A NONPARAMETRIC TEST FOR EQUALITY OF DISTRIBUTIONS WITH
MIXED CATEGORICAL AND CONTINUOUS DATA

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Abstract. In this paper we consider the problem of testing for equality of two density or two
conditional density functions defined over mixed discrete and continuous variables. We smooth both
the discrete and continuous variables, with the smoothing parameters chosen via least-squares cross-
validation. The test statistics are shown to have (asymptotic) normal null distributions. However,
we advocate the use of bootstrap methods in order to better approximate their null distribution
in finite-sample settings. Simulations show that the proposed tests have better power than both
conventional frequency-based tests and smoothing tests based on ad hoc smoothing parameter
selection, while a demonstrative empirical application to the joint distribution of earnings and
educational attainment underscores the utility of the proposed approach in mixed data settings.

Keywords: Mixed discrete and continuous variables; Density testing; Nonparametric smoothing;
Cross-validation.

Date: March 8, 2007.

We would like to thank three anonymous referees, the associate editor, and Peter Robinson for their numerous helpful
comments that collectively led to a much improved version of this paper. Li’s research is partially supported by the
Private Enterprise Research Center, Texas A&M University. Racine would like to gratefully acknowledge support
from Natural Sciences and Engineering Research Council of Canada (NSERC:www.nserc.ca), the Social Sciences and
Humanities Research Council of Canada (SSHRC:www.sshrc.ca), and the Shared Hierarchical Academic Research
Computing Network (SHARCNET:www.sharcnet.ca). We would also like to acknowledge Jose Galdo for his assistance
with data management.

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1. Introduction

It is difficult to think of a more ubiquitous test in applied statistics than the test for equality of distributions and conditional distributions, sometimes conditioned on discrete covariates. The most popular variants are derivative since they involve testing equality of moments, such as means and/or variances, or perhaps quantiles. Examples include tests for ‘regime change’, heteroscedasticity, and ‘symmetry’. Also, comparing distributions, or reconstructing indirectly observed distributions (such as the counterfactuals in program evaluation) is implicit and ever present in almost all statistical/econometric work. However, moment-based tests, which only compare a finite number of moments from two distributions, are not consistent tests. The same can be said for parametric tests which require specification of the null distribution. When the parametric null distribution is misspecified, parametric tests can lead to erroneous conclusions. Generally, interest truly lies in detecting any potential difference between two distributions without having to specify a parametric family, not just their means or variances. When this is the case, nonparametric tests have obvious appeal.

A number of kernel-based tests of equality of distribution functions exist; however, existing kernel-based tests presume that the underlying variable is continuous in nature; see Ahmad & van Belle (1974), Mammen (1992), Fan & Gencay (1993), Li (1996), Fan & Ullah (1999), and the references therein. It is widely known that a traditional ‘frequency-based’ kernel approach could be used to consistently estimate a joint probability function in the presence of mixed continuous and categorical variables, and hence one could readily construct a kernel-based test for the equality between two unknown density functions by simply employing the conventional frequency kernel method. In contrast we consider kernel ‘smoothing’ the discrete variables as well, following a rich literature in statistics on smoothing discrete variables and its potential benefits; see Aitchison & Aitken (1976), Hall (1981), Grund & Hall (1993), Scott (1992), Simonoff (1996), Li & Racine (2003), Hall, Racine & Li (2004), Hall, Li & Racine (forthcoming), among others. Though smoothing discrete variables in an appropriate manner may introduce some finite-sample bias, it simultaneously reduces finite-sample variance substantially, and leads to a reduction in the finite-sample mean square error of the nonparametric estimator relative to the frequency-based estimator. It turns out that, for
testing purposes, this is also highly desirable. The tests developed herein are extensions of existing frequency-based ‘smooth’ kernel tests. ‘Non-smooth’ (i.e., empirical cumulative distribution function (CDF)) tests of distributional differences have recently been examined and reviewed in Anderson (2001).

In this paper we first propose a kernel-based test for equality of distributions mounted on a square integral metric defined over mixed continuous/discrete variables. We then extend our result to the case of testing the equality of two conditional distributions. Conditional distributions, such as that of earnings given gender, age or education categories, are often the main targets of inference and policy analysis. As an alternative to our approach in this paper, entropy metrics have been used for testing equality of distributions, or hypotheses which may be cast as such. For a pioneering paper see Robinson (1991), as well as Hong & White (2000), Ahmad & Li (1997), and Racine & Maasoumi (forthcoming). We use data-driven bandwidth selection methods, smooth both the continuous and discrete variables in a particular manner, and advocate a resampling method for obtaining the statistic’s null distribution, though we also provide its limiting (asymptotic) null distribution and prove that the bootstrap works. It is well known that the selection of smoothing parameters is of crucial importance in nonparametric estimation, and it is now known that the selection of smoothing parameters also affects the size and power of nonparametric tests such as ours. When discrete variables are present, cross-validation has been shown to be an effective method of smoothing parameter selection. Not only is there a large sample optimality property associated with minimizing estimation mean square error, but we also avoid sample splitting in small sample applications. When one smooths both the discrete and continuous variables, cross-validation seems to be the only feasible way of selecting the smoothing parameters. Configuring plug-in rules for mixed data is an algebraically tedious task, and no general formulae are yet available. Additionally, plug-in rules, even after adaption to mixed data, require choice of ‘pilot’ smoothing parameters, and it is not clear how to best make that selection when both continuous and discrete variables are involved.

The paper is organized as follows. Section 2 presents a test for the equality of two unconditional distribution functions and examines the asymptotic distribution of the test statistic, Section
3 proposes a test for the equality of two conditional density functions, Section 4 presents two simulation experiments designed to assess the finite-sample performance of the estimator, while Section 5 presents a demonstrative empirical application that tests for differences in the joint distribution of earnings and educational attainment over time. Section 6 concludes, and all proofs are relegated to the appendix.

2. A Nonparametric Test for the Equality of Unconditional Density Functions with Mixed Categorical and Continuous Data

2.1. Testing the Equality of Two Density Functions. We consider the case where we are faced with a mixture of discrete and continuous data. Let \( X = (X_c, X^d) \in \mathbb{R}^q \times \mathbb{S}^r \), where \( X_c \) is the continuous variable having dimension \( q \), and \( X^d \) is the discrete variable having dimension \( r \) and assuming values in \( S^r = \prod_{s=1}^r \{0, 1, \ldots, c_s - 1\} \). Similarly, \( Y = (Y_c, Y^d) \), which has the same dimension as \( X \). Let \( f(\cdot) \) and \( g(\cdot) \) denote the density functions of \( X \) and \( Y \), respectively, and let \( \{X_i\}_{i=1}^{n_1} \) and \( \{Y_i\}_{i=1}^{n_2} \) be i.i.d. random draws from populations having density functions \( f(\cdot) \) and \( g(\cdot) \), respectively.\(^1\) We are interested in testing the null hypothesis that

\[
H_0 : f(x) = g(x) \text{ for almost all } x \in \mathbb{R}^q \times \mathbb{S}^r
\]

against the alternative hypothesis \( H_1 \) that \( f(x) \neq g(x) \) on a set with positive measure. We first discuss how to estimate \( f(\cdot) \) and \( g(\cdot) \) and then outline the test statistic.

Let \( x^d_s \) and \( X^d_{is} \) denote the \( s \)th components of \( x^d \) and \( X^d_i \) respectively. Following Aitchison & Aitken (1976), for \( x_s, X^d_{is} \in \mathbb{S}^r = \{0, 1, \ldots, c_s - 1\} \) (\( x^d_s \) takes \( c_s \) different values), we define a univariate kernel function

\[
\displaystyle l(X^d_{is}, x^d_s, \lambda_s) = \begin{cases} 
1 - \lambda_s & \text{if } X^d_{is} = x^d_s, \\
\lambda_s / (c_s - 1) & \text{if } X^d_{is} \neq x^d_s,
\end{cases}
\]

\(^1\)For what follows, when we consider a distribution defined over mixed continuous and discrete variables, we shall use the word 'density' to mean that, for any given value of \( x^d \in S^r \), \( f(x^c, x^d) \) is absolutely continuous with respect to \( x^c \).
where the range of the smoothing parameter \( \lambda_s \) is \([0, (c_s - 1)/c_s]\). Note that when \( \lambda_s = 0 \),
\( l(X_{is}^d, x_s^d, 0) = I(X_{is}^d = x_s^d) \) becomes an indicator function. We shall use \( I(\cdot) \) to denote an indicator function, i.e., \( I(A) = 1 \) if the event \( A \) holds true, otherwise \( I(A) = 0 \). Observe that if \( \lambda_s = (c_s - 1)/c_s \), then \( l(X_{is}^d, x_s^d, \frac{c_s - 1}{c_s}) = \frac{1}{c_s} \) which is a constant for all values of \( X_{is}^d \) and \( x_s^d \).

A product kernel function for the discrete variable components \( x^d \) is given by

\[
(2.2) \quad L_{\lambda, x_i, x} = \prod_{t=1}^{r} l(X_{is}^d, x_s^d, \lambda_s) = \prod_{s=1}^{r} \left\{ \lambda_s/(c_s - 1) \right\}^{I_{x_{is}^d \neq x_s^d}} (1 - \lambda_s)^{I_{x_{is}^d = x_s^d}},
\]

where \( I_{x_{is}^d \neq x_s^d} = I(X_{is}^d \neq x_s^d) \), and \( I_{x_{is}^d = x_s^d} = I(X_{is}^d = x_s^d) \).

Let \( w \left( \frac{x_c^i - X_{is}^c}{h_s} \right) \) be a univariate kernel function associated with the continuous variable \( x_c^i \), where \( h_s \) is the associated smoothing parameter. The product kernel for the continuous variable components \( x_c^i \) is given by

\[
(2.3) \quad W_{h, x_i, x} = \prod_{s=1}^{q} \frac{1}{h_s} w \left( \frac{X_{is}^c - x_s^c}{h_s} \right).
\]

The ‘generalized’ product kernel defined over all components, discrete and continuous, is given by

\[
(2.4) \quad K_{\gamma, x_i, x} = W_{h, x_i, x} L_{\lambda, x_i, x},
\]

where \( \gamma = (h, \lambda) \), and where \( L_{\lambda, x_i, x} \) and \( W_{h, x_i, x} \) are defined in (2.2) and (2.3), respectively. Note that we define \( \prod_{i=1}^{n1} = 1 \) so that the generalized product kernel is well defined even if \( q \) or \( r \) equal zero.

We estimate the joint density of \( f(x) \) by

\[
(2.5) \quad \hat{f}(x) = \frac{1}{n1} \sum_{i=1}^{n1} K_{\gamma, x_i, x},
\]

Similarly, we estimate the joint density of \( g(x) \) by

\[
(2.6) \quad \hat{g}(x) = \frac{1}{n2} \sum_{i=1}^{n2} K_{\gamma, y_i, x},
\]

where \( K_{\gamma, y_i, x} = W_{h, y_i, x} L_{\lambda, y_i, x} \).
A test statistic can be constructed based on the integrated squared density difference given by

\[ I = \int [f(x) - g(x)]^2 \, dx = \int [f(x)dF(x) + g(x)dG(x) - f(x)dG(x) - g(x)dF(x)], \]

where \( F() \) and \( G() \) are the cumulative distribution functions for \( X \) and \( Y \), respectively, and where \( \int dx = \sum_{x^d \in \Xi_d} \int dx^c \).

Replacing \( f() \) and \( g() \) by their kernel estimates, and replacing \( F() \) and \( G() \) by their empirical distribution functions, we obtain the following test statistic,

\[
I^a_n = \frac{1}{n_1} \sum_{i=1}^{n_1} \hat{f}(X_i) + \frac{1}{n_2} \sum_{i=1}^{n_2} \hat{g}(Y_i) - \frac{1}{n_1} \sum_{i=1}^{n_1} \hat{g}(X_i) - \frac{1}{n_2} \sum_{i=1}^{n_2} \hat{f}(Y_i)
\]

\[= \frac{1}{n_2} \sum_{i=1}^{n_1} \sum_{j=1}^{n_1} K_{\gamma,x_i,x_j} + \frac{1}{n_2} \sum_{i=1}^{n_2} \sum_{j=1}^{n_2} K_{\gamma,y_i,y_j} - \frac{1}{n_1 n_2} \left[ \sum_{i=1}^{n_1} \sum_{j=1}^{n_1} K_{\gamma,x_i,y_j} + \sum_{i=1}^{n_2} \sum_{j=1}^{n_2} K_{\gamma,x_j,y_i} \right]. \tag{2.7}
\]

It can be shown that the test statistic \( I^a_n \) has a non-zero center term, say \( c_n \), even under the null hypothesis. Therefore, in practice one needs to estimate the center term, say by \( \hat{c}_n \), and subtract \( \hat{c}_n \) from \( I^a_n \) to obtain an asymptotic zero mean test statistic (under \( H_0 \)). In order to remove the non-zero center term when testing the equality of two density functions with only continuous variables \( x^c \), Li (1996) proposed a center-free test statistic which is obtained by removing the \( i = j \) terms in the double summations appearing in \( I^a_n \), i.e.,

\[
I^b_n = \frac{1}{n_1(n_1 - 1)} \sum_{i=1}^{n_1} \sum_{j \neq i} K_{\gamma,x_i,x_j} + \frac{1}{n_2(n_2 - 1)} \sum_{i=1}^{n_2} \sum_{j \neq i} K_{\gamma,y_i,y_j} \\
- \frac{1}{n_1(n_2 - 1)} \sum_{i=1}^{n_1} \sum_{j \neq i} K_{\gamma,x_i,y_j} - \frac{1}{(n_1 - 1)n_2} \sum_{i=1}^{n_1} \sum_{j \neq i} K_{\gamma,x_j,y_i}. \tag{2.8}
\]

It can be shown that \( I^b_n \) is an asymptotic zero mean test statistic (under \( H_0 \)). However, the \( I^b_n \) test has a new problem in that it depends on the ordering of the data. To see this, note that the third term of \( I^b_n \) differs from the third term of \( I^a_n \) by \( \sum_{i=1}^{\min\{n_1, n_2\}} K_{\gamma,x_i,y_j} \), which depends on how one orders the data \( \{X_i\}_{i=1}^{n_1} \) and \( \{Y_i\}_{i=1}^{n_2} \). Below we propose a test statistic which does not have a non-zero center term (under \( H_0 \)) and is also invariant to the ordering of the data. The test statistic
we propose is given by

\[ I_n = \frac{1}{n_1(n_1 - 1)} \sum_{i=1}^{n_1} \sum_{j \neq i}^{n_1} K_{\gamma,x_i,x_j} + \frac{1}{n_2(n_2 - 1)} \sum_{i=1}^{n_2} \sum_{j \neq i}^{n_2} K_{\gamma,y_i,y_j} \]

\[ - \frac{1}{n_1 n_2} \left[ \sum_{i=1}^{n_1} \sum_{j=1}^{n_2} K_{\gamma,x_i,y_j} + \sum_{i=1}^{n_2} \sum_{j=1}^{n_1} K_{\gamma,y_i,x_j} \right] \]

(2.9)

Note that the double summation of the first two terms of \( I_n \) removes the \( i = j \) terms, while the third term in \( I_n \) does not remove the \( i = j \) terms. Therefore, the \( I_n \) test statistic differs from both \( I_n^a \) and \( I_n^b \). Clearly, the test statistic \( I_n \) is invariant to the ordering of the data because the terms removed from \( I_n^a \) (i.e., \( \sum_{i=1}^{n_1} K_{\gamma,x_i,x_j} \) and \( \sum_{i=1}^{n_2} K_{\gamma,y_i,y_j} \)) are both invariant with respect to the ordering of the data. We will show in Theorem 2.1 below that \( I_n \) is an asymptotic zero mean test statistic (under \( H_0 \)).

