

A Distribution Free Version of the Smirnov Two Sample Test in the p-Variate Case

P. J. Bickel

The Annals of Mathematical Statistics, Vol. 40, No. 1. (Feb., 1969), pp. 1-23.

Stable URL:

http://links.jstor.org/sici?sici=0003-4851%28196902%2940%3A1%3C1%3AADFVOT%3E2.0.CO%3B2-5

The Annals of Mathematical Statistics is currently published by Institute of Mathematical Statistics.

Your use of the JSTOR archive indicates your acceptance of JSTOR's Terms and Conditions of Use, available at http://www.jstor.org/about/terms.html. JSTOR's Terms and Conditions of Use provides, in part, that unless you have obtained prior permission, you may not download an entire issue of a journal or multiple copies of articles, and you may use content in the JSTOR archive only for your personal, non-commercial use.

Please contact the publisher regarding any further use of this work. Publisher contact information may be obtained at <u>http://www.jstor.org/journals/ims.html</u>.

Each copy of any part of a JSTOR transmission must contain the same copyright notice that appears on the screen or printed page of such transmission.

The JSTOR Archive is a trusted digital repository providing for long-term preservation and access to leading academic journals and scholarly literature from around the world. The Archive is supported by libraries, scholarly societies, publishers, and foundations. It is an initiative of JSTOR, a not-for-profit organization with a mission to help the scholarly community take advantage of advances in technology. For more information regarding JSTOR, please contact support@jstor.org.

A DISTRIBUTION FREE VERSION OF THE SMIRNOV TWO SAMPLE TEST IN THE *p*-VARIATE CASE¹

BY P. J. BICKEL

University of California, Berkeley

1. Introduction. One of the classic problems of the theory of nonparametric inference is testing whether two samples come from the same or different populations. If the observations are univariate and we suppose only that the parent populations are governed by a continuous distribution function, genuinely distribution free and asymptotically consistent tests have been proposed by Smirnov, Von Mises, and Lehmann among others. For a review of the enormous literature that has grown up about these and related tests we refer to [1] and [3]. The multivariate case seems to have been studied far less fully. The obvious generalization of the univariate methods leads to procedures which are not distribution free (see [17]). Lehmann [12] has proposed a genuinely distribution free test consistent against all alternatives in this instance also, but his method involves post experimental randomization as an intrinsic factor and is therefore not fully satisfactory. Rosenblatt in [14] has given an ingenious solution to the goodness of fit problem in the multivariate case but his device does not seem to carry over to the two sample problem. Recently, David and Fix [2] and Sen and Chatterjee [16] have remarked that the classical permutation principle of Fisher, as used for instance by Wald and Wolfowitz in [19], when applied to rank tests in multivariate problems leads to procedures which are distribution free and which can, in principle at least, be tabled.

Our approach consists of combining this remark with the generalization of the classical Smirnov test to obtain a procedure which we show is distribution free and consistent against all alternatives. To achieve this aim we require a theorem on the "conditional" convergence of the empirical distribution function which is analogous to known results on permutation distributions (cf. [19], [8]). This theorem is then applied to yield known results of Kiefer and Wolfowitz [11] and Dudley [6].

The paper is organized as follows: Section 2 contains a formal statement of the model and our proposed procedure. In Section 3 we introduce the stochastic process considerations we need, state the main theorems on convergence of stochastic processes ((3.1) and (3.2)) and derive the consistency of our procedure from these as well as the previously mentioned results of Dudley and Kiefer and Wolfowitz. Section 4 deals with some auxiliary results on sampling from a finite population necessary for the proof of Theorem 3.1. Two of these results are due to Hájek ([8] and [9]), and one to Wald, Wolfowitz, Noether and Hoeffding

Received 11 September 1967.

¹ This paper was prepared with the partial support of the Office of Naval Research Contract N00014-67-A-0114-0004.

(see e.g. [7]). The others, Lemmas 4.3, 4.6, and 4.7, are new and may be of independent interest. The proof of Theorem 3.1 is completed in Sections 5, 6, and 7.

2. The model and the procedure. Let $X_1, X_2, \dots; Y_1, Y_2, \dots$, be two sequences of *p*-vectors (observations) defined on a measurable space (Ω, α) . For simplicity we take Ω to be $[R^p]^{\infty} \times [R^p]^{\infty}$ where $[R^p]^{\infty}$ is the space of ordered sequences of ordered *p* tuples of real numbers. Let α be the usual product Borel field on Ω , and let X_i, Y_i be the usual projection maps.

We suppose one of a family of probability measures P_{θ} to hold on (Ω, \mathfrak{A}) , where θ ranges over $\Theta = \{(F, G): F, G \text{ probability distribution functions on } R^p\}$ and if, $\theta = (F, G), P_{\theta}$ is the product probability measure on (Ω, \mathfrak{A}) which makes the X_i independent and identically distributed according to F and the Y_i similarly independent and identically distributed according to G. Let $\Theta_0 = \{\theta \in \Theta: \theta = (F, F) \text{ for some } F\}$. The problem we are interested in is testing the hypothesis, $H: \theta \in \Theta_0$ vs. $K: \theta \notin \Theta_0$, consistently at level α .

Formally we seek a sequence of (critical) functions $\varphi_{m,n} : [R^p]^m \times [R^p]^n \to [0, 1]$ which are Borel measurable such that for any given $0 < \alpha \leq 1$

(2.1)
$$E_{\theta}(\varphi_{m,n}(\mathbf{X}_1,\cdots,\mathbf{X}_m,\mathbf{Y}_1,\cdots,\mathbf{Y}_n)) \leq \alpha$$

for all $\theta \in \Theta_0$, and

(2.2)
$$\lim_{m,n\to\infty} E_{\theta}(\varphi_{m,n}(\mathbf{X}_1,\cdots,\mathbf{X}_m,\mathbf{Y}_1,\cdots,\mathbf{Y}_n)) = 1$$

for $\theta \not\in \Theta_0$, where E_{θ} indicates that expectation is taken under the assumption that θ holds.

We begin with some definitions. Let $F_m(\mathbf{t}, \omega)$ denote the empirical cumulative distribution of $\mathbf{X}_1(\omega), \cdots, \mathbf{X}_m(\omega)$ for $\omega \in \Omega$, and $\mathbf{t} \in \mathbb{R}^p$. Thus

(2.3)
$$F_m(\mathbf{t},\omega) = m^{-1} \sum_{j=1}^m I_{A(\mathbf{t})}(\mathbf{X}_j(\omega))$$

where $A(\mathbf{t}) = \{\mathbf{s} \in \mathbb{R}^p : \mathbf{s} \leq \mathbf{t}\}$ and $I_A(\mathbf{s})$ is the indicator function of a set A in \mathbb{R}^p . As is customary $\mathbf{s} \leq \mathbf{t}$ refers to coordinatewise inequality between $\mathbf{s} = (s_1, \dots, s_p)$ and $\mathbf{t} = (t_1, \dots, t_p)$.

Similarly, define

(2.4)
$$G_n(\mathbf{t},\omega) = n^{-1} \sum_{j=1}^n I_{A(\mathbf{t})}(\mathbf{Y}_j(\omega))$$

(2.5)
$$H_N(\mathbf{t},\omega) = \lambda_N F_m(\mathbf{t},\omega) + (1-\lambda_N)G_n(\mathbf{t},\omega),$$

where N = m + n, $\lambda_N = m N^{-1}$.

In the sequel we shall frequently refer to conditional probabilities and conditional expectations given H_N , written $P_{\theta}[\cdot | H_N]$ and $E_{\theta}(\cdot | H_N)$. We use this convenient notation to refer to the (regular) conditional probabilities of events in α and conditional expectations of functions measurable α when the conditioning sigma field

$$\mathfrak{G} = \{ B \varepsilon \mathfrak{A} \colon (\mathbf{x}_1, \cdots, \mathbf{x}_m, \mathbf{x}_{m+1}, \cdots; \mathbf{y}_1, \cdots, \mathbf{y}_n, \mathbf{y}_{n+1}, \cdots) \varepsilon B \\ \Leftrightarrow (P(\mathbf{x}_1), \cdots, P(\mathbf{x}_m), \mathbf{x}_{m+1}, \cdots; P(\mathbf{y}_1), \cdots, P(\mathbf{y}_n), \mathbf{y}_{n+1}, \cdots) \varepsilon B \}$$

for any P which is a map of the set $\{\mathbf{x}_1, \dots, \mathbf{x}_m, \mathbf{y}_1, \dots, \mathbf{y}_n\}$ onto itself. This is clearly the information contained in H_N , viz. the "values" of the first m and nobservations from our two samples without any identification of their "sample" origin. Clearly we can choose the regular conditional probability $P^{\mathfrak{B}}(\omega, \cdot)$ so that it depends on ω only through $H_N(\cdot, \omega)$.

It is evident that H_N is a sufficient statistic for $\theta \in \Theta_0$ if only $(\mathbf{X}_1, \dots, \mathbf{X}_m, \mathbf{Y}_1, \dots, \mathbf{Y}_n)$ has been observed, in the sense that if A is a Borel set in $[(R^p)]^m \times [(R^p)]^n$

$$P_{\theta}[(\mathbf{X}_1, \cdots, \mathbf{X}_m, \mathbf{Y}_1, \cdots, \mathbf{Y}_n) \varepsilon A \mid H_N]$$

can be chosen independent of θ for $\theta \in \Theta_0$. In fact if $\theta \in \Theta_0$ given H_N there are a finite number of equally likely values that $(X_1, \dots, X_m, Y_1, \dots, Y_n)$ can assume.

In what follows we shall always assume that $P_{\theta}[\cdot | H_N]$ is computed under $\theta \in \Theta_0$ and we shall accordingly drop the subscript θ in all subsequent conditional laws, probabilities, and expectations.

Finally for any two distribution functions F and G on \mathbb{R}^p we define the Kolmogorov (L_{∞}) distance between F and G as usual by,

(2.6)
$$d(F,G) = \sup_{\mathbf{t}} |F(\mathbf{t}) - G(\mathbf{t})|.$$

We can now define $\varphi_{m,n}$ by,

(2.7)

$$\varphi_{m,n}(\mathbf{X}_{1}(\omega), \cdots, \mathbf{X}_{m}(\omega), \mathbf{Y}_{1}(\omega), \cdots, \mathbf{Y}_{n}(\omega)) = 1$$
if $d(F_{m}(\cdot, \omega), G_{m}(\cdot, \omega)) > c_{N}(H_{N}(\cdot, \omega))$

$$= \gamma_{N}(H_{N}(\cdot, \omega))$$
if $d(F_{n}(\cdot, \omega), G_{n}(\cdot, \omega)) = c(H_{N}(\cdot, \omega))$

$$= 0 \text{ otherwise,}$$

where c and γ are uniquely determined by the condition that they be the smallest such numbers satisfying,

(2.8)
$$E\{\varphi(\mathbf{X}_1, \cdots, \mathbf{X}_m, \mathbf{Y}_1, \cdots, \mathbf{Y}_n) \mid H_N(\cdot, \omega)\} = \alpha.$$

This sequence clearly satisfies (2.1) and it follows from Theorems 3.1 and 3.2 of the next section that it also satisfies (2.2).

