting $J_N(t) = c_{Ni}^*$ for $(i-1)/N < t \le i/N$ for $i=1, \dots, N$ and $J_N(0) = J_N(0+)$ so that $T_N = \int_0^1 J_N dF_m H_N^{-1}$ by (1.7). Let -J denote a non-constant function of bounded variation on $(\epsilon, 1 - \epsilon)$ for all $\epsilon > 0$ which induces the Lebesgue-Stieltjes measure ν . Chernoff and Savage consider the statistic $N^{\frac{1}{2}}[T_N - \int_{-\infty}^{\infty} J(H) dF] = T_N^* + \gamma_N$ where $\gamma_N = N^{\frac{1}{2}} \int_0^1 [J_N(H) - J(H)] dF$; this equality is obtained by summing μ_N by parts.

Proposition 5.1. Suppose

- (1) $N^{\frac{1}{2}} \int_{0}^{1} |J_{N}(t) J_{N}(t-)| dH_{N}H_{N}^{-1}(t) = o(1),$ (2) $J_{N}(0) = o(N^{\frac{1}{2}})$ and $J_{N}(1) = o(N^{\frac{1}{2}}),$
- (3) $|J(t)| \leq K[t(1-t)]^{-\frac{1}{2}+\delta}$ for some constants $K, \delta > 0$ and
- (4) $N^{\frac{1}{2}} \int_{0}^{1} |J_{N}(t) J(t)| dt = o(1).$

Then

- (a) $\gamma_N = o(1)$, (b) $\int_{1/N}^1 L_N d(\nu_N \nu) = o_p(1)$ uniformly in all continuous F and G and all $\lambda_N \in \Lambda$, and
 - (c) there exists $q \in \mathbb{Q}$ such that $\int_0^1 q \, d|\nu| < \infty$.

PROOF. (a) $\gamma_N \leq \lambda_*^{-1} N^{\frac{1}{2}} \int_0^1 |J_N(H) - J(H)| dH = o(1)$ by (1). (b) $I_N \equiv \int_{1/N}^1 L_N d(\nu_N - \nu) = -\int_{1/N}^{1-1/N} L_N d(J_N - J) + o(1)$ by (3) and the bound on L_N in [1 - 1/N, 1]. Thus $I_N = -\int_{1/N}^{1-1/N} L_N (t+) d(J_N(t) - J(t))$ + o(1) using the integration by parts formula on page 419 of Hewitt and Stromberg in conjunction with (2), (3), (4) and the fact that $|L_N(t) - L_N(t+)|$ $= O(N^{\frac{1}{2}})$ for all t. A second integration by parts used in conjunction with (2), (3) and the bound on L_N at 1/N and 1 - 1/N shows that

$$I_{N} = \int_{1/N}^{1-1/N} \left[J_{N}(t) - J(t-) \right] dL_{N}(t) + o(1)$$

$$\leq \lambda_{*}^{-1} N^{\frac{1}{2}} \left[\int_{0}^{1} \left| J_{N}(t) - J(t-) \right| dH_{N} H_{N}^{-1}(t) + \int_{0}^{1} \left| J_{N}(t) - J(t) \right| dt \right]$$

$$= o(1)$$

where the last equality follows from (1) and (4).

(c) Let $q(t) = [t(1-t)]^{\frac{1}{2}-\delta/2}$ and use integration by parts. \square

Corollary 5.1. Proposition 5.1 remains true if (4) is replaced by

 $(4')J = J_d + J_c$ where J_d is a saltus function taking only a finite number of jumps and where J_c has a continuous derivative J_c' on intervals $(0, a_1), (a_1, a_2), \cdots$ $(a_s, 1)$ which satisfies

$$|J_c'(t)| \leq K[t(1-t)]^{-\frac{3}{2}+\delta}$$
 for $t \neq any a_i$.