The following conditions will be used to derive the asymptotic distribution of \( I_n \).

(C1) The data \( \{X_i\}_{i=1}^{n_1} \) and \( \{Y_i\}_{i=1}^{n_2} \) are independent and identically distributed (i.i.d.) as \( X \) and \( Y \) respectively.

(C2) For all \( x^d \in S^r \), both \( f(\cdot, x^d) \) and \( g(\cdot, x^d) \) are bounded (from above by some positive constants) and continuous functions (continuous with respect to \( x^c \)). The kernel function \( w(\cdot) \) is a bounded, non-negative second order kernel, \( \int w(v)v^4dv \) is finite, and it satisfies a Lipschitz condition: \( |w(u) - w(v)| \leq \xi(v)|u - v| \), where \( \xi(\cdot) \) is a bounded smooth function with \( \int \xi(v)v^4dv < \infty \).

(C3) Letting \( \delta_n = n_1/n_2 \), then as \( n = \min\{n_1, n_2\} \to \infty \), \( \delta_n \to \delta \in (0,1) \), \( n h_1 \ldots h_q \to \infty \), \( h_s \to 0 \) for \( s = 1, \ldots, q \) and \( \lambda_s \to 0 \) for \( s = 1, \ldots, r \).

Note that in (C1) we assume that \( X_i \) (\( Y_i \)) is independent of \( X_j \) (\( Y_j \)) for \( j \neq i \). When \( n_1 = n_2 = n \), however, we allow for the possibility that \( X_i \) and \( Y_i \) are correlated, as would be the case in panel or longitudinal settings where we have repeated measures on individuals. The i.i.d. assumption can be relaxed to weakly dependent (\( \beta \)-mixing) data processes, in which case one needs to apply the central limit theorem for degenerate U-statistics with weakly dependent data as given in Fan & Li (1999) in order to derive the asymptotic distribution of the test statistic. Of course, with dependent data, the bootstrap procedure (see Theorem 2.3 below) will also need to be modified;
block or stationary bootstrapping or subsampling methods would be appropriate. In the remaining part of this paper, we will only consider i.i.d. data as stated in (C1).

The other conditions under which Theorem 2.1 holds are quite weak. (C2) only requires that \( f(\cdot) \) and \( g(\cdot) \) are bounded and continuous, and (C3) is the minimum condition placed upon the smoothing parameters required for consistent estimation of \( f(\cdot) \) and \( g(\cdot) \). In addition, (C3) requires that the two sample sizes have the same order of magnitude.

The following theorem provides the asymptotic null distribution of the test statistic \( I_n \).

**Theorem 2.1.** Assuming that conditions (C1) through (C3) hold, we have, under \( H_0 \), that

\[
T_n = (n_1 n_2 h_1 \ldots h_q)^{1/2} \frac{I_n}{\sigma_n} \rightarrow N(0, 1) \text{ in distribution,}
\]

where

\[
\sigma_n^2 = 2(n_1 n_2 h_1 \ldots h_q) \left[ \frac{1}{n_1^2(n_1 - 1)^2} \sum_{i=1}^{n_1} \sum_{j \neq i} (K_{\gamma, x_i, x_j})^2 + \frac{1}{n_2^2(n_2 - 1)^2} \sum_{i=1}^{n_2} \sum_{j \neq i} (K_{\gamma, y_i, y_j})^2 \right]
\]

\[
+ \frac{1}{n_1 n_2^2} \sum_{i=1}^{n_1} \sum_{j=1}^{n_2} (K_{\gamma, x_i, y_j})^2 + \frac{1}{n_1^2 n_2^2} \sum_{i=1}^{n_1} \sum_{j=1}^{n_2} (K_{\gamma, y_j, x_i})^2 \right],
\]

which is a consistent estimator of \( \sigma_0^2 = 2[\delta^{-1} + \delta + 2][E[f(X_i)][f W^2(v)]dv] \), the asymptotic variance of \( (n_1 n_2 h_1 \ldots h_q)^{1/2} I_n \), where \( \delta = \lim_{\min(n_1, n_2) \rightarrow \infty} (n_1/n_2) \).

The proof of Theorem 2.1 is given in the appendix.

It can also be shown that, when \( H_0 \) is false, the test statistic \( T_n \) will diverge to \( +\infty \) at the rate of \( (n_1 n_2 h_1 \ldots h_q)^{1/2} \). To see this, note that when \( H_0 \) is false, one can show that \( I_n \rightarrow \int [f(x) - g(x)]^2 dx \equiv C > 0 \) (in probability), where \( \sigma_n = O_p(1) \). Hence, \( T_n \) will diverge to \( +\infty \) at the rate of \( (n_1 n_2 h_1 \ldots h_q)^{1/2} \), and therefore it is a consistent test.

It is well known that the selection of smoothing parameters is of crucial importance in nonparametric estimation, and it is now known that the selection of smoothing parameters also affects the performance (particularly the power) of nonparametric tests such as the \( I_n \) test. Given the reasons outlined in the introduction as to why cross-validation methods seem to be the only feasible way of selecting the smoothing parameters in the presence of mixed discrete and continuous variables, we suggest using cross-validations method for selecting \((h, \lambda)\).
The cross-validation method we consider involves selecting smoothing parameters by minimizing a sample analogue of the integrated square error (ISE) of the density estimator. The ISE is defined by
\[
ISE = \int \left[ \hat{f}(x) - f(x) \right]^2 dx = \int \hat{f}(x)^2 dx - 2 \int \hat{f}(x)f(x) + \int f(x)^2,
\]
where \( \int dx = \sum_{x^c} \int dx^c \). The third term on the right-hand side of ISE does not depend on the smoothing parameters. Therefore, in practice one chooses the smoothing parameters to minimize an estimator of \( \int \hat{f}(x)^2 dx - 2 \int \hat{f}(x)f(x) \). Let \( \{Z_i\}_{i=1}^N \) denote the pooled sample \( (N = n_1 + n_2) \), i.e., \( Z_i = X_i \) for \( 1 \leq i \leq n_1 \) and \( z_{n1+i} = y_i \) for \( 1 \leq i \leq n_2 \). Let \( \hat{f}(Z_i) = (N-1)^{-1} \sum_{j \neq i} K_{\gamma,z_i,z_j} \) be the leave-one-out estimator of \( f(Z_i) \). Then \( \int \hat{f}(x)^2 dx - 2 \int \hat{f}(x)f(x) \) can be consistently estimated by the following cross-validation function:
\[
CV(h, \lambda) = \frac{1}{N^2} \sum_{i=1}^N \sum_{j=1}^N \hat{K}_{\gamma,z_i,z_j} - \frac{2}{N(N-1)} \sum_{i=1}^N \sum_{j \neq i} K_{\gamma,z_i,z_j},
\]
where \( K_{\gamma,z_i,z_j} = W_{h,z_i,z_j} L_{\lambda,z_i,z_j}, L_{\lambda,z_i,z_j} \) and \( W_{h,z_i,z_j} \) are defined in (2.2) and (2.3), respectively, and where \( \hat{K}_{\gamma,z_i,z_j} = \hat{W}_{h,i,j} \hat{L}_{\lambda,i,j} \), \( \hat{W}_{h,i,j} = \int W_{h,z_i,z_j} W_{h,z_j,z} \, dz \) and \( \hat{L}_{\lambda,i,j} = \sum_{z \in \mathbb{R}^2} L_{\lambda,z_i,z} L_{\lambda,z_j,z_j} \). It can be shown that \( \hat{W}_{h,z_i,z_j} = \prod_{s=1}^q h_s^{-1} w((X_{is} - X_{js})/h_s) \), where \( \tilde{w}(v) = \int w(u)w(v-u) \, du \) is the two-fold convolution kernel derived from \( w(\cdot) \), which is also a standard second order kernel function. For example, if \( w(v) = e^{-v^2/2}/\sqrt{2\pi} \), i.e., a standard normal kernel, then \( \tilde{w}(v) = e^{-v^2/4}/\sqrt{4\pi} \), a normal kernel with mean zero and variance two, which follows from the fact that two independent \( N(0,1) \) random variables sum to a \( N(0,2) \) random variable. Therefore, we select the smoothing parameters by minimizing the CV function defined in (2.11).

Letting \( \hat{(h_1, \ldots, h_q)} \) and \( \hat{(\lambda_1, \ldots, \lambda_r)} \) denote the cross-validated values of \( (h_1, \ldots, h_q) \) and \( (\lambda_1, \ldots, \lambda_r) \), Li & Racine (2003) and Ouyang, Li & Racine (2006) have proved the following results, we summarize them in a condition below for ease of reference.

\[
(C4) \frac{\hat{h}_s}{h_s^0} - 1 \to 0 \text{ in probability, and } \frac{\hat{\lambda}_s}{\lambda_s^0} - 1 \to 0 \text{ in probability, where } h_s^0 = a_s^0 n^{-\zeta_s}, \text{ and } \lambda_s^0 = b_s^0 n^{-2\zeta} \text{ for some } \zeta > 0, \text{ where } a_s^0 \text{ and } b_s^0 \text{ are some finite constants.}
\]

\( h_s^0 \) and \( \lambda_s^0 \) in (C4) above are the non-stochastic optimal smoothing parameters that minimize the integrated mean squared difference \( \int E[(\hat{f}(z) - f(z))^2] \, dz \). To establish (C4), Li & Racine (2003)

\[\text{Li & Racine (2003) only consider the case for which } h_1 = \cdots = h_q = h \text{ and } \lambda_1 = \cdots = \lambda_r = \lambda. \text{ It is straightforward to generalize the result of Li & Racine (2003) to the vector } h \text{ and } \lambda \text{ case, and the result should be modified as given here.}\]
and Ouyang et al. (2006) assumed that (i) $f(x)$ is four times differentiable with respect to $x^c$; (ii) there exists $x^d, z^d \in S$, such that $f(x^c, x^d) \neq f(x^c, z^d)$ for all $x^c$ in a subset of the support of $X^c$ with positive measure; (iii) $Pr(X^d_s = x^d_s)$ is not constant for all $x^d_s \in \{0, 1, \ldots, c_s - 1\}$. For details on the regularity conditions that ensure (C4) holds, see Li & Racine (2003) and Ouyang et al. (2006).

When $f(x)$ and $g(x)$ have unbounded support and non-vanishing second derivative functions with respect to $x^c_s$ for all $s = 1, \ldots, q$, Li & Racine (2003) and Ouyang et al. (2006) show that
\[ \hat{\zeta} = 1/(4+q), \]

i.e., $\hat{h}_s = O_p(n^{-1/(4+q)})$ and $\hat{\lambda} = O_p(n^{-2/(4+q)})$. When the support of $x^c$ is bounded, the kernel estimator may suffer from the boundary bias problem. However, even in this case, the cross-validated smoothing parameters $\hat{h}_s$ still converge to the optimal smoothing parameter values in the sense that $\hat{h}_s/h^0_s \to 1$ in probability. In fact, Stone (1984) showed that as long as the marginal density function for $X_s$ is bounded (for all $s = 1, \ldots, q$), then $\hat{h}_s/h^0_s \to 1$ almost surely. However, the rate at which $\hat{h}_s$ (or $h^0_s$) converges to zero may be different when the support of $x^c$ is bounded. For example, for the case where $q = 1$ and where $x^c$ is uniformly distributed, Ouyang et al. (2006) (their Lemma 3.1) showed that $\zeta = 1/2$ so that $\hat{h} = O_p(n^{-1/2})$. The reason why $\hat{h} \to 0$ (or $\hat{h}_s/h^0_s \to 1$) in probability even when $x^c$ has bounded support (i.e., uniformly distributed) is as follows. Consider the case where the support of $X$ is $[0, 1]$. Then, at the boundary regions $x \in [0, h] \cup [1-h, 1]$, $\hat{f}(x) - f(x)$ usually does not converge to zero in probability due to the boundary bias problem. In this case, only when $h \to 0$ will the integrated (mean) squared error converge to zero since the boundary regions shrink to zero length as $h \to 0$. Since $\hat{h}_s$ asymptotically minimizes the integrated squared error, we know that $\hat{h}_s$ must converge to zero whether or not $X$ has bounded support.

Let $\hat{T}_n (\hat{I}_n)$ denote the test statistic $T_n (I_n)$ but with $(h, \lambda)$ being replaced by $(\hat{h}, \hat{\lambda})$, the cross-validated smoothing parameters. The next theorem shows that the test statistic $\hat{T}_n$ has the same asymptotic distribution as $T_n$.

**Theorem 2.2.** Assuming that conditions (C1) through (C3) hold, then under $H_0$ we have

\[ \hat{T}_n = (n_1n_2\hat{h}_1 \ldots \hat{h}_q)^{1/2}\hat{I}_n/\hat{\sigma}_n \to N(0, 1) \text{ in distribution}, \]

where $\hat{\sigma}_n$ is defined the same way as in $\sigma_n$ but with $(h, \lambda)$ replaced by $(\hat{h}, \hat{\lambda})$. 


The proof of Theorem 2.2 is given in the appendix.