We remark that, at least in principle, the distribution of the test statistic under the null hypothesis can be tabled.

Suppose for simplicity that $F(\cdot)$ is continuous. We may assume that under $\theta_0 = (F, F)$ with probability one all the coordinates of the X's and Y's are dis-

tinct. Let

$$\mathbf{X}_{i}(\omega) = (X_{i1}(\omega), \cdots, X_{ip}(\omega))$$

and define \mathbf{Y}_i similarly. Let $R_{ij}(\omega)$ be the rank of $X_{ij}(\omega)$ among the set of numbers $[\{X_{kj}(\omega)\}, 1 \leq k \leq m] \cup [\{Y_{kj}(\omega)\}: 1 \leq k \leq n].$

Similarly define $S_{ij}(\omega)$ as the rank of Y_{ij} . Let $\mathbf{R}_i = (R_{i1}, \dots, R_{ip})$ and define \mathbf{S}_i similarly. Define

$$F_m^{R}(\mathbf{t},\omega) = m^{-1} \sum_{i=1}^m I_{A(\mathbf{t})}(\mathbf{R}_i(\omega)).$$

Then if $H_N^R(\mathbf{t}, \omega)$, $G_n^R(\mathbf{t}, \omega)$ are defined in analogy to H_N and G_n , clearly, under θ_0 , with probability one,

(2.9)
$$\mathfrak{L}(d(F_m, G_n) \mid H_N) = \mathfrak{L}(d(F_m^R, G_n^R) \mid H_N^R))$$

where $\mathfrak{L}(Z)$ denotes the law of a random variable Z. Now H_N^R takes on exactly $(N!)^{p-1}$ distinct values and for each value of H_N^R , F_m^R takes on at most $\binom{N}{m}$ possible distinct values.

3. Limiting behavior of F_n given H_N and statement of principal theorems. We begin this section by introducing a more general framework for our results. Let P_{1N}, \dots, P_{NN} be N points in \mathbb{R}^p , where we allow repetitions among the

Let P_{1N} , ..., P_{NN} be N points in K, where we allow repetitions allong the points.

In accord with our usage in the last section let H_N be the distribution function of the measure which assigns mass N^{-1} to each point P_{iN} .

Consider the experiment of choosing at random m of the points P_{1N}, \dots, P_{NN} . Again in accord with previous usage let F_m be the distribution function of the measure which assigns mass m^{-1} to each of the m selected points, and let G_n be the distribution function of the measure which assigns mass n^{-1} to each of the n remaining points where we suppose m + n = N. As we noted in the previous section if P_{1N}, \dots, P_{NN} denote the values of the "pooled" sample of $X_1, \dots, X_m, Y_1, \dots, Y_n$ where the X's and Y's originate from a common distribution F, then the conditional distribution of the empirical distribution function of the X's is precisely that of F_m constructed above. However, for the first two theorems of this section we shall make no assumptions on how the points P_{1N}, \dots, P_{NN} are obtained. For convenience we introduce a suitable probability space $(\tilde{\Omega}, \tilde{\Omega}, \tilde{P})$ with members $\tilde{\omega}$, and write $F_m(\mathbf{t}, \tilde{\omega})$ to indicate the dependence of F_m on $\mathbf{t} \in \mathbb{R}^p$ as a distribution function, and on the particular subset of P_{1N}, \dots, P_{NN} selected by the random experiment through $\tilde{\omega}$. In this formulation H_N is a function of \mathbf{t} only and is not stochastically determined. Of course,

(3.1)
$$G_n(\mathbf{t}, \tilde{\omega}) = n^{-1} \{ NH_N(\mathbf{t}) - mF_m(\mathbf{t}, \tilde{\omega}) \}.$$

The triple $(\tilde{\Omega}, \tilde{\alpha}, \tilde{P})$ whose formal description we leave to the reader is so chosen that $F_m(\cdot, \cdot)$ is a separable stochastic process on \mathbb{R}^p whose finite dimensional joint laws are determined in accord with the experiment we have just described. We take $\{m\}$ and $\{N - m\}$ to be unbounded sequences of natural numbers whose dependence on N we suppress. In all future theorems when we say $N \to \infty$ we mean that m and n = N - m both tend to ∞ as well. Without loss of generality we suppose that all the processes $F_m(\cdot, \cdot)$ are defined on the same probability space $(\tilde{\Omega}, \tilde{\alpha}, \tilde{P})$.

Define the stochastic process $X_N(\mathbf{t}, \tilde{\omega})$ on $(\tilde{\Omega}, \tilde{\alpha})$ and \mathbb{R}^p by

(3.2)
$$X_{N}(\mathbf{t},\tilde{\omega}) = [mn/N]^{\frac{1}{2}}[F_{m}(\mathbf{t},\tilde{\omega}) - G_{n}(\mathbf{t},\tilde{\omega})]$$
$$= [mn^{-1}]^{\frac{1}{2}}N^{\frac{1}{2}}[F_{m}(\mathbf{t},\tilde{\omega}) - H_{N}(\mathbf{t},\tilde{\omega})].$$

Let *H* be the distribution function of a probability measure on \mathbb{R}^p (which we shall also refer to as *H*). We follow Dudley [5] and define Q_H to be the set of real valued functions $\alpha(\cdot)$ on \mathbb{R}^p which:

(i) Are continuous except possibly on hyperplanes of the form $\{(t_1, \cdots, t_p): t_j = \text{constant}\}$ which have positive *H*-measure.

(ii) On hyperplanes as above are continuous from below with limits from above, where "below" and "above" refer to strict inequality in all coordinates.

(iii) Have limits at points of the boundary of the compactification \bar{R}^p of R^p where $\bar{R} = [-\infty, +\infty]$.

Endow Q_H with the supremum norm making it a separable Banach space and let Σ_H be the topological Borel field of Q_H . Then ([5]) there exists a probability measure μ_H on (Q_H, Σ_H) such that the map $X(\mathbf{t}, \alpha)$ defined for $\alpha \in Q_H$ and $\mathbf{t} \in \mathbb{R}^p$ and given by,

(3.3)
$$X(\mathbf{t}, \alpha) = \alpha(\mathbf{t}),$$

is Gaussian with mean 0 and the stochastic process $X(\mathbf{t}, \cdot)$ has covariance structure given by,

(3.4)
$$\operatorname{Cov} (X(\mathbf{t}, \cdot), X(\mathbf{t}', \cdot)) = H(\mathbf{t} \wedge \mathbf{t}') - H(\mathbf{t})H(\mathbf{t}'),$$

where $\mathbf{t} \wedge \mathbf{t}'$ is the vector $(\min(t_1, t_1'), \cdots, \min(t_p, t_p'))$. Finally let $\mathfrak{M}(\mathbb{R}^p)$ be the set of all bounded real valued functions on \mathbb{R}^p endowed with the supremum norm, $\|\cdot\|$. We require that $(\tilde{\Omega}, \tilde{\alpha})$ have the property that $\{\tilde{\omega}: F_m(\cdot, \tilde{\omega}) = v\} \in \tilde{\alpha}$ for every $v \in \mathfrak{M}(\mathbb{R}^p)$. We shall say that a real valued functional h on $\mathfrak{M}(\mathbb{R}^p)$ is regular for H (a distribution function) if h is continuous in the supremum norm at all points of a Borel subset S of Q_H such that $\mu_H(S) = 1$. Note that for any functional h the function $h(X_N(\cdot, \tilde{\omega}))$ is measurable in $\tilde{\omega}$ in view of our assumptions on $(\tilde{\Omega}, \tilde{\alpha})$. We can now state our main result.

THEOREM 3.1 Suppose that $||H_N - H|| \rightarrow 0$ for some distribution function H. Then, for every h regular for H,

(3.5)
$$\Delta\{\mathfrak{L}(h(X_N(\cdot, \cdot))), \mathfrak{L}_{\mathbf{H}}(h(X(\cdot, \cdot)))\} \to \mathbf{0}$$

as $N \to \infty$, where $\mathfrak{L}(h(X_N(\cdot, \cdot)))$ is the law of $h(X_N(\cdot, \cdot))$ under \tilde{P} , $\mathfrak{L}_H(h(X(\cdot, \cdot)))$ is the law of $h(X(\cdot, \cdot))$ under H and Δ is the Lévy-Prohorov distance between laws (distribution functions) on the real line.

In the terminology of (Dudley [6]) the measures induced by the processes X_N on $\mathfrak{M}(\mathbb{R}^p)$ converge weak^{*} to μ_H . We postpone the proof of this theorem to the next four sections.

Suppose now that $X_1, \dots, X_m, Y_1, \dots, Y_n$ are random vectors defined as

in the previous section and (with a slight abuse of notation) F_m , G_n , and H_N are the empirical cumulative distribution functions of the X's, the Y's and the pooled sample respectively. Define the stochastic process X_N in terms of F_m and G_n by equation (3.2). Note that these are all (including H_N) functions on $\mathbb{R}^p \times \Omega$.

We shall say that a real valued functional h is regular for the two sample problem and H if h is regular for H and furthermore $h(X_N(\cdot, \omega))$ is a measurable function of ω on Ω for all N. Note that ω is permitted to vary in H_N also with this definition. We have (using the notation of the previous section),

THEOREM 3.2. Suppose that $\theta = (F, F)$. Let h be a functional regular for the two sample problem and F. Then,

$$(3.6) \qquad \Delta\{\mathfrak{L}(h(X_N(\cdot, \cdot) | H_N)), \mathfrak{L}_{\mathbb{F}}(h(X(\cdot, \cdot)))\} \to 0 \qquad \text{a.s. } P_{\theta}.$$

PROOF. In view of our characterization of the conditional distribution of F_m given H_N if $\theta = (F, F)$ Theorem 3.2 is an immediate consequence of Theorem 3.1 and,

(3.7)
$$P_{(F,F)}[||H_N - F|| \to 0] = 1.$$

Assertion (3.7) is nothing else than the Glivenko-Cantelli theorem for vector variables. Theorem 3.2 is proved.

Put more prosaically our theorem states that for almost all sample sequences the conditional distributions of $h(X_N(\cdot, \cdot))$ given H_N (which do not depend on θ) converge to the same limiting distribution, that of $h(X(\cdot, \cdot))$, when μ_F holds on (Q_F, Σ_F) . The final dependence on θ of course appears because of the dependence of H_N on θ . It seems worth noting that the restriction of regularity for the two sample problem on h may be dropped if the σ -field α is completed with respect to P_{θ} since we have remarked that given H_N , F_m and G_n can take on only finitely many values and consequently the conditional distribution of h given H_N is well defined. We can now establish,

THEOREM 3.3. The sequence of tests $\{\varphi_{m,n}\}$ satisfies (2.1) and (2.2).

PROOF. As we noted previously, (2.1) clearly holds. Let h(q) = ||q||. *h* is continuous on all of $\mathfrak{M}(\mathbb{R}^p)$ and easily satisfies (i).