Proof. We give the proof for s = 0. Let $t_N = (1 + [Nt])/N$ where $[\cdot]$ is the greatest integer function. Now by (2) and (3)

$$\begin{split} I_N &\equiv N^{\frac{1}{2}} \int_0^1 |J_N - J| \, dt = N^{\frac{1}{2}} \int_{2/N}^{1-2/N} |J_N - J| \, dt + o(1) \\ &\leq N^{\frac{1}{2}} \int_0^1 |J_N(t) - J(t_N)| \, dt + N^{\frac{1}{2}} \int_0^1 |J_d(t_N) - J_d(t)| \, dt \\ &\qquad \qquad + N^{\frac{1}{2}} \int_{2/N}^{1-2/N} |J_c(t_N) - J_c(t)| \, dt + o(1) \\ &= N^{\frac{1}{2}} \int_{2/N}^{1-2/N} |J_c(t_N) - J_c(t)| \, dt + o(1) \end{split}$$

where the last equality follows from (1) and the (4') condition on J_d . But

$$N^{\frac{1}{2}} \int_{\frac{1}{2}}^{1-2/N} |J_c(t_N)| - |J_c(t)| dt \le N^{\frac{1}{2}} \int_{\frac{1}{2}}^{1-1/N} N^{-1} K[t(1-t)]^{-\frac{3}{2}+\delta} dt = o(1)$$

since (4') makes possible an application of the mean value theorem to show $|J_c(t_N) - J_c(t)| \leq N^{-1}K[(t+1/N)(1-t-1/N)]^{-\frac{1}{2}+\delta}$ when $t \leq t_N \leq t+1/N$; and where the factor t+1/N accounts for the change in the limits of integration. The integral from 2/N to $\frac{1}{2}$ is similar. Thus $I_N = o(1)$, so that (4) holds. If $s \neq 0$, then separate off sufficiently small intervals about each a_i and repeat the above proof. \Box

Theorem 2 of Chernoff and Savage shows that (1), (2), (3) and (4') hold when $J = \Psi^{-1}$ where Ψ is a df satisfying

$$|(d^{(i)}/dt^{(i)})\Psi^{-1}(t)| \le K[t(1-t)]^{-i-\frac{1}{2}+\delta}$$
 for $i = 0, 1, 2$

and where c_{Ni}^* is the expectation of the *i*th smallest order statistic in a sample of size N from a population having df Ψ . This shows that Theorem 5.1(b) is broadly applicable. Statements that condition (4) in Theorem 1 of Chernoff and Savage, which puts bounds on J, J' and J'', need hold only a.e. appear frequently in the literature; though the validity of that statement is suspect (see Govindarajulu, et al. (1967)). Proposition 5.1 and its corollary are weak enough to cover all the cases where Chernoff and Savages' (4) failed at a finite number of points; and this is weak enough to cover all cases of known interest.

REMARK. Let $q \, \varepsilon \, \mathbf{Q}$ be fixed. Consider a class \mathbf{J} of sequences (J, J_1, J_2, \cdots) for which $\int_0^1 q \, d|J|$ is uniformly bounded and for which the hypotheses of Proposition 5.1 are uniform in \mathbf{J} . Then the conclusion of Theorem 5.1(b) is uniform in \mathbf{J} . To see this, note that the term $\|L_N' - L_0\|_{\nu}$ in (5.7) is $o_p(1)$ uniformly in \mathbf{J} by the uniform boundedness of $\int_0^1 q \, d|J|$ in \mathbf{J} . The term $|\int_{1/N}^1 L_N \, d(\nu_N - \nu)|$ in (5.7) is $o_p(1)$ uniformly in \mathbf{J} by the proof of Proposition 5.1.

6. The c-sample problem. Let X_{j1} , \cdots , X_{jn_j} , $j=1, \cdots, c$, be independent random samples from populations having continuous df's F_j . Let $N=n_1+\cdots+n_c$; when used as a subscript N will denote the c-tuple (n_1, \cdots, n_c) . Let $\lambda_N=(\lambda_{N1}, \cdots, \lambda_{Nc})$ where $\lambda_{Nj}=n_j/N$. Suppose $\lim_{N\to\infty}\lambda_{Nj}=\lambda_{0j}$ exists and $\lambda_*<\lambda_{jN}<1-\lambda_*$ for some $\lambda_*>0$. Let $\lambda_0=(\lambda_{01}, \cdots, \lambda_{0c})$. Let F_{Nj} denote the empirical df of the jth sample. Then $H_N=\sum_{j=1}^c\lambda_{Nj}F_{Nj}$ is the empirical df of the combined sample. Let $H_0=\sum_{j=1}^c\lambda_{0j}F_j$, $H_\lambda=\sum_{j=1}^c\lambda_{Nj}F_j$ where $\lambda=(\lambda_1, \cdots, \lambda_c)$ and $H=H_{\lambda_N}$. For $j=1, \cdots, c$ define stochastic processes $\{L_{Nj}(t): 0 \leq t \leq 1\}$ and $\{L_{0j}(t): 0 \leq t \leq 1\}$ by