2.2. Comparison with Non-Smoothing Tests. In this section we discuss the local power property of our $\hat{T}_n$ test and compare it with some non-smoothing tests. For ease of exposition we only consider the case where $x$ (y) is a continuous variable of dimension $q$.\(^3\) One class of non-smoothing tests involves fixing the value of $h_s$ in a smoothing test, say letting $h_s = 1$ for all $s = 1, \ldots, q$; see Anderson, Hall & Titterington (1994), Fan (1998), Fan & Li (2000) and the references therein. Another class of non-smoothing tests involves testing for the equality of two CDFs. There is a rich literature on testing the equality of two CDFs, i.e., where one tests the null hypothesis that $F(x) = G(x)$ for all $x$ where $F(x)$ and $G(x)$ are two unknown CDFs.

Anderson et al. (1994) show that one can set $h_1 = \ldots = h_q = 1$ in the $p_{n,1} I_{n}^{a}$ test to obtain a non-smoothing test of the form

$$(2.12)$$

$$I_{n, h=1} = \sqrt{n_{1} n_{2}} \left[ \sum_{i=1}^{n_{1}} \sum_{j=1}^{n_{2}} \frac{W(X_i - X_j)}{n_{1}^2} + \sum_{i=1}^{n_{1}} \sum_{j=1}^{n_{2}} \frac{W(Y_i - Y_j)}{n_{2}^2} - \sum_{i=1}^{n_{1}} \sum_{j=1}^{n_{2}} \frac{W(X_i - Y_j)}{n_{1} n_{2}} - \sum_{i=1}^{n_{2}} \sum_{j=1}^{n_{1}} \frac{W(X_j - Y_i)}{n_{1} n_{2}} \right],$$

where $W(X_i - X_j) = \prod_{s=1}^{q} w(X_{is} - X_{js})$ which is obtained from $W_{h, x_i, x_j}$ by setting $h_1 = \ldots = h_q = 1$. It can be shown that, for a wide class of kernel functions $w(\cdot)$, $I_{n, h=1}$ leads to a consistent test for the null hypothesis of equality of $f(x)$ and $g(x)$ for almost all $x \in \mathbb{R}^q$. It is well known that a non-smoothing test such as $I_{n, h=1}$ does not have an asymptotic normal distribution. It can be shown that $n I_{n, h=1}$ has an asymptotic weighted $\chi^2$ distribution of the form $\sum_{l=1}^{\infty} c_l \chi_l(1)$, where the $c_l$’s are some constants, and the $\chi_l(1)$’s are independent $\chi$ random variable with one degree of freedom. The weight $c_l$ depends on the unknown density functions $f(x)$ and $g(x)$. Therefore, it is impossible to tabulate this asymptotic distribution. However, bootstrap methods may be used to approximate the null distribution of $I_{n, h=1}$.

One can also test the null hypothesis of equality of two distributions based upon estimation of the unknown CDFs. For example, a Kolmogorov-Smirnov type test can constructed based on $\sup_{x \in \mathbb{R}^q} |F(x) - G(x)|$. Let $F_{n_1}(\cdot)$ and $G_{n_2}(\cdot)$ be the empirical CDFs of $\{X_i\}_{i=1}^{n_1}$ and $\{Y_i\}_{i=1}^{n_2}$.

\(^3\)Adding discrete components to $x$ (y) will require more complex notation, but will not affect the result of the local power analysis.
respectively. Formally, we have

\[
KS_n = \sup_{x \in \mathbb{R}} \left| \sqrt{\frac{2n_1n_2}{n_1 + n_2}} \left[ \frac{1}{n_1} \sum_{i=1}^{n_1} I(X_i \leq x) - \frac{1}{n_2} \sum_{i=1}^{n_2} I(Y_i \leq x) \right] \right|.
\]

It is easy to check that \(\sqrt{2n_1n_2/(n_1 + n_2)}[F_{n_1}(\cdot) - G_{n_2}(\cdot)]\) converges to a zero mean Gaussian process (under \(H_0\), say \(GP(\cdot)\). Then it follows from the continuous mapping theorem that \(KS_n \to \sup_{x \in \mathbb{R}} |GP(x)|\) in distribution under \(H_0\).

A Cramér-von Mises (CM) type test statistic (based on \(\int [F(x) - G(x)]^2 dx\)) can be constructed by

\[
CM_n = \frac{2n_1n_2}{n_1 + n_2} \int \left[ \frac{1}{n_1} \sum_{i=1}^{n_1} I(X_i \leq x) - \frac{1}{n_2} \sum_{i=1}^{n_2} I(Y_i \leq x) \right] \left[ \frac{1}{n_1} \sum_{j=1}^{n_1} I(X_j \leq x) - \frac{1}{n_2} \sum_{j=1}^{n_2} I(Y_j \leq x) \right] dx.
\]

It can be shown that \(CM_n \to \int GP(x)^2 dx\) in distribution under \(H_0\).

Under the null hypothesis, all the non-smoothing tests, \(I_{n,h=1}\), \(KS_n\) and \(CM_n\) have non-standard asymptotic distributions. Therefore, in practice, some bootstrap methods are often used to approximate the null distributions of \(I_{n,h=1}\), \(KS_n\) and \(CM_n\). A simple bootstrap method involves resampling \(\{Z_i\}_{i=1}^{n_1+n_2}\) with replacement, where \(Z_i = X_i\) for \(i = 1, \ldots, n_1\), and \(Z_{i+n_1} = Y_i\) for \(i = 1, \ldots, n_2\). One then uses the bootstrap sample \(\{X_i^*\} = \{Z_i^*\}_{i=1}^{n_1}\) and \(\{Y_i^*\}_{i=1}^{n_2} = \{Z_{i+n_1}^*\}_{i=1}^{n_2}\) to compute \(I_{n,h=1}^*, KS_n^*\) and \(CM_n^*\), respectively.

Note that since both the \(KS_n\) and the \(CM_n\) tests involve indicator functions, therefore, the sup operator in \(KS_n\) can be replaced by maximization over the \(n_1 + n_2\) sample realizations as follows:

\[
(2.13) \quad KS_n = \max_{1 \leq j \leq n_1+n_2} \left| \sqrt{\frac{2n_1n_2}{n_1 + n_2}} \left[ \frac{1}{n_1} \sum_{i=1}^{n_1} I(X_i \leq Z_j) - \frac{1}{n_2} \sum_{i=1}^{n_2} I(Y_i \leq Z_j) \right] \right|.
\]

Similarly, all integration required for the computation of \(CM_n\) can be computed easily leading to the following result:

\[
(2.14) \quad CM_n = \frac{2n_1n_2}{n_1 + n_2} \left\{ \frac{1}{n_1n_2} \left[ \sum_{i=1}^{n_1} \sum_{j=1}^{n_2} \max\{X_i, Y_j\} + \sum_{i=1}^{n_1} \sum_{j=1}^{n_2} \max\{Y_i, X_j\} \right] \right\}
\]

\[
- \frac{1}{n_1} \sum_{i=1}^{n_1} \sum_{j=1}^{n_1} \max\{X_i, X_j\} - \frac{1}{n_2} \sum_{i=1}^{n_2} \sum_{j=1}^{n_2} \max\{Y_i, Y_j\} \right\}.
\]
We now compare the local power properties of the smoothing test and the non-smoothing tests. We consider two types of local alternatives. One is a sequence of ‘regular’ or ‘Pitman’ alternatives given by

\[ LH_r : f(x) = g(x) + \alpha_n \Delta(x), \]

where \( \int \Delta(x) dx = 0 \) and \( \alpha_n \to 0 \) as \( n \to \infty \). The second is a sequence of so-called ‘singular’ local alternatives which was first introduced by Rosenblatt (1975) and is given by

\[ LH_s : f(x) = g(x) + \alpha_n \Delta_n(x), \]

where \( \int \Delta_n(x) dx = 0 \), \( \int \Delta_n^2(x) dx \to 0 \), and \( \alpha_n \to 0 \) as \( n \to \infty \). For example, one can have \( \Delta_n(x) = \sum_{j=1}^p d_j((x - l_j)/\beta_n) \), where \( p \) is a positive integer, \( l_1, \ldots, l_p \) are constant vectors in \( \mathbb{R}^q \), \( d_1(\cdot), \ldots, d_p(\cdot) \) are bounded smooth functions satisfying \( \int d_j(x) dx = 0 \) for \( j = 1, \ldots, p \), and \( \beta_n \to 0 \) as \( n \to \infty \). Then it is easy to see that \( \int \Delta_n(x)^2 dx = O(\beta_n) = o(1) \).

In finite-sample applications, the ‘singular’ alternative corresponds to a ‘rapidly changing’ or a ‘high frequency’ density function. In the simulations reported in Section 4, we use some mixture normal distributions (density functions having multiple peaks) to represent ‘high frequency’ density functions.

It is well established that non-smoothing tests can detect both the Pitman and the Rosenblatt local alternatives that approach the null at the rate of \( O(n^{-1/2}) \). In contrast, smoothing tests can detect Pitman local alternatives converging to the null at rate \( O((n^2 h_1 \ldots h_q)^{-1/4}) \), which is slower than \( n^{-1/2} \) because \( h_s \to 0 \) for all \( s = 1, \ldots, q \). Therefore, for Pitman local alternatives, a non-smoothing test is asymptotically more powerful than a smoothing test. However, it is also known that, for the class of ‘singular’ local alternatives, a smoothing test can detect local alternatives that approach the null at an rate of \( O(n^{-1/2}) \), see Ghosh & Huang (1991), Fan (1998), and Fan & Li (2000). Hence, a smoothing test is more powerful than a non-smoothing test for ‘singular’ local alternatives. Indeed the simulation evidence reported in Fan & Li (2000) reveals strong support for the above theoretical local power analysis.

The existing simulation comparisons between \( I_{n,h=1} \) and \( I_n \) usually use some ad-hoc selection of \( h \) such as \( h_{s,ad-hoc} = z_{s,ad}(n_1 + n_2)^{-1/(4+q)} \) when computing \( I_n \), where \( z_{s,ad} \) is the sample standard deviation of \( \{Z_is\}_{i=1}^{n_1+n_2} \) \( s = 1, \ldots, q \). In Section 4 we show that by cross-validated (CV) selection
of $h_s$, the resulting test is often more powerful (in finite-sample applications) than either using $h_{s,ad-hoc}$ or using $h = 1$. The superior performance of the CV-based test arises because the CV method can automatically adapt to the smoothness of the underlying density functions. When $f(x)$ ($g(x)$) is a relatively smooth (i.e., unimodal and slowly changing) function of $x_s$, the CV method will select a relatively large value for $h_s$; when $f(x)$ ($g(x)$) is a relatively high frequency function of $x_s$ (i.e., multimodal and peaked), the CV method will select a small value for $h_s$ resulting in a test having high power against either low or high frequency alternatives. The simple ad-hoc rule of selecting $h_{s,ad-hoc}$ or even fixing $h = 1$ cannot possess such flexibility which can harm their power as will be seen.

The CDF-based tests have local power properties similar to non-smoothing tests $I_{n,h=1}$. Therefore, they are asymptotically more powerful than a smoothing test against Pitman local alternatives, and they may be less powerful against ‘singular’ local alternatives. In Section 4 we report simulation results that examine the finite-sample performance of our smoothing test versus some non-smoothing tests including the Kolmogorov-Smirnov test and the Cramer-von Mises test discussed above.

The finite-sample implication of the above local power analysis is that, for density functions $f(x)$ and $g(x)$ that change slowly as $x$ changes (the so-called ‘low frequency’ density), a non-smoothing test may be more powerful than a smoothing test. While for ‘high frequency’ density functions, i.e., $f(x)$ and $g(x)$ changes rapidly as $x$ changes, a smoothing test is expected to be more powerful. We will examine the finite-sample performances of non-smoothing test $T_n$ and some non-smoothing tests in Section 4.

2.3. A Bootstrap Procedure. Theorems 2.1 and 2.2 show that $T_n$ and $\hat{T}_n$ have asymptotic standard normal null distributions. However, existing simulation results suggest that this limiting normal distribution is in fact a poor approximation to the finite-sample distribution of $T_n$. Our experience also shows that the same holds true for the $\hat{T}_n$ statistic. Therefore, in order to better approximate the null distribution of $\hat{T}_n$, in applied settings we advocate the use of the following bootstrap procedure.

Randomly draw $n_1$ observations from the pooled sample $\{Z_j\}_{j=1}^{n_1+n_2}$ with replacement, and call the resulting sample $\{X_i^*\}_{i=1}^{n_1}$; then randomly draw another $n_2$ observations from $\{Z_j\}_{j=1}^{n_1+n_2}$ with
replacement, and call the resulting sample \( \{Y_i^*\}_{i=1}^{n_2} \). Compute the bootstrap test statistic given by 
\[ \hat{T}_n^* = (n_1 n_2 \hat{h}_1 \ldots \hat{h}_q)^{1/2} \hat{I}_n^* / \hat{\sigma}_n^* , \]
where \( \hat{I}_n^* \) and \( \hat{\sigma}_n^* \) are defined the same way as \( \hat{I}_n \) and \( \hat{\sigma}_n \) except that 
\( X_i \) and \( Y_i \) are replaced by \( X_i^* \) and \( Y_i^* \), respectively. We repeat this procedure a large number of 
times, say \( B = 399 \) times (Davidson & MacKinnon (2000)), and we use the empirical distribution of 
the \( B \) bootstrap statistics \( \{\hat{T}_{n,l}^*\}_{l=1}^{B} \) to approximate the null distribution of \( \hat{T}_n \). Empirical P-values 
can be computed via 
\[ \hat{P} = B^{-1} \sum_{l=1}^{B} I(\hat{T}_{n,l}^* > \hat{T}_n), \]
where \( I(\cdot) \) is an indicator function, which is 
simply the proportion of resampled test statistics under the null that are more extreme than the 
statistic itself.

Note that we use the same smoothing parameters \( (\hat{h}, \hat{\lambda}) \) when computing \( \hat{T}_n^* \), i.e., we do not re 
cross-validate for each bootstrap replication. Therefore, this bootstrap procedure is computationally less costly than the computation of \( \hat{T}_n \), which involves a cross-validation procedure. The next 
theorem proves the validity of the proposed bootstrap method.

**Theorem 2.3.** Define 
\[ \hat{T}_n^* = (n_1 n_2 \hat{h}_1 \ldots \hat{h}_q)^{1/2} \hat{I}_n^* / \hat{\sigma}_n^* , \]
Assuming that the same conditions given in 
**Theorem 2.2** hold, but under the null hypothesis, then we have 
\[ \sup_{z \in \mathbb{R}} \left| P \left( \hat{T}_n^* \leq z \left| \{X_i, Y_i\}_{i=1}^{n} \right. \right) - \Phi(z) \right| = o_p(1) , \]
where \( \Phi(\cdot) \) is the cumulative distribution function of a standard normal random variable.