Now,

(3.8)
$$h(X_N(\cdot, \omega)) = \{mnN^{-1}\}^{\frac{1}{2}} \|F_m - G_n\|.$$

Then for every distribution function H on \mathbb{R}^p , $\alpha > 0$ there exists a unique finite smallest $c(H, \alpha)$, $0 \leq \gamma(H, \alpha) \leq 1$, such that,

$$(3.9) \quad \mu_{\mathrm{H}}\{h(X(\cdot, \cdot)) > C(H, \alpha)\} + \gamma(H, \alpha)\mu_{\mathrm{H}}\{h(X(\cdot, \cdot)) = c(H, \alpha)\} = \alpha.$$

From (3.9) and (3.6) we see that if $\theta = (F, F)$,

(3.10)
$$\{mn(m+n)^{-1}\}^{\frac{1}{2}}c_N(H_N) \to c(F,\alpha) \quad \text{a.s. } P_{\theta}$$

where $c_N(H_N)$ is given by (2.7), (2.8).

Now, suppose $\theta = (F, G), F \neq G$. Without loss of generality let us assume that

 $mN^{-1} \rightarrow \lambda$, where $0 \leq \lambda \leq 1$. Define,

(3.11)
$$H = \lambda F + (1 - \lambda)G.$$

By the Glivenko-Cantelli theorem in this case,

(3.12)
$$P_{\theta}[\|F_m - F\| \to 0] = P_{\theta}[\|G_n - G\| \to 0] = 1$$

and hence,

$$(3.13) P_{\theta}[||H_N - H|| \rightarrow 0] = 1.$$

An application of Theorem 3.1 identical to that used to obtain Theorem 3.2 and (3.10) now enables us to conclude that,

(3.14)
$$\{mn(m+n)^{-1}\}^{\frac{1}{2}} c_N(H_N) \to c(H, \alpha) \quad \text{a.s. } P_{\theta}.$$

On the other hand by (3.12),

(3.15)
$$||F_m - G_n|| \to ||F - G|| > 0$$
 a.s. P_{θ} .

Of course, (3.14) and (3.15) imply that,

(3.16)
$$\varphi_{m,n}(\mathbf{X}_1, \cdots, \mathbf{Y}_1, \cdots) \to 1$$
 a.s. P_{θ}

if $F \neq G$. Theorem 3.3 follows.

There are many possible choices of $\varphi_{m,n}$ satisfying (2.1) and (2.2), some perhaps based on statistics having a limiting distribution independent of F if $\theta = (F, F)$. Certainly plausible statistics whose distribution is more invariant than that of $d(F_m, G_n)$ can readily be defined (see for instance Weiss [20] and Vincze [18]). We hope to investigate this question in a later paper. We remark that Theorems 3.2 and 3.3 are of some interest even in the case p = 1 since the test $\varphi_{m,n}$ which for F continuous is the usual Smirnov test is also distribution free when the underlying distribution is discrete, viz. when there are ties. Unfortunately, and this is, of course, the principal difficulty in higher dimensions, the limiting value of the cut-off point of the test depends on F if F is not continuous.

From Theorem 3.3 we immediately deduce,

THEOREM 3.4. If h satisfies the hypotheses of Theorem 3.2 and $\theta = (F, F)$ then

$$(3.17) \qquad \Delta\{\mathfrak{L}_{\theta}(h(X_{N}(\cdot, \cdot))), \mathfrak{L}_{F}(h(X(\cdot, \cdot)))\} \to 0.$$

PROOF. Let,

(3.18)
$$\Psi_N(s, H_N) = E(\exp is h(X_N(\cdot, \cdot)) | H_N).$$

Then,

$$(3.19) P_{\theta}[\Psi_N(s, H_N) \to E_F\{\exp is h(X(\cdot, \cdot))\} \forall s] = 1.$$

Theorem 3.4 follows upon applying the dominated convergence theorem to the variables $\Psi_N(s, H_N)$. From Theorem 3.4 we can derive Dudley's Theorem 2 [6]. Let,

(3.20)
$$\widetilde{X}_m(\mathbf{t},\omega) = m^{\frac{1}{2}}(F_m(\mathbf{t},\omega) - F(\mathbf{t},\omega)).$$

We shall say h is regular for the one sample problem and F if h is regular for F and $h(X_m(\cdot, \omega))$ is measurable in ω .

THEOREM 2. If h is regular for the one sample problem and F and $\theta = (F, F)$ then,

$$(3.21) \quad \mathfrak{L}_{\theta}\{h(\tilde{X}_{m}(\cdot, \cdot))\} \rightarrow \mathfrak{L}_{\mu_{F}}\{h(X(\cdot, \cdot))\} \quad as \quad m \rightarrow \infty.$$

PROOF. By Proposition 1 of [6] we need only consider h continuous on $\mathfrak{M}(\mathbb{R}^p)$. For such an h, if $\theta = (F, F)$,

$$(3.22) \qquad \qquad \mathfrak{L}_{\theta}\{h(X_{N}(\cdot, \cdot))\} \to \mathfrak{L}_{\theta}\{h(X_{m}(\cdot, \cdot))\} \qquad \text{as } n \to \infty,$$

(m fixed) by the Glivenko-Cantelli theorem. Therefore any limit law stochastic or substochastic of the sequence $\mathfrak{L}_{\theta}\{h(\tilde{X}_m(\cdot, \cdot))\}$ as $m \to \infty$ is a limit law of the sequence $\mathfrak{L}_{\theta}\{h(X(\cdot, \cdot))\}$ as $m, n \text{ both } \to \infty$. Theorem 2 follows. Of course, Theorem 3.3 could have been derived from Theorem 2.

From Theorem 2 one can, of course, derive Theorem 2 of [11] by applying the former to the functionals h(q) = ||q|| and $h(q) = ||\max(q, 0)||$. The limiting law of our test statistics $\{mn(m+n)^{-1}\}^{\frac{1}{2}} ||F_m - G_n||$ under $\theta = (F, F)$ given H_N is then identified as being the same as the distribution given in [11], and Theorem 1 of [11] and 1-m of [10] may be employed to give some information about $c(F, \alpha)$. It seems plausible that one can obtain bounds similar to those of [10] and [11] for the conditional distribution of $\{mn(m+n)^{-1}\}^{\frac{1}{2}}d$ (F_m, G_n) in terms of H_N but we have not attempted to do so. Finally, we note that Theorems 2 and 3.2 imply that the expressions for the limit law of $m^{\frac{1}{2}}||F_m - F||$ if p = 1 given by (Schmid [15]) apply in our case also.

4. Some results on sampling without replacement. The lemmas of this section play a fundamental role in the proof of Theorem 3.1 given in the next three sections. In order to place Lemmas 4.6 and 4.7 in the scheme of things it may be worthwhile to read Sections 6 and 7 first.

Let V_{1N}, \dots, V_{NN} be a sample taken in order from a finite population c_{1N}, \dots, c_{NN} without replacement, viz. $P[(V_{1N}, \dots, V_{NN}) = (c_{i_1N}, \dots, c_{i_NN})] = [N!]^{-1}$ for every permutation (i_1, \dots, i_N) of $(1, \dots, N)$. Suppose $\sum_{i=1}^{N} c_{i_N} = 0$. Then, we have,

LEMMA 4.1. (Hájek [8]). Let V_{1N} , \cdots , V_{NN} , $\{c_{iN}\}$ be as above. Let $\{d_{iN}\}$, $1 \leq i \leq N$, be another sequence of constants such that $\sum_{i=1}^{N} d_{iN} = 0$. Suppose $N^{-1} \sum_{i=1}^{N} c_{iN}^2 \to \sigma^2$, $\sum_{i=1}^{N} d_{iN}^2 \to K$. Then,

(4.1)
$$\mathfrak{L}(\sum_{i=1}^{N} d_{iN} V_{iN}) \to \eta(0, \sigma^2 K)$$

if,

(4.2) (a)
$$N^{-\frac{1}{2}} \max |c_{iN}| \to 0$$
, (b) $\max |d_{iN}| \to 0$,

and

(4.3)
$$N^{-1} \sum_{\{|\delta_{N,ij}| > \tau\}} \delta_{N,ij}^2 \to 0$$

for every $\tau > 0$, where $\delta_{Nij} = c_{iN} d_{jN}$.

Another useful result also due to (Hájek [9]) is,

LEMMA 4.2. (Hájek). Let $\{V_{iN}\}$, $\{c_{iN}\}$ be as above. Suppose $S_j = \sum_{i=1}^{j} V_{iN}$. Then,

(4.4)
$$P[\max_{1 \le j \le N} |S_j| \ge \epsilon] \le E |S_n|^r \{\epsilon (1 - nN^{-1})\}^{-r}$$

for every $r \geq 0, 1 \leq n \leq N, \epsilon > 0$.

This lemma is of little value if n is close to N. To deal with that situation we need inequalities of the type given in Lemma 4.3.

LEMMA 4.3. Under the assumption of Lemma 4.2 suppose $kl \leq n < k(l+1)$ where $k \geq 2$, $N \geq n$, $l \geq 1$. Then there exists c(k, r) such that,

(4.5)
$$P[\max_{1 \le j \le n} |S_j| \ge \epsilon] \le c(k, r) \epsilon^{-r} E |S_l|^r.$$

PROOF. Suppose n = kl + d. Then,

$$(4.6) \quad P[\max_{1 \le j \le n} |S_j| \ge \epsilon] \le P[\max_{1 \le j \le kl} |S_j| \ge \epsilon/2]$$

$$+ P[\max_{kl < j \leq n} |S_{kl} + (S_j - S_{kl})| \geq \epsilon/2].$$

But,

$$(4.7) \quad P[\max_{kl < j \le n} |S_{kl} + (S_j - S_{kl})| \ge \epsilon/2]$$
$$\leq P[|S_{kl}| \ge \epsilon/4] + P[\max_{1 \le j \le d} |S_j| \ge \epsilon/4].$$

Hence, if $0 \leq d < k \leq kl$ we have,

(4.8)
$$P[\max_{1 \le j \le n} |S_j| \ge \epsilon] \le 3P[\max_{1 \le j \le kl} |S_j| \ge \epsilon/4],$$

and we need establish the lemma only for n = kl. We proceed by induction. If k = 2,

$$(4.9) \quad P[\max_{1 \le j \le 2l} |S_j| \ge \epsilon] \le P[\max_{1 \le j \le l} |S_j| \ge \epsilon/2] + P[|S_l| \ge \epsilon/4] \\ + P[\max_{1 \le j \le l} |S_j| \ge \epsilon/4] \le 3P[\max_{1 \le j \le l} |S_j| \ge \epsilon/4] \le 3.8^r \epsilon^{-r} E|S_l|^r$$

by Lemma 4.2. The induction step is given by (4.6) and (4.7) with d = l.

We will also need,

LEMMA 4.4. (Wald-Wolfowitz-Noether-Hoeffding.) Suppose the c_{iN} satisfy the hypotheses of Lemma 4.1 as well as (4.2) (a), S_n is as in Lemma 4.2, and $n/N \rightarrow \lambda$, $0 < \lambda < 1$. Then as $N \rightarrow \infty$,

(4.10)
$$E(N^{-k}S_n^{2k}) \to \lambda^k \sigma^{2k} \mu_{2k}$$

where μ_{2k} is the 2kth central moment of the standard normal distribution.