(6.1)
$$L_{Nj}(t) = N^{\frac{1}{2}} [F_{Nj} H_N^{-1}(t) - F_j H^{-1}(t)]$$

and

(6.2)
$$L_{0j}(t) = -a_{0j}(t) \sum_{i \neq j} \lambda_{0i}^{\frac{1}{2}} U_{0i}(F_i H_0^{-1}(t)) + \left[\sum_{i \neq j} (\lambda_{0i} / \lambda_{0j}^{\frac{1}{2}}) a_{0i}(t) \right] U_{0j}(F_j H_0^{-1}(t))$$

where the processes $\{U_{0j}(t): 0 \le t \le 1\}, j = 1, \dots, c$, are independent tied-down Wiener processes.

We consider the c-sample analog of Theorem 5.1(b). Let $Z_{Ni}^{(j)}$ equal 1 or 0 depending on whether or not the ith smallest observation from the combined sample of size N is from the jth sample. Consider the linear rank statistics

$$(6.3) T_{Nj} = n_j^{-1} \sum_{i=1}^{N} C_{Ni} R_{Ni}^{(j)} = n_j^{-1} \sum_{i=1}^{N} C_{Ni}^* Z_{Ni}^{(j)}$$

where $R_{Ni}^{(j)}=Z_{N1}^{(j)}+\cdots+Z_{Ni}^{(j)}, \{C_{Ni}:1\leq i\leq N\}$ is a given set of constants and $C_{Ni}^* = C_{Ni} + \cdots + C_{NN}$. Let ν_N denote the signed measure which puts measure C_{Ni} on the point i/N for $1 \le i \le N$ and puts zero measure elsewhere. Then, for $j = 1, \dots, c$,

(6.4)
$$T_{Nj} = \int_0^1 F_{Nj} H_N^{-1} d\nu_N.$$

For $j = 1, \dots, c$ let $\mu_{Nj} = \int_0^1 F_j H^{-1} d\nu_N$ and set $T_{Nj}^* = N^{\frac{1}{2}} (T_{Nj} - \mu_{Nj})$. Let ν be a signed Lebesgue-Stieltjes measure on (0, 1) for which $|\nu|([\epsilon, 1 - \epsilon]) < \infty$ for all $\epsilon > 0$.

THEOREM 6.1. Suppose $\int_0^1 q \, d|\nu| < \infty$ for some $q \in \mathbb{Q}$, and for $j = 1, \dots, c$ we have $\lambda_{Nj} = \lambda_{0j} + O(N^{-\frac{1}{2}})$ and $K_{0j} = F_j H_0^{-1}$ is differentiable a.e. ν . If

(i) $N^{\frac{1}{2}} \int_{1/N}^1 L_{Nj} \, d(\nu_N - \nu) \to_p 0$ for $j = 1, \dots, c$,

then $(T_{N1}^*, \dots, T_{Nc}^*) \to (\int_0^1 L_{01} \, d\nu, \dots, \int_0^1 L_{0c} \, d\nu)$, a $N(0, \Sigma_0)$ random vector where Σ_0 is straightforwardly determined from (6.2).

Proof. Guided by (6.2) and (3.2) represent each L_{Nj} as the sum of c empirical processes (say U_{jn_j}) multiplied by random coefficients (say A_{Nj}) plus a negligible term (say δ_{Nj}). The proof is completed by applying the lemmas of Sections 2, 4, and 5 in a fashion analogous to their use in Theorem 5.1 to obtain that for each j, $T_{Nj} \rightarrow_p \int_0^1 L_{0j} d\nu$.

Let D^c denote the c-fold product $D \times \cdots \times D$. On this space consider the metrics $d^c(f, g) = \sum_{j=1}^c d(f_j, g_j)$ and $\rho^c(f, g) = \sum_{j=1}^c \rho(f_j, g_j)$ for all $f = (f_1, \dots, f_c), g = (g_1, \dots, g_c)$ in D^c . Clearly the analogue of Theorem 5.1(a) holds on D^c under the hypotheses of Theorem 6.1. Moreover, this result clearly holds for any metric $\rho^c(f,g)$ for which $\rho(f_{jN},g_{jN})\to 0$ as $N\to\infty$ for $j=1,\cdots,c$ implies $\rho^c(f_N, g_N) \to 0$ as $N \to \infty$; and similarly for d^c .

It is also clear that c-sample results analogous to Theorem 4.1 hold.

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