The proof of Theorem 2.3 is given in the appendix.

In words, Theorem 2.3 states that \( \hat{T}_n^* \) converges in distribution to a \( N(0, 1) \) in distribution in 
probability. In the bootstrap hypothesis testing literature, the notion of ‘convergence in distribution 
with probability one’ is often used to describe the asymptotic behavior of bootstrap tests, where 
one states that the left-hand-side of (2.16) is of small order \( o(1) \) with probability one. ‘Convergence 
in distribution in probability’ is much easier to establish than ‘convergence in distribution with 
probability one’, and runs parallel to that of ‘convergence in probability’ and ‘convergence with 
probability one’.

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3. A Nonparametric Test for the Equality of Conditional Density Functions with Mixed Categorical and Continuous Data

In this section we consider the problem of testing the equality of two conditional density functions. We will only consider the case for which the conditioning variable is categorical by nature. There are two reasons for this. First, technically it is difficult to handle the continuous conditioning variable case when the density function is not bounded below by a positive constant. The second consideration is a practical one. In empirical applications it is likely that one is interested in knowing the distribution of a continuous variable, say the distribution of income conditional on a discrete variable such as a person’s sex, or perhaps their level of education.

Given that we only consider a discrete conditioning variable in this section, we shall employ slightly different notation for what follows. We shall continue to use \( x = (x^c, x^d) \in \mathbb{R}^q \times \mathbb{S}^r \) to denote a mixture of continuous and discrete variables, and we use \( w \) to denote the conditioning discrete variable. \( w \) can be a multivariate discrete variable. We use \( S_w \) to denote the support of \( W \), and we assume that \( P(w) = Pr(W = w) \) is bounded below by a positive constant for all \( w \in S_w \). Suppose we have i.i.d data, \( \{X_i, U_i\}_{i=1}^{n_1} \), which are random draws from the joint density function \( f(x, w) \) along with i.i.d. draws of \( \{Y_i, V_i\}_{i=1}^{n_2} \) from the joint density function \( g(x, w) \). We use \( f(x|w) \) \( (g(x|w)) \) to denote the conditional density function of \( X \) \( (Y) \) conditional on \( U = w \) \( (V = w) \). We use \( S_w \) to denote a subset of the support of \( w \) such that one is interested testing for

\[
H_c^0 : f(x|w) = g(x|w) \quad \text{for all} \quad w \in S_w, x^d \in \mathbb{S}^r \quad \text{and for almost all} \quad x^c \in \mathbb{R}^q.
\]

against the alternative hypothesis \( H_c^1 \) that \( f(x|w) \neq g(x|w) \) on a set with positive measure.

Define \( p_f(w) = Pr(U = w) \) and \( p_g(w) = Pr(V = w) \). Note that \( p_f(w) \) can differ from \( p_g(w) \). For example, consider the case where \( x \) is income and \( w \) is a dummy variable equal to one for males, zero otherwise. Clearly the percentage of males in two populations can differ, i.e., \( p_f(w) \) may not equal \( p_g(w) \).
Using $f(x|w) = f(x, w)/p_f(w)$ and $g(x|w) = g(x, w)/p_g(w)$, we will construct a test statistic based on

$$J = \sum_{w \in S_w} \int [f(x|w) - g(x|w)]^2 dx$$

$$= \sum_{w \in S_w} \int \left[ \frac{f(x, w)^2}{p_f(w)^2} + \frac{g(x, w)^2}{p_g(w)^2} - \frac{2f(x, w)g(x, w)}{p_f(w)p_g(w)} \right] dx,$$

where $\int dx = \sum_{x \in \mathbb{R}^r} \int dx^c$.

Let $I_{u_i,w} = I(U_i = w)$ denote an indicator function which equals one if $U_i = w$ and zero otherwise. $I_{v_i,w}$ is similarly defined. We estimate the joint density of $f(x, w)$ and $g(x, w)$ by

$$\hat{f}(x, w) = \frac{1}{n_1} \sum_{i=1}^{n_1} K_{\gamma_i, x} I_{u_i,w},$$

$$\hat{g}(x, w) = \frac{1}{n_2} \sum_{i=1}^{n_2} K_{\gamma_i, x} I_{v_i,w},$$

where $K_{\gamma_i, x} = W_{h,x_i,x}L_{\lambda_i, x, i}$ with $W_{h,x_i,x}$ and $L_{\lambda_i, x, i}$ as defined in Section 2. Also, we estimate $p_f(w)$ and $p_g(w)$ by

$$\hat{p}_f(w) = \frac{1}{n_1} \sum_{i=1}^{n_1} I(U_i = w)$$

$$\hat{p}_g(w) = \frac{1}{n_2} \sum_{i=1}^{n_2} I(V_i = w).$$

Define the leave-one-out empirical functions by $F_{n,-i}(x) = (n_1 - 1)^{-1} \sum_{j \neq i}^{n_1} I(X_j = x)$ and $G_{n,-i}(x) = (n_2 - 1)^{-1} \sum_{j \neq i}^{n_2} I(Y_j = x)$. Replacing $f$, $g$, $p_f$ and $p_g$ by their estimators in (3.18), and using the short-hand notation $\hat{p}_f = \hat{p}_f(w)$ and $\hat{p}_g = \hat{p}_g(w)$, we obtain a feasible test statistic given
by

\[ J_n = \sum_{w \in S_w} \int \left[ \frac{\hat{f}(x, w)}{\hat{p}_f^2} dF_{n,-i}(x) + \frac{\hat{g}(x, w)}{\hat{p}_g^2} dG_{n,-i}(x) - \frac{\hat{f}(x, w)}{\hat{p}_f \hat{p}_g} dF_n(x) - \frac{\hat{g}(x, w)}{\hat{p}_f \hat{p}_g} dG_n(x) \right] \]

\[ = \sum_{w \in S_w} \left\{ \sum_{i} \sum_{j \neq i} \left[ \frac{\hat{K}_{\gamma,i,y} I_{u_i,w} I_{u_j,z}}{n_1(n_1 - 1)\hat{p}_f^2} + \frac{\hat{K}_{\gamma,i,y} I_{v_i,w} I_{v_j,u}}{n_2(n_2 - 1)\hat{p}_g^2} \right] \right\} - \frac{1}{n_1 n_2 \hat{p}_f \hat{p}_g} \sum_{i} \sum_{j} \left[ \hat{K}_{\gamma,i,y} I_{u_i,w} I_{v_j,u} + \hat{K}_{\gamma,i,y} I_{u_i,w} I_{v_j,u} \right] \]

(3.21)

where \( \hat{K}_{h,x_i,y_j} = \hat{W}_{h,x_i,x} L_{h,x_i,x} \), and \( \hat{W}_{h,x_i,x} \) and \( \hat{L}_{h,x_i,x} \) are defined in Section 2. Also, as in Section 2, \( j_i \) is \( \sum_{i=1}^{n_1} \) if the summand has \( (X_i, U_i) \) as its argument, and \( j_i \) is \( \sum_{i=1}^{n_2} \) if the summand has \( (Y_i, V_i) \) as its argument. We choose \((h_1, \ldots, h_q, \lambda_1, \ldots, \lambda_q)\) by the cross-validation method discussed in Section 2, and we use \((\hat{h}_1, \ldots, \hat{h}_q, \hat{\lambda}_1, \ldots, \hat{\lambda}_q)\) to denote the cross-validated smoothing parameters. We will use \( \hat{J}_n \) to denote our test statistic as defined in (3.21) but with \((h_1, \ldots, h_q, \lambda_1, \ldots, \lambda_q)\) replaced by \((\hat{h}_1, \ldots, \hat{h}_q, \hat{\lambda}_1, \ldots, \hat{\lambda}_q)\).

We make the following additional assumption.

(C5) For all \((x^d, w) (x^d \in S^r)\), both \( f(\cdot, x^d, w) \) and \( g(\cdot, x^d, w) \) are bounded (from above by some positive constants) and continuous functions (continuous with respect to \( x^c \)).

The asymptotic null distribution of our test statistic is given in the next theorem.

**Theorem 3.1.** Assuming that conditions (C1) to (C5) hold, then under \( H_0^* \), we have

\[ \hat{T}_{n,c} \overset{def}{=} (n_1 n_2 \hat{h}_1 \ldots \hat{h}_q)^{1/2} \hat{J}_n / \hat{\sigma}_{n,c} \rightarrow N(0, 1) \] in distribution,

where

\[ \hat{\sigma}_{n,c}^2 = 2(n_1 n_2 \hat{h}_1 \ldots \hat{h}_q) \sum_{w \in S_w} \left[ \sum_{i=1}^{n_1} \sum_{j \neq i} \left( \frac{\hat{K}_{\gamma,i,y} I_{u_i,w} I_{u_j,w}}{n_1^2 \hat{p}_f(w)^4} \right)^2 + \sum_{i=1}^{n_2} \sum_{j \neq i} \left( \frac{\hat{K}_{\gamma,i,y} I_{v_i,w} I_{v_j,w}}{n_2^2 \hat{p}_g(w)^4} \right)^2 \right] + \sum_{i=1}^{n_1} \sum_{j=1}^{n_2} \left( \frac{\hat{K}_{\gamma,i,y} I_{u_i,w} I_{v_j,u}}{n_1^2 n_2^2 \hat{p}_f(w)^2 \hat{p}_g(w)^2} \right)^2 + \sum_{i=1}^{n_2} \sum_{j=1}^{n_1} \left( \frac{\hat{K}_{\gamma,i,y} I_{v_i,w} I_{u_j,u}}{n_1 n_2^2 \hat{p}_f(w)^2 \hat{p}_g(w)^2} \right)^2 \]

The proof of Theorem 3.1 is given in the appendix.
In practice we recommend the use of the following bootstrap procedure to approximate the null distribution of $\hat{T}_{n;c}$.

Let $Z_i = \{X_i, U_i\}$ for $i = 1, \ldots, n_1$, and $Z_{i+n_1} = \{Y_i, V_i\}$ for $i = 1, \ldots, n_2$. Then randomly draw $n_1$ observations from the pooled sample $\{Z_j\}_{j=1}^{n_1+n_2}$ with replacement, and call the resulting sample $\{X^*_i, U^*_i\}_{i=1}^{n_1}$, and then randomly draw another $n_2$ observations from $\{Z_j\}_{j=1}^{n_1+n_2}$ with replacement, and call the resulting sample $\{Y^*_i, V^*_i\}_{i=1}^{n_2}$. Compute a test statistic $\hat{T}_{n;c}^* = \left(\frac{n_1n_2}{n_1+n_2}\right)^{1/2} \hat{J}_n^*/\hat{\sigma}_{n,c}^*$, where $\hat{J}_n$ and $\hat{\sigma}_{n,c}$ are defined the same way as $\hat{J}_n$ and $\hat{\sigma}_{n,c}$ except that $(X_i, U_i)$ and $(Y_i, V_i)$ are replaced by $(X^*_i, U^*_i)$ and $(Y^*_i, V^*_i)$, respectively. We repeat this procedure a large number of times (say $B = 399$), and we use the empirical distribution of the $B$ bootstrap statistics $\{\hat{T}_{n,c,i}^*\}_{i=1}^B$ to approximate the null distribution of $\hat{T}_{n,c}$.

Note that we use the same $(\hat{h}, \lambda)$, and the same $\hat{p}_f(w)$ and $\hat{p}_g(z)$ when computing $\hat{T}_{n,c}^*$. The next theorem states that the above bootstrap method can be used to approximate the null distribution of $\hat{T}_{n,c}$.

**Theorem 3.2.** Define $\hat{T}_{n,c}^* = \left(\frac{n_1n_2}{n_1+n_2}\right)^{1/2} \hat{J}_n^*/\hat{\sigma}_{n,c}^*$. Assume the same conditions as in Theorem 3.1 except that we do not impose the null hypothesis $H^c_0$. Then we have

$$\sup_{z \in \mathbb{R}} |P(\hat{T}_{n,c}^* \leq z|\{X_i, U_i, Y_i, V_i\}_{i=1}^n) - \Phi(z)| = o_p(1),$$

where $\Phi(\cdot)$ is the cumulative distribution function of a standard normal random variable.

The proof of Theorem 3.2 is similar to the proof of Theorem 3.1 and is thus omitted.

Note the in constructing our conditional density test statistic $T_{n,c}$, we smooth both $x^c$ and $x^d$, we do not smooth over the conditional discrete covariate $w$. In practice one can also smooth the conditional discrete variable $w$ when testing $H^c_0$. For expositional simplicity, we discuss the case where $w$ is a scalar below. In this case, one replaces the indicator function, say, $I_{u_i, w} = I(U_i = w)$ by $I_{\lambda_0, u_i, w} = I(\lambda_0, U_i, w)$ which is defined in (1), and $\lambda_0$ is the smoothing parameter associated with
w. The modified test statistic becomes

\[ J_{n,\lambda_0} = \sum_{w \in \mathcal{S}_w} \left\{ \sum_i \sum_{j \neq i} \left[ \frac{\tilde{K}_{\gamma, x, j} l_{\lambda_0, u, w} l_{\lambda_0, u, z}}{n_1(n_1 - 1) \tilde{p}_{j}} + \frac{\tilde{K}_{\gamma, y, j} l_{\lambda_0, v, w} l_{\lambda_0, v, w}}{n_2(n_2 - 1) \tilde{p}_{g}} \right] \right\} - \frac{1}{n_1 n_2 \tilde{p}_{f} \tilde{p}_{g}} \sum_i \sum_j [\tilde{K}_{\gamma, x, j} l_{\lambda_0, u, w} l_{\lambda_0, v, w} + \tilde{K}_{\gamma, y, j} l_{\lambda_0, u, w} l_{\lambda_0, v, w}] \]

(3.23)

where \( \tilde{p}_{f} = n_1^{-1} \sum_{i=1}^{n} l_{\lambda_0, u, w} \) and \( \tilde{p}_{g} = n_2^{-1} \sum_{i=1}^{n} l_{\lambda_0, v, w} \). For the test statistic \( J_{n,\lambda_0} \) we also need a different method for selecting the smoothing parameters. In the framework of estimating a conditional density function, Hall, Racine and Li (2004) propose to select the smoothing parameters by minimizing a sample analogue of \( \int [f(y|x) - f(y|x)]^2 \mu(x) dx \), where \( \mu(x) \) is a weight function and \( \int dx \) should be interpreted as \( \sum_{x \in \mathcal{X}} \sum_{y \in \mathcal{Y}} \int dx \mu(x) \). We suggest to use the least squares cross-validation method proposed by Hall, Racine and Li (2004) in selecting the smoothing parameters \( (\lambda_0, \lambda_1, \ldots, \lambda_r) \) and \( (h_1, \ldots, h_q) \). According to Hall, Racine and Li (2004), the resulting smoothing parameters have the order of \( \lambda_s \sim n^{-2/(4+q)} \) for \( s = 0, 1, \ldots, r \); and \( h_s \sim n^{-1/(4+q)} \), see Hall et al (2004) for details. In Section 4 we also compute the test statistic \( J_{n,\lambda_0} \) and compare it with the test statistic \( J_n \). The bootstrap procedure for obtaining critical values for the \( J_{n,\lambda_0} \) test is similar to that for \( J_n \), except that one replaces the indicator functions by the (discrete variable) corresponding kernel functions. The simulations reported there show that \( J_{n,\lambda_0} \) test has better power performance than that of \( J_n \).