A proof of Lemma 4.4 may be found in [7] pp 237–239.

Before giving the main results of this section we state and prove an auxiliary lemma. Our techniques in this and the succeeding lemma are similar to those used by Hájek in [9] and by Kiefer and Wolfowitz in [11].

LEMMA 4.5. Let $\{V_{iN}\}, 1 \leq i \leq N$, be as above. Then,

$$P[\sum_{k=r+1}^{r+m} V_{kN}$$

$$(4.11) < g(V_{1N}, \dots, V_{rN}) | V_{1N}, \dots, V_{rN}] \leq [g(V_{1N}, \dots, V_{rN})$$

$$+ mS_r(N-r)^{-1}]^{-2}m(1-m(N-r)^{-1})(N-r)^{-1}\sum_{i=1}^{N} c_{iN}^2$$

for any $g(V_{1N}, \dots, V_{rN})$ such that $g(V_{1N}, \dots, V_{rN}) + mS_r(N-r)^{-1} < 0$.

PROOF. Given V_{1N} , \cdots , V_{rN} the V_{kN} with k > r are distributed as a sample without replacement from the population of size N - r whose members c_{iN} , $i = 1, \dots, N - r$, are the c_{iN} which have not appeared among the V_{iN} with $i \leq r$. Then

(4.12)
$$E(V_{(r+1)N} | V_{iN}, 1 \leq i \leq r) = -S_r(N-r)^{-1}.$$

Chebyshev's inequality applied to the left-hand side of (4.11) gives,

$$P[\sum_{k=i+1}^{r+m} V_{kN} \leq g(V_{1N}, \cdots, V_{rN}) | V_{1N}, \cdots, V_{rN}]$$

$$(4.13) \leq [g(V_{1N}, \cdots, V_{rN}) + mS_r(N - r)^{-1}]^{-2}$$

$$(1 - (m - 1)/(N - r - 1))m \operatorname{Var}\{V_{r+1,N} | V_{1N}, \cdots, V_{rN}\}.$$

Finally,

$$(4.14) \quad \operatorname{Var}\{V_{r+1,N} \mid V_{1N}, \cdots, V_{rN}\} \\ = (N-r)^{-1}\{\sum_{i=1}^{N} (c_{iN} + S_r(N-r)^{-1})^2 - \sum_{i=1}^{r} (V_{iN} + S_r(N-r)^{-1})^2\} \\ = (N-r)^{-1}\{\sum_{i=1}^{N} c_{iN}^2 - \sum_{i=1}^{r} V_{iN}^2 - S_r^2(N-r)^{-1}\} \\ \le (N-r)^{-1} \sum_{i=1}^{N} c_{iN}^2. \qquad \qquad \square$$

The final results of this section are further extensions of Lemmas 4.2 and 4.3 which are central to our argument. Let S be a fixed subset with n members $(n \leq N)$ of the set of all ordered p-tuples (i_1, \dots, i_p) where the i_j may range independently over 1, \dots , N. Now, relabel V_{1N}, \dots, V_{nN} as $V_{(i_1,\dots,i_p)}$ where (i_1, \dots, i_p) ranges over S. We shall call such sets S, configurations. Define,

(4.15)
$$S_{(j_1,\dots,j_p)} = \sum_{(i_1,\dots,i_p) \in [A(j)] \cap S} V_{(i_1,\dots,i_p)}$$

where $\mathbf{j} = (j_1, \dots, j_p)$ is a vector with natural number entries and $A(\mathbf{j})$ is defined in Section 2.

Define,

$$(4.16) a(\gamma, \epsilon) = [1 - n(N - n)^{-1} - \gamma]\epsilon$$

and

(4.17)
$$b(\gamma, \epsilon) = 1 - n(N-n)^{-1}[a(\gamma, \epsilon)]^{-2} \sum_{i=1}^{N} c_{iN}^{2}.$$

Choose γ so that $\gamma \leq 1 - n(N - n)^{-1}$ and $b(\gamma, \gamma^s \epsilon) > 0$ for $0 \leq s \leq p - 2$.

We have,

LEMMA 4.6 Under the above assumptions, if
$$p \ge 2$$
, $n < N$, and $r > 0$,

(4.18)
$$P[\max\{|S_{(i_1,\dots,i_p)}|: 1 \leq i_j \leq N, j = 1, \dots, p\} \geq \epsilon]$$

 $\leq \{\prod_{s=0}^{p-2} b^{-1}(\gamma, \gamma^s \epsilon)\} E|S_n|^r [\gamma^{p-1} \epsilon]^{-r}.$

PROOF. We prove the lemma for p = 2. The general case follows by an induction which we leave to the reader.

Let t_1 be the smallest *i* such that $S_{(i,j)} \ge \epsilon$ for some $1 \le j \le N$ if such an *i* exists. Otherwise let $t_1 = \infty$. Let t_2 be the largest *j*, $1 \le j \le N$, such that $S_{(t_1,j)} \ge \epsilon$, if $t_1 < \infty$. Otherwise let $t_2 = \infty$. Then,

$$P[\max\{S_{(i,j)} : 1 \leq i \leq N, 1 \leq j \leq N\} \geq \epsilon]$$

$$(4.19) = \sum_{k=1}^{N} P[t_1 = k]$$

$$\leq \sum_{k=1}^{N} P[t_1 = k, S_{(N,t_2)} < \gamma \epsilon] + P[S_{(N,t_2)} \geq \gamma \epsilon, t_1 < \infty]$$

$$\leq \sum_{k=1}^{N} P[S_{(N,t_2)} < \gamma \epsilon, t_1 = k] + P[\max_{1 \leq j \leq N} S_{(N,j)} \geq \gamma \epsilon].$$

But,

(4.20)
$$P[S_{(N,t_2)} < \gamma \epsilon, t_1 = k]$$

= $\int_{[t_1=k]} P[S_{(N,t_2)} < \gamma \epsilon \mid V_{(i,j)}, 1 \le i \le k, 1 \le j \le N, (i,j) \varepsilon S] dP$

where $P[\cdot | \cdot]$ denotes as usual a conditional probability. For convenience denote the conditioning event in (4.20) by A.

Let n(i, j) be the number of summands in $S_{(i, j)}$ and define,

(4.21)
$$\rho(i,j) = (n(N,j) - n(i,j))[N - n(i,N)]^{-1}.$$

By definition of t_2 we have $S_{(t_1,N)} \leq S_{(t_1,t_2)}$ and therefore,

$$(4.22) \quad P[S_{(N,t_2)} < \gamma \epsilon \,|\, A] \leq P[(S_{(N,t_2)} - S_{(t_1,t_2)} + S_{(t_1,N)}\rho(t_1,t_2) < \gamma \epsilon - (1 - \rho(t_1,t_2))S_{(t_1,t_2)} \,|\, A].$$

Applying Lemma 4.5 to the right-hand side of (4.22) we obtain,

$$P[S_{(N,t_2)} < \gamma \epsilon | A]$$

$$\leq [\gamma \epsilon - (1 - \rho(t_1, t_2))S_{(t_1,t_2)}]^{-2}[1 - (n(N, t_2)$$

$$(4.23) - n(t_1, t_2) - 1)(N - n(t_1, N) - 1)^{-1}](N - n(t_1, N))^{-1}$$

$$\{\sum_{i=1}^{N} c_{iN}^2\}(n(N, t_2) - n(t_1, t_2))$$

if $\gamma \epsilon - (1 - \rho(t_1, t_2))S_{(t_1, t_2)} < 0.$

But,

(4.24)
$$0 \leq \rho(i,j) \leq n[N-n]^{-1}$$

and $S_{(t_1, t_2)} \geq \epsilon$. From these remarks we conclude,

$$P[S_{(N,t_2)} < \gamma \epsilon \mid A]$$

$$(4.25) \leq \{ [\gamma - (1 - n(N - n)^{-1})] \epsilon \}^{-2} n(N - n)^{-1} \sum_{i=1}^{N} c_{iN}^2,$$
for $\alpha \leq 1 - n(N - n)^{-1}$

Substituting (4.25) in (4.20) we get,

(4.26)
$$P[S_{(N,t_2)} < \gamma \epsilon, t_1 = k]$$

 $\leq P[t_1 = k] \{ [\gamma - (1 - n(N - n)^{-1})] \epsilon \}^{-2} n(N - n)^{-1} \sum_{i=1}^{N} c_{iN}^2.$
Finally, from (4.19) and (4.26)

$$P[\max\{S_{(i,j)}: 1 \leq i \leq N, 1 \leq j \leq N\} \geq \epsilon]$$

$$(4.27) \leq P[\max_{1 \leq j \leq N} S_{(N,j)} \geq \gamma \epsilon] + [a(\gamma, \epsilon)n (N-n)^{-1} \sum_{i=1}^{N} c_{i,i}^2] \cdot P[\max\{S_{(i,j)} : 1 \leq i \leq N, 1 \leq j \leq N\} \geq \epsilon].$$

Applying (4.27) to sampling from the population $\{-c_{IN}, \cdots, -c_{NN}\}$ adding the resulting inequalities and transposing we obtain,

(4.28)
$$P[\max\{|S_{(i,j)}|: 1 \leq i \leq N, 1 \leq j \leq N\} \geq \epsilon]$$
$$\leq b^{-1}(\gamma, \epsilon)P[\max_{1 \leq j \leq N}|S_{(N,j)}| \geq \gamma\epsilon].$$

The proof for p = 2 is finished upon applying Lemma 4.2 to the right hand side of (4.28).

It is important to note that the right hand side of (4.18) depends only on the c_{iN} and n and is independent of the configuration chosen.