The asymptotic analysis of \( J_{n,\lambda_0} \) is much more involved than that of \( J_n \). This is because that the \( J_n \) statistic uses \( \tilde{p}_{f}(w) \) and \( \tilde{p}_{g}(w) \) to estimate \( p_{f}(w) \) and \( p_{g}(w) \); and that \( \tilde{p}_{f}(w) - p_{f}(w) = O_p(n_1^{-1/2}) \) and \( \tilde{p}_{g}(w) - p_{g}(w) = O_p(n_2^{-1/2}) \); they both have the parametric root-n convergence rate. In the Appendix A we show that the asymptotic distribution of \( T_{n,c} \) is unaffected if one replaces \( \tilde{p}_{f}(w) \) and \( \tilde{p}_{g}(w) \) by \( p_{f}(w) \) and \( p_{g}(w) \). In contrast, the \( J_{n,\lambda_0} \) statistic uses kernel-smoothed probability estimators \( \tilde{p}_{f}(w) \) and \( \tilde{p}_{g}(w) \) to estimate \( p_{f}(w) \) and \( p_{g}(w) \); and that \( \tilde{p}_{f}(w) - p_{f}(w) = O_p(\sum_{s=1}^{r} h_s^2 + \sum_{s=0}^{r} \lambda_s + (nh_1 \ldots h_q)^{-1/2}) = O_p(n_1^{-2/(4+q)}) \), and \( \tilde{p}_{g}(w) = O_p(n_2^{-1/(4+q)}) \). They both have (slow) nonparametric rate of convergence. Estimating of \( p_{f}(w) \) and \( p_{g}(w) \) by nonparametric estimators \( \tilde{p}_{f}(w) \) and \( \tilde{p}_{g}(w) \) may affect the asymptotic distribution of \( J_{n,\lambda_0} \), rendering the asymptotic analysis to be much more complex than that of \( J_n \). We leave the asymptotic analysis of \( J_{n,\lambda_0} \) as a future
research topic. We use simulations to examine the finite-sample performance of $J_{n,\lambda_0}$ based on bootstrap critical values.

4. Monte Carlo Simulations

In this section we consider the finite-sample performance of the proposed tests in a variety of settings. We begin by comparing the performance of the proposed unconditional density test ($T_n$) with its frequency-based and ad-hoc smoothing parameter selection counterparts when there exist both continuous and discrete variables. Next, we compare the proposed unconditional density test with non-smooth tests under both ‘high frequency’ and ‘low frequency’ alternatives. The local power analysis of Section 2.2 suggests that a non-smoothing test is likely to be more powerful against low frequency (i.e., slowly changing) density functions, while a smoothing test is expected to be more powerful for high frequency (i.e., rapidly changing) density functions. Finally, we examine the finite-sample performance of the proposed conditional density ($J_n$ and $J_{n,\lambda_0}$) tests.

4.1. Testing Equality of Unconditional Density Functions with Mixed Data. We consider a range of mixed data DGPs for $f(x,z)$ and $g(x,z)$ designed to examine empirical size and power of the proposed test. We allow $x$ and $z$ to be correlated, and vary the degree of correlation. We let $z \in \{0,1,\ldots,3\}$ with probabilities $(0.125,0.375,0.375,0.125)$. We let $y$ be a mixture of normals drawn from $N(-2,\sigma^2)$ and $N(2,\sigma^2)$ with equal probability. We first draw $z$ and then let $x = \alpha z + y$. When $\alpha = 0$, $z$ and $x$ are independent, while when $\alpha = 3/4$, $\rho_{x,z} = 1/4$. As the bootstrap test is correctly sized, we report size only once, and report power for a range of alternative DGPS. The null DGP is DGP0 for independent $x$ and $z$ where $f(\cdot) = g(\cdot)$. For DGP1, we let the continuous components of $f(\cdot)$ and $g(\cdot)$ differ in their means under the alternative with the difference in means equal to 1/2, with $\rho_{x,z} = 0$. For DGP2, we again let the continuous components of $f(\cdot)$ and $g(\cdot)$ differ in their means under the alternative with the difference in means equal to 1/2, with $\rho_{x,z} = 1/4$. For DGP3, the marginals of $x$ and $z$ are identical under the null and alternative, but the degree of correlation differs under the alternative ($\rho_{x,z} = 1/2$ under the null, $\rho_{x,z} = 1/4$ under the alternative). Finally, for DGP4 the standard deviation of $x$ differs by 1/2 under the null and alternative, while $\rho_{x,z} = 1/4$.\footnote{We are grateful to an anonymous referee for suggesting this rich range of DGPs.}
We consider three tests of the hypothesis $H_0 : g(x, z) = f(x, z) a.e.$: i) the proposed test with least-squares cross-validated $h$ and $\lambda$ ($T_n$), ii) the conventional frequency test with cross-validated $h$ and $\lambda = 0$ ($T_{n,\lambda=0}$), and iii) the conventional ad hoc test with $h = 1.06 \sigma n^{-1/5}$ and $\lambda = 0$ ($T_{n,h=1.06\sigma n^{-1/5},\lambda=0}$). Least-squares cross-validated bandwidth selection is used to obtain $h$ and $\lambda$ for each of the $M = 1,000$ Monte Carlo replications, except where noted. The second order Gaussian kernel is used throughout. For each Monte Carlo replication we conduct $B = 399$ bootstrap replications, and then compute empirical $P$-values for each statistic. We then summarize the empirical rejection frequencies for each test at the 1%, 5%, and 10% levels. We vary the sample size from $n = 50$ through $n = 400$. Empirical size and power over the $M = 1,000$ Monte Carlo replications is summarized in Table 1.

Table 1 suggests the following; i) our test is correctly sized, while the other test sizes are reasonable as well, (ii) the proposed method often possesses substantial power gains, especially in small sample situations relative to the conventional frequency test ($\lambda = 0$) and relative to the ad-hoc test in particular ($h = 1.06 \sigma n^{-1/5}$, $\lambda = 0$), and (iii) the consistency of the tests is evident in the large sample experiments with power approaching one. As the ad-hoc test appears to be slightly undersized, we also computed size-adjusted power and the ranking of estimators in terms of power remains unchanged (the results of size-adjusted power are not reported here to save space).

4.2. Testing Equality of Unconditional Density Functions Under ‘Low Frequency’ Alternatives. First, we consider a Monte Carlo simulation designed to demonstrate how, under a ‘low frequency’ alternative, non-smoothing tests such as the $KS_n$ and $CM_n$ tests can perform better than smoothing tests such as the $T_n$ test proposed in this paper. For what follows, we consider the case where the marginal density for the continuous variable ($x^c$) is a simple univariate normal density function. Here we treat a unimodal normal distribution as a low frequency (i.e., slowly changing) density function.

Specifically, for this experiment the null DGP is $N(0, 1)$ while the alternative DGP is $N(1/2, 1)$. Under the null, both $X$ and $Y$ are drawn from the $N(0, 1)$, while under the alternative $X$ is drawn the $N(0, 1)$ while $Y$ is drawn from the $N(1/2, 1)$. Least-squares cross-validated bandwidth selection is used for the $T_n$ test, and is computed for each of the $M = 1,000$ Monte Carlo replications. For each Monte Carlo replication we conduct $B = 399$ bootstrap replications, and then compute
Table 1. Unconditional density $T_n$ test mixed data Monte Carlo.

<table>
<thead>
<tr>
<th>Size (DGP0)</th>
<th>$T_n$ $T_{n,\lambda=0}$ $T_{n,\lambda=1.05\sigma_n^{-1/5},\lambda=0}$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n$</td>
<td>$\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$</td>
</tr>
<tr>
<td>50</td>
<td>0.013 0.058 0.106 0.011 0.050 0.105 0.010 0.034 0.073</td>
</tr>
<tr>
<td>100</td>
<td>0.006 0.059 0.102 0.006 0.053 0.113 0.008 0.043 0.089</td>
</tr>
<tr>
<td>200</td>
<td>0.009 0.052 0.106 0.007 0.045 0.104 0.005 0.036 0.083</td>
</tr>
<tr>
<td>400</td>
<td>0.013 0.052 0.103 0.010 0.055 0.104 0.008 0.042 0.097</td>
</tr>
</tbody>
</table>

Power (DGP1), mean of $x$ differs under null and alternative, $\rho_{x,z} = 0$

| $T_n$ $T_{n,\lambda=0}$ $T_{n,\lambda=1.05\sigma_n^{-1/5},\lambda=0}$ |
|--------------| -------------------------------------------------------------------------------------------------|
| $n$          | $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ |
| 50           | 0.045 0.146 0.257 0.026 0.102 0.187 0.016 0.067 0.148 |
| 100          | 0.093 0.255 0.389 0.062 0.195 0.316 0.042 0.146 0.265 |
| 200          | 0.221 0.499 0.646 0.174 0.424 0.592 0.143 0.374 0.531 |
| 400          | 0.633 0.867 0.936 0.564 0.830 0.920 0.534 0.850 0.925 |

Power (DGP2), mean of $x$ differs under null and alternative, $\rho_{x,z} = 1/4$

| $T_n$ $T_{n,\lambda=0}$ $T_{n,\lambda=1.05\sigma_n^{-1/5},\lambda=0}$ |
|--------------| -------------------------------------------------------------------------------------------------|
| $n$          | $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ |
| 50           | 0.036 0.135 0.218 0.029 0.120 0.213 0.017 0.085 0.162 |
| 100          | 0.066 0.216 0.343 0.064 0.208 0.340 0.039 0.154 0.280 |
| 200          | 0.176 0.437 0.605 0.158 0.435 0.604 0.121 0.389 0.551 |
| 400          | 0.532 0.830 0.913 0.508 0.828 0.909 0.485 0.805 0.906 |

Power (DGP3), marginals identical under null and alternative, correlation differs under alternative

| $T_n$ $T_{n,\lambda=0}$ $T_{n,\lambda=1.05\sigma_n^{-1/5},\lambda=0}$ |
|--------------| -------------------------------------------------------------------------------------------------|
| $n$          | $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ |
| 50           | 0.059 0.192 0.308 0.042 0.165 0.276 0.023 0.118 0.215 |
| 100          | 0.111 0.325 0.492 0.105 0.301 0.466 0.058 0.247 0.386 |
| 200          | 0.335 0.647 0.767 0.316 0.629 0.768 0.250 0.575 0.731 |
| 400          | 0.811 0.964 0.989 0.815 0.968 0.986 0.776 0.967 0.985 |

Power (DGP4), standard deviations of $x$ differ under the null and alternative, $\rho_{x,z} = 1/4$

| $T_n$ $T_{n,\lambda=0}$ $T_{n,\lambda=1.05\sigma_n^{-1/5},\lambda=0}$ |
|--------------| -------------------------------------------------------------------------------------------------|
| $n$          | $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ $\alpha = 0.01$ $\alpha = 0.05$ $\alpha = 0.10$ |
| 50           | 0.028 0.122 0.221 0.030 0.116 0.203 0.014 0.067 0.131 |
| 100          | 0.068 0.217 0.343 0.067 0.219 0.328 0.036 0.140 0.246 |
| 200          | 0.164 0.413 0.569 0.165 0.409 0.567 0.089 0.304 0.490 |
| 400          | 0.461 0.788 0.895 0.485 0.785 0.892 0.388 0.732 0.867 |

empirical $P$-values for each statistic. We then summarize the empirical rejection frequencies for each test at the 1%, 5%, and 10% levels. We vary the sample size from $n = 50$ through $n = 400$ at which point power is equal to one for all three tests considered. Results are reported in Table 2 below.

Table 2 reveals that each test is correctly sized while power is highest for the non-smoothed tests as expected under ‘low frequency’ alternatives with the $CM_n$ test being most powerful for this DGP. We also computed the non-smoothing test $I_{n,h=1}$ (in fact $I_{n,h=\sigma}$) in our simulations, the
Table 2. Monte Carlo comparison of the $T_n$, $CM_n$, and $KS_n$ tests (low frequency data).

<table>
<thead>
<tr>
<th>Size (DGP0)</th>
<th>$T_n$</th>
<th>$CM_n$</th>
<th>$KS_n$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n = 0$</td>
<td>0.01</td>
<td>0.01</td>
<td>0.01</td>
</tr>
<tr>
<td>50</td>
<td>0.013</td>
<td>0.011</td>
<td>0.018</td>
</tr>
<tr>
<td>100</td>
<td>0.014</td>
<td>0.011</td>
<td>0.011</td>
</tr>
<tr>
<td>200</td>
<td>0.015</td>
<td>0.006</td>
<td>0.009</td>
</tr>
<tr>
<td>400</td>
<td>0.007</td>
<td>0.014</td>
<td>0.011</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>Power (DGP1)</th>
<th>$T_n$</th>
<th>$CM_n$</th>
<th>$KS_n$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n = 0$</td>
<td>0.01</td>
<td>0.01</td>
<td>0.01</td>
</tr>
<tr>
<td>50</td>
<td>0.170</td>
<td>0.344</td>
<td>0.297</td>
</tr>
<tr>
<td>100</td>
<td>0.440</td>
<td>0.719</td>
<td>0.643</td>
</tr>
<tr>
<td>200</td>
<td>0.795</td>
<td>0.999</td>
<td>0.936</td>
</tr>
<tr>
<td>400</td>
<td>0.999</td>
<td>1.000</td>
<td>0.999</td>
</tr>
</tbody>
</table>

results show that the $I_{n,h=1}$ test is corrected sized and has power similar to that of the $KS_n$ test, the detailed results are not reported here due to space limitation.

4.3. Testing Equality of Unconditional Density Functions Under ‘High Frequency’ Alternatives. Next, we consider a Monte Carlo simulation designed to demonstrate how, under ‘high frequency’ alternatives, smoothing tests such as the $T_n$ proposed in this paper can perform better than non-smoothing tests such as the $KS_n$ and $CM_n$ tests, a fact that may not be appreciated by many readers.

We shall draw data from a mixture of normal distributions each having different location and scales. Under the null, $f(x)$ (and $g(x)$) is a mixture of two normal distributions: $N(-1/2, 1)$ and $N(1/2, 4)$, with equal probability. That is, we draw data for $X$ and $Y$ from a $N(-1/2, 1)$ and $N(1/2, 4)$ with equal probability. It can be seen that the PDF of data drawn from this mixture has a bimodal and asymmetric distribution, the left peak being higher than the right. Under the alternative, however, $f(x)$ remains the same as above, while $g(x)$ is again a mixture but of two different normal distributions: $N(-1/2, 4)$ and $N(1/2, 1)$ with equal probability. That is, we reverse the peaks and draw data for $Y$ from a $N(-1/2, 4)$ and $N(1/2, 1)$ with equal probability. All remaining particulars are the same as for the Monte Carlo setting considered above. Results are reported in Table 3.
Table 3. Monte Carlo comparison of the $T_n$, $CM_n$, and $KS_n$ tests (high frequency data).

<table>
<thead>
<tr>
<th>Size (DGP0)</th>
<th>$T_n$</th>
<th>$CM_n$</th>
<th>$KS_n$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n$</td>
<td>$\alpha = 0.01$</td>
<td>$\alpha = 0.05$</td>
<td>$\alpha = 0.10$</td>
</tr>
<tr>
<td>50</td>
<td>0.012</td>
<td>0.053</td>
<td>0.100</td>
</tr>
<tr>
<td>100</td>
<td>0.006</td>
<td>0.047</td>
<td>0.091</td>
</tr>
<tr>
<td>200</td>
<td>0.006</td>
<td>0.059</td>
<td>0.106</td>
</tr>
<tr>
<td>400</td>
<td>0.009</td>
<td>0.050</td>
<td>0.095</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>Power (DGP1)</th>
<th>$T_n$</th>
<th>$CM_n$</th>
<th>$KS_n$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n$</td>
<td>$\alpha = 0.01$</td>
<td>$\alpha = 0.05$</td>
<td>$\alpha = 0.10$</td>
</tr>
<tr>
<td>50</td>
<td>0.107</td>
<td>0.269</td>
<td>0.395</td>
</tr>
<tr>
<td>100</td>
<td>0.214</td>
<td>0.452</td>
<td>0.586</td>
</tr>
<tr>
<td>200</td>
<td>0.539</td>
<td>0.756</td>
<td>0.850</td>
</tr>
<tr>
<td>400</td>
<td>0.895</td>
<td>0.986</td>
<td>0.994</td>
</tr>
</tbody>
</table>

Table 3 reveals that each test is correctly sized while power is highest for the smoothing tests as expected under ‘high frequency’ alternatives with the $T_n$ test being most powerful for this DGP. Note that our so-called ‘high frequency’ density function only has two modes. Simulation results (not reported here to save space) show that, for density functions with more than two modes, a smoothing test enjoys even more power gains relative to non-smoothing tests. Therefore, smoothing tests are complement to non-smoothing tests and should be part of all applied researchers’ standard toolkit.

4.4. Testing Equality of Conditional Density Functions with Mixed Data. Finally, we consider a Monte Carlo simulation designed to examine the finite-sample performance of the proposed $J_{n,\lambda_0}$ and $J_n$ tests and compare them to a counterpart that smooths the continuous variable with a ad-hoc $h$ and uses the frequency indicator function for the discrete conditional variable.

Under the null of $f(x|z) = g(x|z)$ we let the discrete variable, $z$, assume four values with equal probability, $\{0, 1, \ldots, 3\}$. Next, we create $Y = z/4 + N(0, 1)$ under the null so that $f(y|z) \sim N(z/4, 1)$. Under the alternative $Y = z/4 + N(1/2, 1)$ so that $g(Y|z) \sim N(z/4 + 1/2, 1)$. All remaining particulars are the same as for the Monte Carlo setting considered above. Results are reported in Table 4.

Table 4 reveals that the proposed tests $J_{n,\lambda_0}$ and $J_n$ are correctly sized and both are more powerful than its ad-hoc/frequency-based counterpart $J_{n,h}=1.06\sigma n^{-1/5}$. Also, the $J_{n,\lambda_0}$ test is more
powerful than the $J_n$ test due to the fact that $J_{n,0}$ also smooths the conditional discrete variable $z$.

5. Earnings, Educational Attainment, and Wage Gaps

There exists a large literature in labour economics regarding how returns to a college education is instrumental in understanding the widening wage gap in the US economy. The first step of such analysis, however, would involve determining whether or not statistically significant differences between joint distributions defined over both continuous (income) and discrete (educational attainment) variables exist.

For what follows, we consider data spanning the years 1980-2000 constructed from the US Current Population Survey (CPS) March supplement on real incomes for white non-Hispanic workers aged 25 to 55 years who were full-time workers working at least 30 hours a week and at least 40 weeks a year. Self-employed, farmers, unpaid family workers, and members of the Armed Forces are excluded. We consider the distribution of income for high school versus college graduates. Wage income is the income category considered, and figures are expressed in 2000 dollars.

Figure 1 presents kernel smoothed PDF estimates for income by year and educational attainment. Table 5 presents various moments for the income data, namely, measures of location and scale by year and educational attainment. We observe that average/median income for both high school
and college graduates is lower in 2000 than it was in 1980. The interquartile range (IQR) has fallen for both groups from 1980 to 2000, while standard deviations have increased.

**Table 5.** Income location and scale summaries by year and educational attainment.

<table>
<thead>
<tr>
<th></th>
<th>1980</th>
<th>2000</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>High School</td>
<td>College</td>
</tr>
<tr>
<td>Mean</td>
<td>$20,637.72</td>
<td>$22,838.88</td>
</tr>
<tr>
<td>Median</td>
<td>$18,880.17</td>
<td>$20,661.16</td>
</tr>
<tr>
<td>Stdev</td>
<td>$10,720.27</td>
<td>$11,767.74</td>
</tr>
<tr>
<td>IQR</td>
<td>$14,709.37</td>
<td>$14,876.03</td>
</tr>
</tbody>
</table>

As noted in Section 1, moment-based tests, which only compare a finite number of moments from two distributions, are not consistent tests. By way of example, we test whether the joint distribution of earnings and educational attainment differ over time. We select two random samples, one for the year 1980 and one for the year 2000, each of size $n_1 = n_2 = 1,000$, and apply the unconditional $T_n$ and conditional $J_n$ tests. We obtain $\hat{T}_n = 87.15$ with an associated bootstrap $P$-value of $P < 0.001$, while $\hat{J}_n = 54.63$ with an associated bootstrap $P$-value of $P < 0.001$. This suggests that there are
indeed significant differences in the joint distribution of income and educational attainment between 1980 and 2000, and that there are significant differences in the distribution of income conditional upon educational attainment between 1980 and 2000.

6. Conclusion

We consider the problem of testing for equality of two density or two conditional density functions defined over mixed discrete and continuous data. Smoothing parameters are chosen via least squares cross-validation, and we smooth both the discrete and continuous variables in a particular manner. We advocate the use of bootstrap methods for obtaining the statistic’s null distribution in finite-sample settings. Simulations show that the proposed tests enjoy power gains relative to both a conventional frequency-based test and a smoothing test based on ad hoc smoothing parameter selection. An application to testing for the equality of the joint distribution of income and educational attainments underscores the novelty and flexibility of the proposed approach in mixed data settings.

Our approach can be extended to testing the equality of two residual distributions. Hall et al. (2004) have shown that the cross-validation method has the remarkable ability of potentially removing irrelevant conditioning variables. In the testing framework we expect that this will lead to a more powerful test relative to peers that lack this ability. In this paper we only consider the case where the discrete variable has finite support. Extension of our approach to allow for the discrete variable to be countably infinite will be a fruitful avenue for further investigation. We leave the exploration of these topics for future research.

Appendix A. Proofs of Theorems

Proof of Theorem 2.1. The test statistic $I_n$ can be written as $I_n = I_{1n} + I_{2n}$, where

$$I_{1n} = -\frac{2}{n_1 n_2} \min\{n_1, n_2\} \sum_{i=1}^{\min\{n_1, n_2\}} K_{\gamma,x_i,y_i}$$

and

$$I_{2n} = \sum_i \sum_{j \neq i} \left[ \frac{1}{n_1^2} K_{\gamma,x_i,x_j} + \frac{1}{n_2^2} K_{\gamma,y_i,y_j} - \frac{1}{n_1 n_2} (K_{\gamma,x_i,y_j} + K_{\gamma,x_j,y_i}) \right],$$
where \( \sum_i = \sum_{i=1}^{n_1} \) if the summand contains \( x_i \), \( \sum_1 = \sum_{i=1}^{n_2} \) if the summand contains \( y_i \). For example, \( \sum_i \sum_{j \neq i} K_{\gamma_i,j} x_i y_j = \sum_{i=1}^{n_1} \sum_{j \neq i}^{n_2} K_{\gamma_i,j} x_i y_j \), and \( \sum_i \sum_{j \neq i} K_{\gamma_j,i} x_i y_j = \sum_{i=1}^{n_2} \sum_{j \neq i}^{n_1} K_{\gamma_j,i} x_i y_j \).

Let \( n = \min\{n_1, n_2\} \) and let \( m(x,y) \) denote the joint density of \((X_i^c, Y_i^c)\). By noting that

\[
J_{L_i;j} x_i y_i j = \frac{K_{\gamma_i,j} x_i y_j}{}\]

we have that

\[
E[J_{K_{\gamma_i,j} x_i y_j}] = (h_1 \ldots h_q)^{-1} \int W((y_i^c - x_i^c)/h)m(x_i^c, y_i^c)dx_i^c dy_i^c = \int W(v)m(x_i^c, x_i^c + hv)dv dx_i^c = O(1).
\]

Then it follows that

\[
E[I_{1n}] = O_p(n^{-1}).
\]

Note that we obtain the above result by allowing for arbitrary correlation between \( X_i \) and \( Y_i \).

For example, for panel data with two time periods, \( X_i \) and \( Y_i \) can be repeated measures from the same individual \( i \) over two different time periods.

Next, we consider \( I_{2n} \). We will write \( A_n = B_n + (s.o.) \) to mean that \( B_n \) is the leading term of \( A_n \), while \( (s.o.) \) denotes terms having orders smaller than \( B_n \). Let \( Z_i = (X_i, Y_i) \) and define

\[
H_n(Z_i, Z_j) = K_{\gamma_i,j} x_i x_j + K_{\gamma_i,j} y_i y_j - K_{\gamma_i,j} x_i y_j - K_{\gamma_i,j} y_i x_j.
\]

For \( i \neq j \), we have

\[
E[H_n(Z_i, Z_j)|Z_i] = E[K_{\gamma_i,j} x_i x_j | X_i] - E[K_{\gamma_i,j} y_i y_j | X_i] + E[K_{\gamma_i,j} y_i x_j | Y_i] - E[K_{\gamma_i,j} x_i y_j | Y_i] = 0,
\]

where \( \int dx = \sum_x f(x)dx \), which follows because

\[
E[K_{\gamma_i,j} x_i x_j | X_i] = \int K_{\gamma_i,j} x_j f(x_j)dx_j = \int K_{\gamma_i,j} y_j g(y_j)dy_j = E[K_{\gamma_i,j} | X_i]
\]

since \( f(\cdot) = g(\cdot) \) under \( H_0 \).
Therefore, $I_{2n}$ is a degenerate U-statistic. Defining $H = h_1 \ldots h_q$, then it is easy to show that

$$\text{var}(I_{2n}) = E[(I_{2n})^2]$$

$$= 2 \sum_i \sum_{j \neq i} \{ n_1^{-4} E[(K_{\gamma,x_i,x_j})^2] + n_2^{-4} E[(K_{\gamma,y_i,y_j})^2] \}
+(n_1n_2)^{-2} E[(K_{\gamma,x_i,y_j})^2] + (n_1n_2)^{-2} E[(K_{\gamma,x_j,y_i})^2] + (s.o.) \}
$$

$$= \frac{2}{n_1n_2} \{ \delta_n^{-1} E [(K_{\gamma,x_i,x_j})^2] + \delta_n E [(K_{\gamma,y_i,y_j})^2] + E [(K_{\gamma,x_i,y_j})^2] + E [(K_{\gamma,x_j,y_i})^2] + (s.o.) \}
$$

$$= \frac{2}{n_1n_2} \left\{ \delta_n^{-1} + \delta_n + 2 \left[ E[f(X_i)] \left[ \int W^2(v)dv \right] + o(1) \right] \right\}
$$

$$\equiv (n_1n_2H)^{-1} \{ \sigma_0^2 + o(1) \},$$

where $\delta_n = n_1/n_2$, and $\sigma_0^2 = 2[\delta^{-1} + \delta + 2][E[f(X_i)][\int W^2(v)dv]$, and we have used

$$E[(K_{\gamma,x_i,x_j})^2] = H^{-2} \sum_{x_i^c, x_j^c} \int W^2((x_j^c - x_i^c)/h)L^2(x_i^d, x_j^d, \lambda)f(x_i)f(x_j)dx_i^c dx_j^c$$

$$= H^{-1} \sum_{x_i^d, x_j^d} \int W^2(v)L^2(x_i^d, x_j^d, \lambda)f(x_i)f(x_j + hv, x_j^d)dx_i^d dv$$

$$= H^{-1} \left\{ \left[ \sum_{x_i^d} \int f(x_i^c, x_i^d)^2 dx_i^d \right] \left[ \int W^2(v)dv \right] + O(|h|^2 + |\lambda|) \right\}
$$

$$= H^{-1} \left\{ E[f(X_i)] \left[ \int W^2(v)dv \right] + o(1) \right\},$$

with $H = h_1 \ldots h_q$, where $|h|^2 = \sum_{s=1}^q h_s^2$ and $|\lambda| = \sum_{s=1}^q \lambda_s$. Similarly, $E[(K_{\gamma,y_i,y_j})^2], E[(K_{\gamma,x_i,y_j})^2], E[(K_{\gamma,x_j,y_i})^2]$ all equal $H^{-1} \{ E[f(X_i)][\int W^2(v)dv] + o(1) \}$ under $H_0$ (since $E[f(X_i)] = E[g(Y_i)]$ under $H_0$).