Finally Lemma 4.7 extends Lemma 4.3. We use the notation of Lemma 4.6. **LEMMA 4.7.** Under the assumptions of Lemma 4.6 there exist c(k, r), $\lambda(k)$ such that, if $kl \leq n < k(l+1)$, $k \geq 3$, r > 0, and $l \geq 1$, then

(4.29)
$$P[\max\{|S_{(i_1,\dots,i_p)}|: 1 \leq i_j \leq N, j = 1, \dots, p\} \geq \epsilon]$$

 $\leq c(k, r)(1 - \lambda(k)\epsilon^{-2} \sum_{i=1}^{N} c_{iN}^2)^{-(p-1)}\epsilon^{-r}E|S_l|^r.$

PROOF. With our new notation the argument is in fact simpler than that given for Lemma 4.3. Suppose n = kl + d, with $0 \le d < k \le l$. Divide S into (k + 1) subsets $S^{(1)}, \dots, S^{(k+1)}$ such that $S^{(1)}, \dots, S^{(k)}$ have l members each while $S^{(k+1)}$ has d members. Define,

(4.30)
$$S_{(j_1,\dots,j_p)}^{(t)} = \sum_{(i_1,\dots,i_p) \in [A(j)] \cap S^{(t)}} V_{(i_1,\dots,i_p)}$$

where $j = (j_1, \dots, j_p)$ for $t = 1, \dots, k + 1$. Then,

(4.31)
$$S_{(i_1,\dots,i_p)} = \sum_{j=1}^{k+1} S_{(i_1,\dots,i_p)}^{(j)},$$

and

$$(4.32) \quad P[\max\{|S_{(i_1,\dots,i_p)}|: 1 \le i_j \le N, j = 1, \dots, p\} \ge \epsilon] \\ \le \sum_{i=1}^{k+1} P[\max\{|S_{(i_1,\dots,i_p)}^{(t)}|: 1 \le i_j \le N, j = 1, \dots, p\} \ge \epsilon (k+1)^{-1}].$$

Since each $S^{(l)}$ is a configuration with $\leq l$ members, we have for suitable γ , by Lemma 4.6,

$$P[\max\{|S_{(i_1,\dots,i_p)}|: 1 \leq i_j \leq N, j = 1, \dots, p\} \geq \epsilon]$$

$$(4.33) \qquad \leq (k+1) \left[\prod_{s=0}^{p-2} (1-(k-1)^{-1}) \\ \left[(1-(k-1)^{-1}-\gamma) \gamma^s \epsilon/(k+1)\right]^{-2} \sum_{i=1}^{N} c_{iN}^2\right]^{-1} \\ \cdot \{E|S_l|^r\} \gamma^{-r(p-1)} \epsilon^{-r}(k+1)^r.$$

(It is clear that the bound on the right hand side of (4.18) continues to be valid if the number of members of S is $\leq n$.) Since $k \geq 3$ we may take $\gamma = \frac{1}{4}$ and hence $\lambda(k) = [16(k-1)(k+1)^2]^{-1}$ and $c(k, r) = (k+1)^{r+1}4^{r(p-1)}$ will do. If l < k a further obvious decomposition of $S^{(k+1)}$ will complete the proof.

Both Lemma 4.3 and this lemma will be used for k = 3 only.

5. Proof of Theorem 3.1: Initial considerations and enumeration of cases. It follows from the work of Dudley (Theorem 1 and Proposition 1 of [6]) that to establish Theorem 3.1 we need only show,

(i) Finite dimensional convergence:

(5.1)
$$\mathfrak{L}(X_N(\mathbf{t}^{(1)}, \cdot), \cdots, X_N(\mathbf{t}^{(r)}, \cdot)) \to \mathfrak{L}_H(X(\mathbf{t}^{(1)}, \cdot), \cdots, X(\mathbf{t}^{(r)}, \cdot))$$

as $N \to \infty$, for all $\mathbf{t}^{(1)}, \dots, \mathbf{t}^{(r)} \in \mathbb{R}^p$. (As usual, $\mathcal{L}(X_N(\mathbf{t}^{(1)}, \cdot), \dots)$ is the joint law of the variable $X_N(\mathbf{t}^{(1)}, \cdot), \dots, X_N(\mathbf{t}^{(r)}, \cdot)$ when \tilde{P} holds while $\mathcal{L}_H(X(\mathbf{t}^{(1)}, \cdot), \dots)$ is the law of the variables $X(\mathbf{t}^{(1)}, \cdot), \dots, X(\mathbf{t}^{(r)}, \cdot)$ when μ_H is the measure on Q_H .)

(ii) Tightness: For every $\epsilon > 0$, $\delta > 0$, there exists a compact subset K of Q_H such that,

(5.2)
$$\tilde{P}[\tilde{\omega}:X_N(\cdot,\tilde{\omega}) \varepsilon K^{\circ}] \ge 1 - \epsilon$$

for all N sufficiently large (m, n large), where

 $K^{\delta} = \{ f \varepsilon \mathfrak{M}(R^{p}) \colon ||f - g|| \leq \delta \text{ for some } g \varepsilon K \}.$

(We need not introduce outer probabilities because our processes X_N take on only a finite number of values.)

PROOF OF (i). If A is a Borel subset of R^{p} let us define

(5.3)
$$F_m(A, \tilde{\omega}) = \iint_A dF_m(\mathbf{t}, \tilde{\omega}).$$

Similarly, we shall talk about $G_n(A, \cdot)$, $H_N(A, \cdot)$ and H(A). Then, (5.1) will follow if we show that for any s disjoint rectangles A_1, \dots, A_s of \mathbb{R}^p such that $0 < H(A_i) < 1, 1 \leq i \leq s$, we have,

(5.4)
$$\mathfrak{L}(X_N)(A_1, \cdot), \cdots, X_N(A_s, \cdot)) \to \mathfrak{L}(X(A_1), \cdots, X(A_s))$$

as $m, n \to \infty$,

where $X(A_1), \dots, X(A_s)$ are Gaussian with mean 0 and Cov $(X(A_i), X(A_j)) = H(A_i \cap A_j) - H(A_i)H(A_j)$.

Let

(5.5)
$$Z_{i} = (nm^{-1})^{\frac{1}{2}} (m+n)^{-\frac{1}{2}}$$

if P_{iN} has been selected among the sample of m, = $-(mn^{-1})^{\frac{1}{2}}(m+n)^{-\frac{1}{2}}$, otherwise.

Then,

(5.6)
$$\tilde{P}[X_N(A_i, \cdot) = q_i, 1 \leq i \leq s]$$

= $\tilde{P}[\sum_{P_j \in A_i} Z_j = q_i, 1 \leq i \leq s].$

Let

(5.7)
$$c_{iN} = (nm^{-1})^{\frac{1}{2}}, \qquad 1 \leq i \leq m,$$

 $= -(mn^{-1})^{\frac{1}{2}}, \qquad m+1 \leq i \leq m+n.$

Let V_{1N} , \cdots , V_{NN} be obtained as in Lemma 4.2. Then,

$$P[\sum_{P_{jN} \in A_{i}} Z_{j} = q_{i}, 1 \leq i \leq s]$$

$$(5.8) = P[N^{-\frac{1}{2}} \sum_{i=1}^{NH_{N}(A_{i})} V_{iN} = q_{1}, \cdots,$$

$$N^{-\frac{1}{2}} \sum_{i=NH_{N}(\Sigma_{1}^{k}A_{j})+1}^{NH_{N}(\Sigma_{1}^{k}A_{j})} V_{iN} = q_{k}, \cdots,$$

$$N^{-\frac{1}{2}} \sum_{i=NH_{N}(\Sigma_{1}^{*-1}A_{j})+1}^{NH_{N}(\Sigma_{1}^{*-1}A_{j})+1} V_{iN} = q_{s}].$$

From (5.8)

(5.9)
$$\mathfrak{L}\left(\sum_{i=1}^{s} \lambda_{i} X_{N}(A_{i}, \cdot)\right) = \mathfrak{L}\left(\sum_{i=1}^{N} d_{iN} V_{iN}\right)$$

where $d_{iN} = N^{-\frac{1}{2}}(\lambda_{k} - \sum_{k=1}^{s} \lambda_{k} H_{N}(A_{k}))$ if $NH_{N}\left(\sum_{j=1}^{k-1} A_{j}\right)$

 $\leq i \leq NH_N(\sum_{j=1}^k A_j)$ for $1 \leq k \leq s$,

= 0 otherwise.

It readily follows that c_{iN} , d_{iN} as given satisfy (4.2) and (4.3). Then (5.4) is a consequence of (5.9) and Lemma 4.2.

PROOF OF (ii). We prove assertion (ii) in four steps. We shall show that the assertion holds,

(a) For H which are continuous and such that at least one marginal of H is the uniform distribution on [0, 1];

(b) For H which are continuous and such that at least one marginal of H is strictly increasing in [0, 1];

(c) For H which are continuous;

(d) For arbitrary H.

The next section deals with cases (a) and (b) while the final section has cases (c) and (d).

6. Proof of Theorem 3.1: Cases (iia) and (iib). We note first that in all cases it suffices to prove the result for sequences H_N such that all the P_{iN} fall in $I^p = \{t: ||t|| \leq 1\}$ where $||t|| = \max_{1 \leq i \leq p} |t_i|$. To see this let q be any 1 - 1 map of R onto (0, 1) which is a homeomorphism.

If $\{P_{iN}\}, \{H_N\}, H$ are given let $\{P_{iN}^*\}$ be the sequence of points in I^p obtained by transforming each coordinate of P_{iN} by q, and let H_N^* be the associated sequence of distribution functions. Clearly if $||H_N - H|| \to 0$, then $||H_N^* - H^*|| \to 0$ where H^* is the distribution function of a measure concentrated on I^p such that, if $q^{-1}(\mathbf{t}) = (q^{-1}(t_1), \cdots, q^{-1}(t_p))$,

(6.1)
$$H^*(\mathbf{t}) = H(q^{-1}(\mathbf{t}))$$

for t in the interior of I^{p} . It is easy to see now that if h is any functional regular for H, then there exists a functional h^{*} regular for H^{*} such that,

(6.2)
$$\mathfrak{L}(h(X_N(\cdot, \cdot))) = \mathfrak{L}(h^*(X_N^*(\cdot, \cdot)))$$

and

(6.3)
$$\mathfrak{L}_{H}(h(X(\cdot, \cdot))) = \mathfrak{L}_{H^{\bullet}}(h^{\bullet}(X(\cdot, \cdot))).$$

Our assertion follows and in the sequel we shall limit ourselves to sequences $\{P_{iN}\}$ falling in I^{p} . For cases (a), (b) and (c), we conclude from Proposition 2 of [6] that, in view of (i), to establish (ii) we need only show that,

(6.4) $\limsup_{\delta \to 0} \limsup_{N \to 0} P[\sup \{ |X_N(\mathbf{s}, \cdot) - X_N(\mathbf{t}, \cdot)| : ||\mathbf{s} - \mathbf{t}||$

$$\leq \delta, \mathbf{s}, \mathbf{t} \, \varepsilon \, I^p \geq \epsilon = 0$$
 for every $\epsilon > 0$.

PROOF OF (a). We need a lemma.

LEMMA 5.1 Let $A(\mathbf{s}, \delta) = \{\mathbf{t} \in I^p : \mathbf{t} \ge \mathbf{s}, ||\mathbf{t} - \mathbf{s}|| \le \delta\}$. Then (6.4) holds if

(6.5) $\lim_{\delta \to 0} \delta^{-p} \limsup_{\mathbf{s} \in I^p} \widetilde{P}[\sup \{ |X_N(\mathbf{s}, \cdot)| \}$

$$-X_{N}(\mathbf{t}, \cdot)|: \mathbf{t} \varepsilon A(\mathbf{s}, \delta)\} \geq \epsilon] = 0.$$

PROOF. Without loss of generality consider the particular sequence $\delta_m = m^{-1} \to 0$. Decompose I^p into the m^p cubes whose vertices are lattice points of the form $m^{-1}(j_1, \dots, j_p), 0 \leq j_i \leq m$. Then,

$$(6.6) \quad \widetilde{P}[\sup \{ |X_N(\mathbf{s}, \cdot) - X_N(\mathbf{t}, \cdot)| \colon ||\mathbf{s} - \mathbf{t}|| \leq \delta, \mathbf{s}, \mathbf{t} \in I^p] \geq \epsilon]$$

$$\leq \widetilde{P}[\mathbf{U}_{j=1}^{m^p}[\sup |X_N(\mathbf{e}_j, \cdot) - X_N(\mathbf{s}, \cdot)| \colon \mathbf{s} \in A(\mathbf{e}_j, \delta) \} \geq \epsilon/2]$$

$$\leq \sum_{j=1}^{m^p} \widetilde{P}[\sup \{ |X_N(\mathbf{e}_j, \cdot) - X_N(\mathbf{s}, \cdot)| \colon \mathbf{s} \in A(\mathbf{e}_j, \delta) \} \geq \epsilon/2]$$

where \mathbf{e}_1 , ..., \mathbf{e}_{m^p} are the "smallest" vertices of the m^p cubes we are considering.