It is straightforward, though tedious, to check that the conditions of Hall’s (1984) central limit theorem for degenerate U-statistics holds. Thus, under $H_0$ we have

(A.2) \hspace{1cm} (n_1n_2H)^{1/2}I_{2n}/\sigma_0 \rightarrow N(0, 1) \text{ in distribution.}
Note that $E(\sigma_n^2) = \sigma_0^2 + o(1)$ ($\sigma_n^2$ is defined in Theorem 2.1), and by the U-statistic H-decomposition, it follows that $\sigma_n^2 = E(\sigma_n^2) + o_p(1) = \sigma_0^2 + o_p(1)$. Therefore, from (A.2) we obtain

(A.3) \( (n_1 n_2 H)^{1/2} I_{2n}/\sigma_n \to N(0, 1) \) in distribution.

In the above proof we implicitly assumed that the support of $x^c$ is unbounded since we did not address the possible boundary bias problem. Below we show that, in fact, the above result holds true even when $x^c$ has bounded support and the density is bounded below by a positive constant in its support. First, it is obvious that $E[H_n(Z_i, Z_j)|Z_i] = 0$ under $H_0$, because this follows from $f(\cdot) = g(\cdot)$, regardless of whether $x^c$ has bounded or unbounded support. Next, we show that $\text{var}(I_{2n}) = (n_1 n_2 H)^{-1}\{\sigma_0^2 + o(1)\}$ also holds true. It suffices to show that $E[(K_{i,x_i,x_j})^2] = \{E[f(X_i)][\int W^2(v)dv + o(1)]\}$. For expositional simplicity, we will only consider the univariate $x^c$ case (and without $x^d$) where $x^c$ is uniformly distributed in $[a, b]$ for some constants $a, b$ with $b > a$. The proof for the general case is similar but much more tedious. Now let $\alpha \in (0, 1)$ be a constant. Then we have

\[
E[(W_{h,x_i,x_j})^2] = H^{-2} \int_a^b \int_a^b W^2((x_j^c - x_i^c)/h) f(x_i^c) f(x_j^c) dx_i^c dx_j^c \\
= H^{-1} \int_a^b \int_{(a-x_i)/h}^{(b-x_i)/h} W^2(v) f(x_i^c) f(x_i^c + hv) dv dx_i^c \\
= H^{-1} \left[ \int_{a-h^\alpha}^{a+|h^\alpha|} f(x_i^c) f(x_i^c + hv) dv dx_i^c + \int_{b-h^\alpha}^{b+|h^\alpha|} f(x_i^c) f(x_i^c + hv) dv dx_i^c \right] \\
= H^{-1} \left[ \int_{a-h^\alpha}^{b-h^\alpha} f(x_i^c) f(x_i^c + hv) dv dx_i^c + \text{(s.o.)} \right] \\
= H^{-1} \left\{ \int_{a-h^\alpha}^{b-h^\alpha} f(x_i^c)^2 dx_i^c \right\} \left[ \int_0^{\infty} W^2(v) dv \right] + o(1) \\
= H^{-1} \left\{ E[f(X_i^c)] \right\} \int_0^{\infty} W^2(v) dv + o(1),
\]

where we have used the fact that for $x_i \in (a + h^\alpha, b - h^\alpha)$, the interval $((a - x_i)/h, (b - x_i)/h) \supset (-h^\alpha/h, h^\alpha/h)$, which expands to $(-\infty, \infty)$ as $h \to 0$ since $0 < \alpha < 1$. Hence, $\int_{(a-x_i)/h}^{(b-x_i)/h} W(v^2) dv \to \int_0^{\infty} W^2(v) dv$ as $h \to 0$. The basic idea underlying the proof above is as follows: We divide $[a, b]$ into three intervals: $[a, a + h^\alpha]$, $(a + h^\alpha, b - h^\alpha)$, and $[b - h^\alpha, b]$. Compared with the second interval, the first and third intervals are negligible since their lengths shrink to zero as $n \to \infty$. 

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The second interval does not have a boundary problem as the boundary regions lie in the first and third intervals. The testing problem we consider here is different from pointwise estimation which may suffer from boundary bias issues. The leading term of test statistic has zero mean and its asymptotic variance has the same expression, regardless of the nature of the support of \( x^c \). Hence, the asymptotic null (normal) distribution presented in Theorem 2.1 is invariant to the nature of the support of \( x^c \).

Summarizing the above, (A.1) and (A.3) complete the proof of Theorem 2.1.

Before we prove Theorem 2.2, we first present a lemma which will be used in the proof of Theorem 2.2.

**Lemma A.1.** Let \( A_n(c) = (n_1n_2h_1 \ldots h_q)^{1/2}I_{2n}(h, \lambda) \), where \( h_s = a_sn^{-\zeta}, \lambda_s = b_sn^{-2\zeta}, c = (a_1, \ldots, a_q, b_1, \ldots, b_r), c_s \in [C_{1s}, C_{2s}] \) with \( 0 < C_{1s} < C_{2s} < \infty \) \((s = 1, \ldots, q + r)\).

Then the stochastic process \( A_n(c) \) indexed by \( c \) is tight under the sup-norm.

**Proof.** Writing \( K_{\gamma,ij} \) as \((h_1 \ldots h_q)^{-1}K_{c,ij}\) with \( h_s = a_sn^{-\zeta} \) and \( \lambda_s = b_2n^{-2\zeta} \), where \( K_{c,ij} = W(\frac{X_i - X_j}{h})L(X_i^d, X_j^d, \lambda), \) and letting \( \delta = q\zeta, H^{-1/2} = (h_1 \ldots h_q)^{-1/2}, C_1 = (a_1, \ldots, a_q)^T, C_2 = (b_1, \ldots, b_r)^T \) (where the superscript \( T \) denotes transpose), \( \tilde{C}_1 = \prod_{s=1}^qa_s, \) and \( \tilde{C}_2 = \prod_{s=1}^rb_s \). Then we have \( H^{-1/2}K_{c,ij} = \tilde{C}_1n^{\delta/2}W_{C_1,ij}L_{C_2,ij} \). Let \( C_2' = (b_1', \ldots, b_r')^T \), then by noting that \( |L_{C_2',ij} - L_{C_2,ij}| \leq d_1\sum_{s=1}^r|b_s - b_s'| \leq d_2||C_2 - C_2'|| \), where \( d_1 \) and \( d_2 \) are some finite positive constants.

Note that \( a_s \) are all bounded below by some positive constants, we have for all \( s = 1, \ldots, q \) that

\[
\left| \frac{X_{is} - X_{js}}{h_s} - \frac{X_{is} - X_{js}}{h_s'} \right| \leq \left| X_{is} - X_{js} \right| \left| \frac{1}{h_s} - \frac{1}{h_s'} \right| \leq \left| X_{is} - X_{js} \right| \left| \frac{h'}{h_s} - \frac{h'}{h_s'} \right| \leq \left| X_{is} - X_{js} \right| \left| \frac{a_s'}{a_s} - \frac{a_s'}{a_s} \right| \leq d_s \left| X_{is} - X_{js} \right| \left| \frac{a_s'}{a_s} - a_s \right|,
\]

(A.4)

where \( d_s \) is a finite positive constant (because \( a_s \) is bounded below and above by some positive constants).
By the Lipschitz condition (see (C2)) on the univariate kernel function that \( |w(u) - w(v)| \leq \xi(v)|u - v| \), it is easy to show that the product kernel also satisfies a Lipschitz condition, namely,

\[(A.5) \quad |W(u) - W(v)| \leq M(v)||u - v||,\]

where \( M(v) = c[\sum_{s=1}^{q} \xi(v_s)] \) and where \( c \) is a positive constant such that \( \sup_{v\in\mathbb{R}} w(v)q^{-1} \leq c \).

(A.4) and (A.5) yield

\[(A.6) \quad |W_{C',x_i,x_j} - W_{C,x_i,x_j}| \leq d M_{C,x_i,x_j} \left[ \left| \frac{X_i - X_j}{h} \right| \times ||C'_1 - C_1||, \right]\]

where \( d \) is a positive constant, and

\[M_{C_i,j} = M((X_j - X_i)/h) = c[\sum_{s=1}^{q} \xi((X_j - X_i)/h_s)].\]

Using (A.6) we have

\[(A.7) \quad \left| (H')^{-1/2} K_{C',ij} - H^{-1/2} K_{C,ij} \right| = \left| n^{\delta/2} \left\{ \left( \tilde{C}'_1 \right)^{-1/2} W_{C'_1,ij} L_{C'_2,ij} - \tilde{C}'_1^{-1/2} W_{C_1,ij} L_{C_2,ij} \right\} \right|\]

\[= \left| n^{\delta/2} \left\{ (\tilde{C}'_1)^{-1/2} W_{C'_1,ij} L_{C'_2,ij} - L_{C'_2,ij} + (\tilde{C}'_1)^{-1/2} W_{C'_1,ij} - \tilde{C}'_1^{-1/2} W_{C'_1,ij} \right\} L_{C_2,ij} \right|\]

\[\leq D_1 \left\{ (H')^{-1/2} W_{C'_1,ij} ||C'_2 - C_2|| + H^{-1/2} M_{C_1,ij} \left[ \left| \frac{X_j - X_i}{h} \right| \times ||C'_1 - C_1|| \right] \right\},\]

where \( D_1 > 0 \) is a finite constant. In the last equality we used \( |L_{C_2,ij}| \leq 1 \) and assumption (C3), and we also replaced one of the \( (\tilde{C}'_1)^{-1/2} \) by \( \tilde{C}'_1^{-1/2} \) because \( a_s \in [C_{1s}, C_{2s}] \) are all bounded from above and below. The difference can be absorbed into \( D_1 \).
By noting that $A_n(c') - A_n(c)$ is a degenerate U-statistic, and using (A.7), we have

$$
E \left\{ \left[ A_n(c') - A_n(c) \right]^2 \right\} \\
= E \left\{ \left[ (H')^{-1/2} K_{c',ij} - H^{-1/2} K_{c,ij} \right]^2 \right\} \\
\leq 4D_2 E \left\{ \left[ (H')^{-1} W_{C_2,ij}^2 ||C'_2 - C_2||^2 + H^{-1} M_{C,i}^2 \left\| \frac{X_j - X_i}{h} \right\| \left\| C'_1 - C_1 \right\| \right] \right\} \\
\leq D_3 \left\{ \int \int f(x_i)f(x_i + hu)W^2(u)dx_i du \right\} \left\| C'_2 - C_2 \right\|^2 + \left\{ \int \int f(x_i)f(x_i + v)M^2(v)||v||^2 dx_i dv \right\} \left\| C'_1 - C_1 \right\|^2 \\
\leq D_4 \sup_x f(x) \left\{ \int W^2(u)du \right\} \left\| C'_2 - C_2 \right\|^2 + \left\{ \int M^2(v)||v||^2 dv \right\} \left\| C'_1 - C_1 \right\|^2 \\
(A.8) \leq D_5 \left\| C' - C \right\|^2,
$$

where $D_j (j = 3, 4, 5)$ are some finite positive constants. Therefore, $A_n(\cdot)$ (hence, $B_n(\cdot)$) is tight by Theorem 15.6 of Billingsley (1968, p. 128), or Theorem 3.1 of Ossiander (1987).

**Proof of Theorem 2.2.** Theorem 2.1 implies that when $h_s = h_s^0 = a_s^0 n^{-\zeta}$ and $\lambda_s = \lambda_s^0 = b_s^0 n^{-2\zeta}$, the test statistic $\hat{T}_n(h^0, \lambda^0) \to N(0,1)$ in distribution. Therefore, it is sufficient to prove that $\hat{T}_n(\hat{h}, \hat{\lambda}) - \hat{T}_n(h_0, \lambda_0) = o_p(1)$. For this, it suffices to show the following:

(i) $(n_1 n_2 \hat{h}_1 \ldots \hat{h}_q)^{1/2} \hat{I}_{2n} = (n_1 n_2 h_1^0 \ldots h_q^0)^{1/2} I_{2n} + o_p(1),$

(ii) $(n_1 n_2 \hat{h}_1 \ldots \hat{h}_q)^{1/2} \hat{I}_{1n} = (n_1 n_2 h_1^0 \ldots h_q^0)^{1/2} I_{1n} + o_p(1),$ and

(iii) $\hat{\sigma}_n^2 = \hat{\sigma}_0^2 + o_p(1)$, $\hat{\sigma}_0^2$ is defined in Theorem 2.1 but with $(h_1, \ldots, h_q, \lambda_1, \ldots, \lambda_r)$ replaced by $(h_1^0, \ldots, h_q^0, \lambda_1^0, \ldots, \lambda_r^0)$.

Below we will only prove (i) since (ii) and (iii) are much easier to establish than (i) (and can be similarly proved).

Write $\hat{h}_s = \hat{a}_s n^{-\zeta}$ and $\hat{\lambda}_s = \hat{b}_s n^{-2\zeta}$. From Theorem 3.1 of Li & Racine (2003), we know that $\hat{h}_s/h_s^0 - 1 \to 0$ and $\hat{\lambda}_s/\lambda_s^0 - 1 \to 0$ (in probability). This implies that $\hat{a}_s \to a_s^0$ and $\hat{b}_s \to b_s^0$ in probability. Let $C = \prod_{s=1}^q [a_{1s}, a_{2s}] \times \prod_{t=1}^r [b_{1t}, b_{2t}]$, where $a_{js}$ and $b_{jt}$ ($j = 1, 2$) are some positive constants with $a_{1s} < a_{s}^0 < a_{2s}$ ($s = 1, \ldots, q$) and $b_{1t} < b_{t}^0 < b_{2t}$ ($t = 1, \ldots, r$). Let $c = (a_1, \ldots, a_q, b_1, \ldots, b_r)$, $\hat{c} = (\hat{a}_1, \ldots, \hat{a}_q, \hat{b}_1, \ldots, \hat{b}_r)$. Then Lemma
A.1 shows that $A_n(c) \equiv (n_1 n_2 h_1 \ldots h_q)^{1/2} I_{2n}(h, \lambda)$ (with $h_s = a_s n^{-c}$ and $\lambda_s = b_s n^{-2c}$) is tight in $c \in \mathcal{C}$.