Let $B(\mathbf{s}, \mathbf{t}) = {\mathbf{v}: \mathbf{s} < \mathbf{v} \leq \mathbf{t}}$ if $\mathbf{t} > \mathbf{s}$, where < corresponds to strict inequality for all coordinates. More generally, let T be the set of all sub collections of the natural numbers $1, \dots, p$. Denote such a collection by ${i_1, \dots, i_r}$. If $\mathbf{t} \geq \mathbf{s}$ let

$$B(\mathbf{s}, \mathbf{t}, \{i_1, \cdots, i_r\}) = \{\mathbf{v} \in I^p : \mathbf{v} = (v_1, \cdots, v_p),$$

 $\mathbf{v} \leq \mathbf{t}, v_k > s_k$ if $k \in \{i_1, \cdots, i_r\}, v_k \leq s_k$ otherwise}.

Thus $B(s, t) = B(s, t, \{1, \dots, p\})$. Then,

(6.7)
$$A(t) - A(s) = \sum_{\{i_1, \dots, i_r\} \in T} B(s, t, \{i_1, \dots, i_r\})$$

It follows that to prove (6.5) we need only establish

(6.8) $\lim_{\delta \to 0} \delta^{-p} \limsup_{m,n} \sup_{s \in I^p}$

 $\tilde{P}[\sup \{|X_N(B(\mathbf{s}, \mathbf{t}, (i_1, \cdots, i_r)), \cdot)| : \mathbf{t} \in A(\mathbf{s}, \delta)\} \geq \epsilon] = 0$

for every $\{i_1, \cdots, i_r\}, \epsilon > 0$.

We give the argument for $\{i_1, \dots, i_r\} = \{1, \dots, p\}$. The same proof holds for the general case.

We return to the notation of Section 4 but suppose throughout that c_{iN} is given by (5.7). We have,

(6.9)
$$\widetilde{P}[\sup\{|X_N(B(\mathbf{s},\mathbf{t}),\cdot)|:\mathbf{t} \in A(\mathbf{s},\delta)\} \ge \epsilon]$$

$$\leq P[\max N^{-\frac{1}{2}}\{|S_{i_1},\dots,i_p|:1 \le i_j \le N, 1 \le j \le p\} \ge \epsilon]$$

where $\delta = (\delta, \dots, \delta)$ and the subset S of Lemma 4.6 is defined as follows.

Move the $NH_N(B(\mathbf{s}, \mathbf{s} + \mathbf{\delta}))$ points in $B(\mathbf{s}, \mathbf{s} + \mathbf{\delta})$ whose position is given by H_N to points on the lattice of points of the form (i_1, \dots, i_p) with $1 \leq i_j \leq$ $NH_N(B(\mathbf{s}, \mathbf{s} + \mathbf{\delta})), 1 \leq j \leq p$, in such a way as to preserve the strict order relationships between the coordinates of these $NH_N(B(\mathbf{s}, \mathbf{s} + \mathbf{\delta}))$ points. It is permissible to move ties to distinct points. Let S be the set of all points (i_1, \dots, i_p) which are obtained in this fashion from the points $\{P_{iN}\}$. (The inequality in (5.18) is an equality if there are no duplicates among the points P_{iN} .)

Now, if $\limsup nN^{-1} \leq \delta$ it is easy to see that,

(6.10)
$$\liminf_{N} b(\gamma, N^{\frac{1}{2}} \epsilon) \geq 1 - \delta(1-\delta)^{-1} [(1-\delta(1-\delta)^{-1}-\gamma)\epsilon]^{-1}.$$

If H has a uniform marginal then $H(B(\mathbf{s}, \mathbf{s} + \mathbf{\delta})) \leq \mathbf{\delta}$ and by assumption,

(6.11) $\sup \{H_N(B(\mathbf{s}, \mathbf{s} + \mathbf{\delta})) : \mathbf{s} \in I^p\} \to \sup \{H(B(\mathbf{s}, \mathbf{s} + \mathbf{\delta})) : \mathbf{s} \in I^p\} \leq \delta.$

Applying Lemma 4.5 to the right hand side of (6.9) and using (6.10), (6.11) we obtain

$$\lim \sup_{N} \sup_{s \in I^{p}} \tilde{P}[\sup \{ |X_{N}(B(\mathbf{s}, \mathbf{t}), \cdot)| : \mathbf{t} \in A(\mathbf{s}, \delta) \} \geq \epsilon]$$

$$(6.12) \qquad \qquad \leq \gamma^{-r(p-1)} \epsilon^{-r} \limsup_{N \to 0} N^{-\frac{1}{2}} E|S_{n_{1}}|^{r}$$

$$\cdot \prod_{s=0}^{p-2} [1 - \delta(1 - \delta)^{-1} [(1 - \delta(1 - \delta)^{-1} - \gamma) \gamma^{s} \epsilon]^{2}]^{-1}$$

for suitable γ , where $n_1 = NH_N(B(\mathbf{s}, \mathbf{s} + \delta))$. But by Lemma 4.4 if (6.11) holds, (6.13) $\lim \sup_N N^{-k} E(S_{n_1})^{2k} \leq \delta^k \mu_{2k}.$

16

From (6.12) and (6.13) we conclude that in case (a),

(6.14)
$$\limsup_{\mathbf{s} \in I^p} P[\sup \{ |X_N(B(\mathbf{s}, \mathbf{t}), \cdot| : \mathbf{t} \in A(\mathbf{s}, \delta) \} \ge \epsilon]$$

$$\leq \epsilon^{-2k} 4^{k(p-1)} \delta^k \mu_{2k} (1 + o(\delta))$$

for δ so small that we can take $\gamma = \frac{1}{2}$. Taking k = p + 1 we see that (6.5) is satisfied. Theorem 3.1 is proved in case (a).

PROOF FOR CASE (b). Without loss of generality suppose $q(t) = H(t, \infty, \dots, \infty)$ is continuous and strictly increasing on [0, 1]. Let P'_{iN} be the point obtained by transforming the first co-ordinate of P_{iN} by q and leaving the others unchanged. If we now define H'_N , F'_m , G'_n , and X'_N suitably it is clear that,

(6.15)
$$||H_N' - H'|| \to 0$$

where $H'(t, \infty, \dots, \infty) = t$ for $0 \le t \le 1$. Therefore, from case (a), the process X_N' satisfies the tightness condition (6.4). But,

(6.16)
$$\sup \{ |X_N(\mathbf{s}, \cdot) - X_N(\mathbf{t}, \cdot)| : ||\mathbf{s} - \mathbf{t}|| \le \delta \}$$

= $\sup \{ |X_N'(\mathbf{q}^{-1}(\mathbf{s}), \cdot) - X_N'(\mathbf{q}^{-1}(\mathbf{t}), \cdot)| : ||\mathbf{s} - \mathbf{t}|| \le \delta \}$

where

$$\mathbf{q}^{-1}((t_1, \cdots, t_p)) = (q^{-1}(t_1), t_2, \cdots, t_p).$$

In view of (6.16) and the continuity of q^{-1} we see that the theorem holds in case (b).

7. Proof of Theorem 3.1: Cases (iic) and (iid). If q(t) is as above we say s is a point of increase of q if $t \neq s \Rightarrow q(t) \neq q(s)$. We say s is a point of increase if its first co-ordinate is a point of increase of q. An examination of the arguments used for case (b) will show that they serve to establish the theorem even if q is not strictly increasing as long as the P_{iN} are all points of increase. We need only note that the sup on the right hand side of (6.16) can be taken just over the set of all s, t whose first coordinates are points of increase of q.

If N is fixed, let $A^{(1)}$ be the set of all P_{iN} which are points of increase and $A^{(2)}$ be the set of all the other P_{iN} . Call the number of members of $A^{(i)}$, $N^{(i)}$, and if $N^{(i)} > 0$ let $H_N^{(i)}$ be the distribution function of the measure assigning mass $[N^{(i)}]^{-1}$ to each point of $A^{(i)}$, while if $N^{(i)} = 0$ let $H_N^{(i)} \equiv 0$. Let $M_N^{(i)}$ be the (random) number of points of the sample which fall in $A^{(i)}$, and if $M_N^{(i)} > 0$ let $F_N^{(i)}$ be the distribution function of the measure which assigns mass $[M_N^{(i)}]^{-1}$ to each point of the sample falling in $A^{(i)}$. Otherwise let $F_N^{(i)} = 0$. Define $G_N^{(i)}$ similarly and let,

(7.1)
$$X_N^{(1)}(\mathbf{t}, \cdot) = \{M_N^{(1)}(N^{(1)} - M_N^{(1)})^{-1}\}^{\frac{1}{2}} [N^{(1)}]^{\frac{1}{2}} [F_N^{(1)}(\mathbf{t}, \cdot) - H_N^{(1)}(\mathbf{t}, \cdot)] \text{ if } M_N^{(1)} \text{ and } N^{(1)} - M_N^{(1)} > \mathbf{0}$$

$$\equiv \mathbf{0} \quad \text{otherwise.}$$

Then,

$$X_{N}(\mathbf{t}, \cdot) = [M_{N}^{(1)}(N^{(1)} - M_{N}^{(1)})/mn]^{\frac{1}{2}} (N^{(1)}N^{-1})^{\frac{1}{2}} X_{N}^{(1)}(\mathbf{t}, \cdot) + (m/n)^{\frac{1}{2}} N^{\frac{1}{2}} (M_{N}^{(1)}m^{-1}) - (N^{(1)}N^{-1}) H_{N}^{(1)}(\mathbf{t}) + (m/n)^{\frac{1}{2}} N^{\frac{1}{2}} (M_{N}^{(2)}m^{-1}) F_{m}^{(2)}(\mathbf{t}, \cdot) - (N^{(2)}N^{-1}) H_{N}^{(2)}(\mathbf{t}, \cdot)).$$

It is easy to see that $(N^{(1)}N^{-1}) \to 1$ as $N \to \infty$. Furthermore since $M_N^{(1)}$ has a hypergeometric distribution we can readily calculate

(7.3)
$$E\{(mn^{-1})^{\frac{1}{2}} N^{\frac{1}{2}} (M_N^{(1)} m^{-1}) - (N^{(1)} N^{-1})\}^2$$

= $(N)(N-1)^{-1} - (N^{(1)} N^{(2)}) N^{-2} \to 0.$

From (7.2) we conclude that to establish (6.4) in case (c) we need only show, (7.4) $\lim_{s} \lim \sup_{N} \tilde{P}[\max \{ |X_{N}^{(1)}(\mathbf{t}, \cdot) \}$

 $-X_{N}^{(1)}(\mathbf{s}, \cdot)|: ||\mathbf{s} - \mathbf{t}|| \leq \delta, \, \mathbf{s}, \, \mathbf{t} \, \varepsilon \, I^{p} \} \geq \epsilon] = 0$

and

(7.5)
$$\lim_{N} \tilde{P}[\max\{(mn^{-1})^{\frac{1}{2}} N^{\frac{1}{2}} | (M_{N}^{(2)}m^{-1}) F_{m}^{(2)}(\mathbf{t}, \cdot) - N^{(2)}N^{-1} H_{N}^{(2)}(\mathbf{t}, \cdot) | \mathbf{t} \varepsilon I^{p}\} \ge \epsilon] = 0 \quad \text{for every} \quad \epsilon > 0.$$

Now note that given the sequence of variables $M_N^{(1)}$, the "distribution" of the stochastic processes $X_N^{(1)}(\cdot, \cdot)$ is the same as that of a sequence of processes $X_N(\cdot, \cdot)$ based on taking a sample of size $M_N^{(1)}$ from the $N^{(1)}$ points $\{P_{iN}\}$ lying in $A^{(1)}$. Since the H_N function corresponding to these points (which is just $H_N^{(1)}$) tends to H in norm and since all of the P_{iN} in $A^{(1)}$ are points of increase we may conclude from our extension of case (b) that for every $\epsilon' > 0$ there exists a $\delta' > 0$ and natural numbers m' and n' such that if $m' \leq r \leq N^{(1)}$, $N^{(1)} - r \geq n'$, $\delta \leq \delta'$ then,

(7.6)
$$\tilde{P}[\max\{|X_N^{(1)}(\mathbf{s},\,\cdot\,) - X_N^{(1)}(\mathbf{t},\,\cdot\,)|: \mathbf{s}, \mathbf{t} \in I^p, \\ ||\mathbf{s} - \mathbf{t}|| \le \delta\} \ge \epsilon |M_N^{(1)} = r] \le \epsilon'.$$

Since both $M_N^{(1)}$ and $N^{(1)} - M_N^{(1)} \to \infty$ in probability as $N \to \infty$ we see that (7.6) implies (7.4). Similarly we see that

$$(7.7) \quad \tilde{P}[\sup\{(mn^{-1})^{\frac{1}{2}} N^{\frac{1}{2}} | (M_N^{(2)}m^{-1}) F_m^{(2)}(\mathbf{t}, \cdot) \\ \quad - N^{(2)}N^{-1} H_N^{(2)}(\mathbf{t}, \cdot) | :\mathbf{t} \varepsilon I^p\} \geq \epsilon | M_N^{(2)} = r] \leq P[\sup\{|S_{(i_1, \cdots, i_p)}^{(2)} \\ \quad + n(i_1, \cdots, i_p) a(r, N) | :1 \leq i_j \leq N^{(2)}, j = 1, \cdots, p\} \geq \epsilon],$$

$$(7.8) \quad a(r, N) = N^{-\frac{1}{2}} [N^{(2)}]^{-1} [r(nm^{-1})^{\frac{1}{2}} - (N^{(2)} - r)(mn^{-1})^{\frac{1}{2}}]$$

and $\{S_{(i_1,\dots,i_p)}^{(2)}\}\$ is the collection of partial sums based on a configuration of $N^{(2)}$ points with selection from a finite population with $N^{(2)}$ members whose

18

values are given by,

(7.9)
$$c_{iN}^{(2)} = (nm^{-1})^{\frac{1}{2}} N^{-\frac{1}{2}} - a(r, N)$$
 for $i = 1, \dots, r$,
= $(mn^{-1})^{\frac{1}{2}} N^{-\frac{1}{2}} - a(r, N)$ for $i = r + 1, \dots, N^{(2)}$.

Then,

(7.10)
$$\sum_{i=1}^{N^{(2)}} c_{iN}^{(2)} = 0$$

and

(7.11)
$$\sum_{i=1}^{N^{(2)}} [c_{iN}^{(2)}]^2 \leq N^{-1} \{ (rnm^{-1}) + (N^{(2)} - r)(mn^{-1}) \}.$$

Let us now apply Lemma 4.7 to the configuration of $A^{(2)}$ with k = 3 and r = 2, if $N^{(2)} \ge 3$.

(7.12)
$$P[\sup\{|S_{(i_1,\dots,i_p)}|: 1 \le i_j \le N^{(2)}, j = 1, \dots, p\} \ge \epsilon] \le c(3,2)(1-\lambda(3)\sum_{i=1}^{N^{(2)}} [c_{iN}^{(2)}]^2)^{-(p-1)} E(S_i^{(2)})^2$$

where $S_l^{(2)}$ is the sum of a sample of *l* from the population $\{c_{lN}^{(2)}\}$. Of course,

(7.13)
$$E(S_l^{(2)})^2 \leq \frac{4}{9} \sum_{i=1}^{N(2)} [c_{iN}^{(2)}]^2.$$

Now, if
$$(rm^{-1}) \to 0$$
 and $(N^{(2)} - r)(n^{-1}) \to 0$ we see that,
(7.14) $\sum_{i=1}^{N^{(2)}} [c_{iN}^{(2)}]^2 \to 0.$

Hence, for $N^{(2)} \ge 3$, it follows that if

(7.15)
$$N^{(2)} a(M_N^{(2)}, N) \to 0,$$

$$(7.16) M_N^{(2)} m^{-1} \to 0$$

in probability, (7.5) will follow from (7.7) and (7.11)–(7.14). However, (7.15) and (7.16) follow by a calculation like (7.3). Of course, if $N^{(2)} \leq 3$, (7.5) follows trivially. The proof of case (c) is complete.

PROOF FOR CASE (d). For every $\{i_1, \dots, i_p\} \in T$ there exist at most a denumerable set of hyperplanes of the form $\{t: t_{ij} = a_{ij}\}$ which are assigned positive measure by H. Let us for simplicity label all of these hyperplanes by E_1, E_2, \dots . Let $E_0 = \mathbb{R}^p$. For every $i \ge 0$, let K_i be the union of all E_j properly included in E_i . Let N_i be the numbers of points P_{jN} which fall in $(E_i - K_i)$ and if $N_i > 0$ let H_{Ni} be the distribution function of the measure which assigns mass $[N_i]^{-1}$ to each point P_{jN} which falls in $(E_i - K_i)$. If $N_i = 0$ let $H_{Ni} \equiv 0$. Similarly, let M_{Ni} be the number of points of the sample falling in $(E_i - K_i)$ and define F_{mi} to be distribution functions of the measure assigning mass $[M_{Ni}]^{-1}$ to each such point if $M_{Ni} > 0$ and to be $\equiv 0$ otherwise. Finally let,

$$(7.17) \quad X_{Ni}(\mathbf{t}, \cdot) = \{M_{Ni}(N_i - M_{Ni})^{-1}\}^{\frac{1}{2}} N_i^{\frac{1}{2}}(F_{mi}(\mathbf{t}, \cdot) - H_{Ni}(\mathbf{t}, \cdot)).$$

We begin by proving case (d) when there are only a finite number say J of the E_i and $H(E_i - K_i) > 0$ for $i = 1, \dots, J$.

Write

(7.18)
$$X_N(\mathbf{t}, \cdot) = \sum_{i=1}^{J} \{ [M_{Ni}(N_i - M_{Ni})(mn)^{-1}]^{\frac{1}{2}} (N_i N^{-1})^{-\frac{1}{2}} X_{Ni}(\mathbf{t}, \cdot) + (mn^{-1})^{\frac{1}{2}} N^{\frac{1}{2}} (M_{Ni} m^{-1} - N_i N^{-1}) H_{Ni}(\mathbf{t}) \}.$$

Note first that,

uniformly in t εI^{p} where H_{i} is the distribution function of the measure assigned by H to $E_{i} - K_{i}$ divided by $H(E_{i} - K_{i})$. Further, by a calculation such as (7.3), since $N_{i}N^{-1} \rightarrow H(E_{i} - K_{i})$ we see that,

(7.20)
$$M_{Ni}m^{-1} \to H(E_i - K_i)$$

in probability and in fact,

(7.21)
$$(mn^{-1})^{\frac{1}{2}} N^{\frac{1}{2}} (M_{Ni}m^{-1} - N_i N^{-1}) = O_p(1),$$

(is bounded in probability).

Since the map $(f_1, \dots, f_J) \to \alpha_1 f_1 + \dots + \alpha_J f_J$ of $\mathfrak{M}(I^p) \times \dots \times \mathfrak{M}(I^p) \to M(I^p)$ is continuous in $|| \cdot ||$ we see that to prove (5.2) we need only show that, for every $\epsilon > 0, \delta > 0, 1 \leq i \leq J$, there exists a compact $K \subset Q_H$ such that,

(7.22)
$$\widetilde{P}[\widetilde{\omega}: X_{Ni}(\cdot, \widetilde{\omega}) \varepsilon K^{\delta}] \geq 1 - \epsilon,$$

and

(7.23)
$$\tilde{P}[\tilde{\omega}: (mn^{-1})^{\frac{1}{2}} N^{\frac{1}{2}} (M_{Ni}(\tilde{\omega})m^{-1} - N_i N^{-1}) H_{Ni}(\cdot) \varepsilon K^{\delta}] \ge 1 - \epsilon$$

for all m, n sufficiently large. Now, (7.19) and (7.21) suffice to establish (7.23).

For simplicity we shall prove (7.22) when $E_i = \{t: t_p = 0\}$. The proof of the general case is analogous. Define

(7.24)
$$\bar{X}_{Ni}(t_1, \cdots, t_{p-1}, \tilde{\omega}) = X_{Ni}(t_1, \cdots, t_{p-1}, 0, \tilde{\omega}).$$

It is easy to see that given $M_{Ni} = m_i$, the stochastic process $\tilde{X}_{Ni}(t_1, \dots, t_{p-1}, \cdot)$ is distributed as an $X_N(\mathbf{t}, \cdot)$ process based on a configuration of N_i points in \mathbb{R}^{p-1} with distribution function \tilde{H}_{Ni} given by

(7.25)
$$\bar{H}_{Ni}(t_1, \cdots, t_{p-1}) = H_{Ni}(t_1, \cdots, t_{p-1}, 0).$$

Of course,

$$(7.26) \quad \tilde{H}_{Ni}(t_1, \cdots, t_{p-1}) \to \tilde{H}_i(t_1, \cdots, t_{p-1}) = H_i(t_1, \cdots, t_{p-1}, 0)$$

and by construction \tilde{H}_i is continuous on I^{p-1} . Since $M_{Ni} \to \infty$ in probability as does $N_i - M_{Ni}$, and since case (c) has been established for all p we conclude that,

(7.27)
$$\lim_{\delta} \lim \sup_{m,n} \tilde{P}[\sup \{ |\tilde{X}_{Ni}(\mathbf{s}, \cdot) - \bar{X}_{Ni}(\mathbf{t}, \cdot)| \colon ||\mathbf{s} - \mathbf{t}|| \le \delta, \mathbf{s}, \mathbf{t} \in I^{p-1} \} \ge \epsilon] = 0$$

 $\mathbf{20}$

and hence by Dudley's Proposition 2 [6], for every $\epsilon > 0$, $\epsilon > 0$ there exists a compact subset K_i of $Q\tilde{H}_i$ such that

(7.28)
$$\widetilde{P}[\widetilde{\omega}:\widetilde{X}_{Ni}(\cdot,\,\widetilde{\omega})\,\varepsilon\,K_i^{\,\delta}] \ge \epsilon$$

for all *m*, *n* sufficiently large. (Note that $K_i^{\delta} \subset \mathfrak{M}(I^{p-1})$.) Now, the map *T* of $\mathfrak{M}(I^p)$ into $\mathfrak{M}(I^{p-1})$ which is defined by,

(7.29)
$$Tf(t_1, \cdots, t_{p-1}) = f(t_1, \cdots, t_{p-1}, 0)$$

is continuous in $\|\cdot\|$ and hence (7.28) implies (7.22).

To complete the proof of case (d) we must finally deal with the possibility that some of the $H(E_i - K_i) = 0$ and that $J = \infty$. Let i_1, i_2, \cdots be the sequence of indices such that $H(E_{i_j} - K_{i_j}) > 0$. Write,

(7.30)
$$X_{N}(\mathbf{t}, \cdot) = \sum_{j=1}^{A} X_{Nij}(\mathbf{t}, \cdot) \left[(M_{Nij}(N_{ij} - M_{Nij}) (mn)^{-1} \right]^{\frac{1}{2}} (N_{ij}N^{-1})^{-\frac{1}{2}} + R_{NA}(\mathbf{t}, \cdot).$$

From our previous assertions it is evident that (5.2) will follow in this general case if we can show that for every $\epsilon > 0$, $\delta > 0$, there exists an $A < \infty$ such that

(7.31)
$$P[\sup\{|R_{NA}(\mathbf{t}, \cdot)|: \mathbf{t} \in I^p\} \leq \delta] \geq 1 - \epsilon$$

for all m, n sufficiently large.

Let,

(7.32)
$$g(A) = \sum_{i=A+1}^{\infty} H(E_{ij} - K_{ij})$$

and $Ng_N(A)$ equal the number of P_{iN} which are not in $\bigcup_{j=1}^{A} (E_{ij} - K_{ij})$. Evidently,

(7.33) $g_N(A) \to g(A)$ as $N \to \infty$ and,

(7.34)
$$g(A) \to 0$$
 as $A \to \infty$.

We argue as in the proof of (7.5). Let M_N^A be the number of points in the sample which are not in $\bigcup_{j=1}^{A} (E_{ij} - K_{ij})$. Then,

$$(7.35) \quad \tilde{P}[\sup\{|R_{NA}(\mathbf{t}, \cdot)|: \mathbf{t} \in I^{p} \ \delta \mid M_{N}^{A} = r] \\ \leq P[\max\{|S_{(i_{1}, \cdots, i_{p})} + n(i_{1}, \cdots, i_{p})a(r, N, A)|: 1 \leq i_{j} \leq Ng_{N}(A), \\ j = 1, \cdots, p\} \geq \delta]$$

where $S_{(i_1,\ldots,i_p)}^A$ are partial sums corresponding to a configuration of $Ng_N(A)$ points and sampling from a population of $Ng_N(A)$ numbers $\{c_{iN}(r)\}$ which are given by,

(7.36)
$$c_{iN}^{A}(r) = (nm^{-1})^{\frac{1}{2}}N^{\frac{1}{2}} - a(r, N, A)$$
 for $i = 1, \dots, r$
$$= -(mn^{-1})^{\frac{1}{2}}N^{\frac{1}{2}} - a(r, N, A)$$
 for $i = r + 1, \dots, Ng_{N}(A)$

and

$$(7.37) \quad a(r, N, A) = N^{-\frac{1}{2}} (Ng_N(A))^{-1} \{r(nm^{-1})^{\frac{1}{2}} - [Ng_N(A) - r](mn^{-1})^{\frac{1}{2}} \}.$$

If $Ng_N(A) \ge 3$ we can apply Lemma 4.7 with k = 3 and r = 2 to the right hand side of (7.35) and get after some computation,

$$\begin{split} \tilde{P}[\sup\{|R_{N}(\mathbf{t},\,\cdot\,)|:\mathbf{t}\,\varepsilon\,I^{p}\} &\geq \delta \mid M_{N}^{A}]\\ (7.38) &\leq \frac{4}{9}c(2,\,3)[1\,-\,\lambda(3)(\delta + Ng_{N}(A\,)|a(M_{N}^{A},\,N,\,A\,)|)^{-2}K(M_{N}^{A},\,A\,)]^{-(p-1)}\\ &\cdot (\delta + Ng_{N}(A\,)|a(M_{N}^{A},\,N,\,A\,)|)^{-2}K(M_{N}^{A},\,A\,), \end{split}$$

where

(7.40)
$$K(r A), = \sum_{i=1}^{NG} (c_{iN}(r))^2.$$

Finally, we have,

(7.41)
$$K(M_N^{A}, A) \leq \max\{M_N^{A}m^{-1}, (Ng_N(A) - M_N^{A})n^{-1}\}$$

and

$$(7.42) \quad E([Ng_N(A)a(M_N^A, N, A)]^2) = E((mn^{-1})^{\frac{1}{2}}N^{\frac{1}{2}}(M_N^Am^{-1} - g_N(A))^2) = g_N(A)(1 - g_N(A))N(N - 1)^{-1} \le 2 g_N(A)$$

for N sufficiently large. From (7.42) we see that,

(7.43)
$$\tilde{P}[|[Ng_N(A)]^{-1}a(M_N^A, N, A)| \ge g_N^{\frac{1}{2}}(A)] \le 2g_N^{\frac{1}{2}}(A),$$

(7.44)
$$\tilde{P}[M_N{}^A m^{-1} \ge g_N(A) + m^{-1}] \le 2g_N(A)$$

and

(7.45)
$$\tilde{P}[(Ng_N(A) - M_N^A)m^{-1} \ge g_N(A) + n^{-\frac{1}{2}}] \le 2g_N(A).$$

By (7.39), (7.41) and (7.43)-(7.45) we obtain,

$$\begin{aligned} \tilde{P}[\sup\{|R_{NA}(\mathbf{t},\,\cdot\,)|:\mathbf{t}\,\varepsilon\,I^{p}\} &\geq \delta] \\ (7.46) &\leq \frac{4}{9}c(2,\,3)[1\,-\,\lambda(3)(\delta\,+\,g_{N}^{\frac{1}{4}}(A\,))^{-2}\,g_{N}(A\,) \\ &+\,\max(m^{-\frac{1}{2}},\,n^{-\frac{1}{2}})]^{-(p-1)}[(g_{N}(A\,) \\ &+\,\max(m^{-\frac{1}{2}},\,n^{-\frac{1}{2}})](\delta\,+\,g_{N}^{\frac{1}{4}}(A\,))^{-2}\,+\,6g_{N}^{\frac{1}{2}}(A\,), \end{aligned}$$

for $Ng_N(A) \ge 3$. Letting $N \to \infty$ and then $A \to \infty$ we have by (7.33) and (7.34),

(7.47)
$$\lim_{A\to\infty} \limsup_{N} \widetilde{P}[\sup\{|R_{NA}(t, \cdot)|: t \in I^p\} \ge \delta] = 0$$

and this is equivalent to requiring (7.31) for N sufficiently large. The same conclusion evidently holds if $Ng_N(A) < 3$. The theorem is proved.

Acknowledgment. I am grateful to Volker Strassen for a helpful conversation.

REFERENCES

 BARTON, C. L. and MALLOWS, D. E. (1965). Some aspects of the random sequence. Ann. Math. Statist. 36 236-260.

 $\mathbf{22}$

- [2] DAVID, F. N. and FIX, E. (1960). Rank correlation and regression in a nonnormal surface. Proc. Fourth Berkeley Symp. Math. Statist. Prob. 6 177-197. Univ. of California Press.
- [3] DARLING, D. A. (1957). The Kolmogorov-Smirnov, Crámer-von Mises tests. Ann. Math. Statist. 28 823-838.
- [4] DOOB, J. L. (1949). Heuristic approach to the Kolmogorov-Smirnov theorems. Ann. Math. Statist. 20 393-403.
- [5] DUDLEY, R. M. (1965). Gaussian processes on several parameters. Ann. Math. Statist. 36 771-788.
- [6] DUDLEY, R. M. (1966). Weak convergence of probabilities on nonseparable metric spaces and empirical measures on Euclidean spaces. *Illinois J. Math.* 10 109-126.
- [7] FRASER, D. A. S. (1957). Nonparametric Methods in Statistics. Wiley, New York.
- [8] HAJEK, J. (1961). Some extensions of the Wald-Wolfowitz-Noether theorems. Ann. Math. Statist. 32 506-523.
- [9] HAJEK, J. (1965). Extension of the Kolmogorov-Smirnov test to regression alternatives. Bernoulli-Bayes-Laplace Seminar. Springer Verlag, Heidelberg.
- [10] KIEFER, J. (1961). On large deviations of the empiric d.f. of vector chance variables and a law of the iterated logarithm. Pacific J. Math. 11 649-660.
- [11] KIEFER, J. and WOLFOWITZ, J. (1958). On the deviations of the empiric distribution function of vector chance variables. *Trans. Amer. Math. Soc.* 87 173-186.
- [12] LEHMANN, E. L. (1951). Consistency and unbiasedness of certain nonparametric tests. Ann. Math. Statist. 22 165-179.
- [13] Рконовоv, Y. (1956). Convergence of random processes and limit theorems in probability theory. Theor. Prob. Appl. 1 157-214.
- [14] ROSENBLATT, M. (1952). Remarks on a multivariate transformation. Ann. Math. Statist.
 23 470-472
- [15] SCHMID, P. (1958). On the Kolmogorov Smirnov limit theorem for discontinuous distribution functions. Ann. Math. Statist. 29 1011-1027.
- [16] SEN, P. K. and CHATTERJEE, S. K. (1964). Nonparametric tests for the bivariate two sample location problem. *Calcutta Statist. Assn. Bull.* 13 18-58.
- [17] SIMPSON, P. B. (1951). Note on estimation of a bivariate distribution function. Ann. Math. Statist. 22 476-478.
- [18] VINIZE, I. (1960). Two sample tests based on order statistics. Proc. Fourth Berkeley Symp. Math. Statist. Prob. 1 697-705. Univ. of California Press.
- [19] WALD, A. and WOLFOWITZ, J. (1944). Statistical tests based on permutations of the observations. Ann. Math. Statist. 15 358-372.
- [20] WEISS, L. (1960). Two sample tests for multivariate distributions. Ann. Math. Statist. **31** 159-164.