Define $B_n(c) = A_n(c) - A_n(c_0)$. Then (i) becomes $B_n(\hat{c}) = o_p(1)$, i.e., we want to show that, for all $\epsilon > 0$,

$$(A.9) \quad \lim_{n \to \infty} Pr[|B_n(\hat{c})| < \epsilon] = 1.$$ 

For any $\delta > 0$, denote the $\delta$-ball centered at $c_0$ by $C_\delta = \{c : |c - c_0| \leq \delta\}$, where $|.|$ denotes the Euclidean norm of a vector. By Lemma A.1 we know that $A_n(\cdot)$ is tight. By the Arzela-Ascoli Theorem (see Theorem 8.2 of Billingsley (1968, p. 55)) we know that tightness implies the following stochastic equicontinuous condition: for all $\epsilon > 0$, $\eta_1 > 0$, there exist a $\delta (0 < \delta < 1)$ and an $N_1$, such that

$$(A.10) \quad Pr\left[\sup_{||c' - c|| < \delta} |A_n(c') - A_n(c)| > \epsilon \right] < \eta_1$$

for all $n \geq N_1$.

(A.10) implies that

$$(A.11) \quad Pr[|B_n(\hat{c})| > \epsilon, \hat{c} \in C_\delta] \leq Pr\left[\sup_{c \in C_\delta} |B_n(c)| > \epsilon \right] < \eta_1$$

for all $n \geq N_1$.

Also, from $\hat{c} \to c_0$ in probability we know that, for all $\eta_2 > 0$ and for the $\delta$ given above, there exists an $N_2$ such that

$$(A.12) \quad Pr[\hat{c} \notin C_\delta] \equiv Pr[||\hat{c} - c_0|| > \delta] < \eta_2$$

for all $n \geq N_2$.

Therefore,

$$Pr[|B_n(\hat{c})| > \epsilon] = Pr[|B_n(\hat{c})| > \epsilon, \hat{c} \in C_\delta] + Pr[|B_n(\hat{c})| > \epsilon, \hat{c} \notin C_\delta]$$

$$(A.13) \quad < \eta_1 + \eta_2$$
for all $n \geq \max\{N_1, N_2\}$ by (A.11) and (A.12), where we have also used the fact that $\{|B_n(\hat{c})| > \epsilon, \hat{c} \not\in C_3\}$ is a subset of $\{\hat{c} \not\in C_3\}$ (If $A$ is a subset of $B$, then $P(A) \leq P(B)$).

(A.13) is equivalent to (A.9). This completes the proof of (i).

Proof of Theorem 2.3. Because the cumulative distribution function for the standard normal random variable is a continuous distribution, by Polyá’s Theorem (Bhattacharya & Rao (1986)), we know that (3.22) is equivalent to (for a given value of $z$),

$$P(\hat{T}_n^* \leq z | \{X_i, Y_i\}_{i=1}^n) - \Phi(z) = o_p(1).$$

First, we can write $\hat{T}_n^* = \hat{T}_{1n}^* + \hat{T}_{2n}^*$, where $\hat{T}_{jn}^*$ is the same as in $\hat{T}_{jn}$ ($j = 1, 2$) except that $X_i$ ($Y_i$) is replaced by $X_i^*$ ($Y_i^*$) and $(h, \lambda)$ is replaced by $(\hat{h}, \hat{\lambda})$. Let $E^*(\cdot)$ denote $E(\cdot | \{X_i\}_{i=1}^n, \{Y_i\}_{i=1}^n)$. By exactly the same arguments as we used in the proof of Theorem 2.1, one can show that $\hat{T}_{1n}^* = O_p(n^{-1})$ by showing that $E^*|\hat{T}_{1n}^*| = O_p(n^{-1})$. Also, one can show that $\hat{T}_{2n}^* - \hat{T}_{2n} = o_p((n^2H)^{-1/2})$ (by showing that $E^*[\hat{T}_{2n}^* - \hat{T}_{2n}]^2 = o_p((n^2H)^{-1}) (H = h_1 \ldots h_q)$, and that $\hat{\sigma}_n^2 - \hat{\sigma}_n^2 = o_p(1)$. Therefore, we have that

$$\left(\sum_{k=1}^{q} \hat{h}_k \ldots \hat{h}_q \right)^{1/2} \hat{T}_{n}^*/\hat{\sigma}_n^* - \left(\sum_{k=1}^{q} \hat{h}_k \ldots \hat{h}_q \right)^{1/2} \hat{T}_{1n}/\hat{\sigma}_{1n} = o_p(1).$$

Thus, (A.14) (and hence, Theorem 2.3) follows from (A.15) Theorem 2.1. \hfill \qed

Proof of Theorem 3.1. We know that $\hat{h}_s = h_s^0 + o_p(h_s^0)$ for $s = 1, \ldots, q$; and $\hat{\lambda}_s = \lambda_s^0 + o_p(\lambda_s^0)$ for $s = 1, \ldots, r$, where $h_s^0 = a_s^0 n^{-\zeta}$ and $\lambda_s^0 = b_s^0 n^{-2\zeta}$ for some $\zeta > 0$. We will only prove Theorem 3.1 for the non-stochastic smoothing parameter case, i.e., for $(h_1, \ldots, h_q, \lambda_1, \ldots, \lambda_r) = (h_1^0, \ldots, h_q^0, \lambda_1^0, \ldots, \lambda_r^0)$, since the cross-validated smoothing parameter case follows by stochastic equicontinuity arguments analogous to those used in the proof of Theorem 2.2.

For expositional simplicity, in the derivations below, we will assume that $n_1 = n_2 = n$.\footnote{The proof for the general $n_1 \neq n_2$ case is similar, but much more tedious notationally.} We write

$$J_n = J_{1n} + J_{2n},$$

where
where

\[ J_{1n} = - \frac{2}{n^2 \hat{p}_f \hat{p}_g} \sum_{w \in S_w} \sum_{i=1}^{n} \hat{K}_{\gamma, x \cdot y_i} I_{u_i, w} I_{v_i, w}, \]

and \( J_{2n} = J_n - J_{1n} \). Using \( n^{-2} = [n(n - 1)]^{-1} + O(n^{-3}) \) we can write \( J_{2n} \) as

\[
J_{2n} = \frac{1}{n(n - 1)} \sum_{i=1}^{n} \sum_{j \neq i} \sum_{w \in S_w} \left\{ \frac{\hat{K}_{\gamma, x \cdot x_j} I_{u_i, w} I_{v_j, w}}{\hat{p}_f} + \frac{\hat{K}_{\gamma, y \cdot y_i} I_{v_i, w} I_{v_j, w}}{\hat{p}_g} \right\} + O_p(n^{-1}).
\]

Using the fact that \( \hat{p}^{-1} = O_p(1) \) (\( \hat{p} = \hat{p}_f(w) \) or \( \hat{p} = \hat{p}_g(w) \)) and by the same argument used in the proof of \( I_{1n} = O_p(n^{-1}) \), one can easily show that \( J_{1n} = O_p(n^{-1}) \).

Next, we consider \( J_{2n} \). Using the fact that \( \hat{p} - p = O_p(n^{-1/2}) \) (\( p = p_f(w) \) or \( p = p_g(w) \)), the following expansion immediately follows:

\[
\frac{1}{\hat{p}} = \frac{1}{p} + \frac{p - \hat{p}}{p^2} + O_p(n^{-1}),
\]

\[
\frac{1}{\hat{p}^2} = \frac{1}{p^2} + \frac{2(p - \hat{p})}{p^3} + O_p(n^{-1}).
\]

Define the leave-two-out estimators \( \hat{p}_{f, -\{ij\}} = (n-2)^{-1} \sum_{i \neq j} I_{u_i, w}, \) and \( \hat{p}_{g, -\{ij\}} = (n-2)^{-1} \sum_{i \neq j} I_{v_i, w}. \) From \( \hat{p}_f = n^{-1} \sum_{i=1}^{n} I_{u_i, w} \) and \( \hat{p}_g = n^{-1} \sum_{i=1}^{n} I_{v_i, w}, \) it is easy to see that \( \hat{p}_f = \hat{p}_{f, -\{ij\}} + O_p(n^{-1}) \) and \( \hat{p}_g = \hat{p}_{g, -\{ij\}} + O_p(n^{-1}) \) (uniformly in \( i, j = 1, \ldots, n \)). Hence, we can replace \( \hat{p}_f \) and \( \hat{p}_g \) in (A.18) by \( \hat{p}_{f, -\{ij\}} \) and \( \hat{p}_{g, -\{ij\}} \) without affecting the asymptotic behavior of \( J_{2n} \).

Hence, substituting (A.19) into (A.18), and using the leave-two-out estimators \( \hat{p}_{f, -\{ij\}} \) and \( \hat{p}_{g, -\{ij\}} \) to replace \( \hat{p}_f \) and \( \hat{p}_g \), we obtain

\[
J_{2n} = J_{2n}^a + J_{2n}^b + O_p(n^{-1}),
\]
where $J_{2n}^a$ is obtained from $J_{2n}$ by replacing $\hat{p}_f$ and $\hat{p}_g$ by $p_f$ and $p_g$ in $J_{2n}$, i.e.,

$$J_{2n}^a = \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j \neq i} \sum_{w \in S_w} \left[ \frac{K_{i,x_i,x_j}I_{u_i,w}I_{u_j,w}}{p_f^2} + \frac{K_{i,y_i,y_j}I_{v_i,w}I_{v_j,w}}{p_g^2} - \frac{1}{p_fp_g} \left[ K_{i,x_i,y_j}I_{u_i,w}I_{v_j,w} + K_{i,x_j,y_i}I_{u_j,w}I_{v_i,w} \right] \right]$$

$$= \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j \neq i} H_{n,ij}^a,$$

where $H_{n,ij}^a = \sum_{w \in S_w} \left\{ K_{i,x_i,x_j}I_{u_i,w}I_{u_j,w}/p_f^2 + K_{i,y_i,y_j}I_{v_i,w}I_{v_j,w}/p_g^2 - [K_{i,x_i,y_j}I_{u_i,w}I_{v_j,w} + K_{i,x_j,y_i}I_{u_j,w}I_{v_i,w}]/(p_fp_g) \right\}$, and

$$J_{2n}^b = \frac{1}{n(n-1)} \sum_{w \in S_w} \left\{ \sum_{i=1}^{n} \sum_{j \neq i} \sum_{l \neq i,j} \frac{2(p_f - \hat{p}_f, -(ij))}{p_f^2} K_{i,x_i,x_j}I_{u_i,w}I_{u_j,w} + \frac{2(p_g - \hat{p}_g, -(ij))}{p_g^2} K_{i,y_i,y_j}I_{v_i,w}I_{v_j,w} \right\}$$

$$- \sum_{i=1}^{n} \sum_{j \neq i} \sum_{l \neq i,j} \left[ \frac{p_g - \hat{p}_g, -(ij)}{p_g^2} \right] \left[ \frac{(p_f - \hat{p}_f, -(ij))}{p_f^2} \right] \left[ K_{i,x_i,y_j}I_{u_i,w}I_{v_j,w} + K_{i,x_j,y_i}I_{u_j,w}I_{v_i,w} \right]$$

$$= \frac{1}{n(n-1)(n-2)} \sum_{i=1}^{n} \sum_{j \neq i} \sum_{l \neq i,j} \sum_{w \in S_w} \left\{ \frac{2(p_f - I_{u_i,w})}{p_f^2} K_{i,x_i,x_j}I_{u_i,w}I_{u_j,w} + \frac{2(p_g - I_{v_i,w})}{p_g^2} K_{i,y_i,y_j}I_{v_i,w}I_{v_j,w} \right\}$$

$$- \sum_{i=1}^{n} \sum_{j \neq i} \sum_{l \neq i,j} \sum_{w \in S_w} \left[ \frac{p_g - I_{v_i,w}}{p_g^2} \right] \left[ \frac{(p_f - I_{u_i,w})}{p_f^2} \right] \left[ K_{i,x_i,y_j}I_{u_i,w}I_{v_j,w} + K_{i,x_j,y_i}I_{u_j,w}I_{v_i,w} \right].$$

We will first analyze $J_{2n}^a$. $J_{2n}^a$ is a second order U-Statistic. Below we show that it is a degenerate U-statistic. Letting $Z_i = (X_i, V_i, Y_i, W_i)$, we have

$$E \left( H_{n,ij}^a \mid Z_i \right) = \sum_{w \in S_w} \left\{ \frac{I_{u_i,w}}{p_f} E \left( \frac{K_{i,x_i,x_j}I_{u_i,w}}{p_f^2} \mid Z_i \right) + \frac{I_{v_i,w}}{p_g} E \left( \frac{K_{i,y_i,y_j}I_{v_i,w}}{p_g^2} \mid Z_i \right) \right\}$$

$$- \frac{1}{p_fp_g} \left[ I_{v_i,w} E \left( K_{i,x_i,y_j}I_{u_i,w} \mid Z_i \right) + I_{v_i,w} E \left( K_{i,x_j,y_i}I_{u_j,w} \mid Z_i \right) \right]$$

$$= \sum_{w \in S_w} \left\{ \frac{I_{u_i,w}}{p_f} \left[ p_f^{-1} E \left( K_{i,x_i,x_j}I_{u_i,w} \mid Z_i \right) - p_g^{-1} E \left( K_{i,x_j,y_i}I_{u_j,w} \mid Z_i \right) \right] \right\}$$

$$+ \frac{I_{v_i,w}}{p_g} \left[ p_g^{-1} E \left( K_{i,y_i,y_j}I_{v_i,w} \mid Z_i \right) - p_f^{-1} E \left( K_{i,x_j,y_i}I_{u_j,w} \mid Z_i \right) \right]$$

\( (A.21) \)

$$= 0.$$
because
\[
p_f^{-1} E(K_{x_i, x_j I_u, I_w} | Z_i) - p_g^{-1} E(K_{x_i, x_j I_v, I_w} | Z_i) = \int [f(x, w)/p_f(w)] K_\gamma((x - x_i)/h) dx \\
- \int [g(x, w)/p_g(w)] K_\gamma((x - x_i)/h) dx \\
= 0
\]
since \( f(x, w)/p_f(w) = g(x, w)/p_g(w) \) under \( H^c_0 \).

By utilizing Hall’s (1984) central limit theorem for degenerate U-statistics, one can show that
\[
(A.22) \quad (n^2 h_1 \ldots h_q)^{-1/2} J_{2n}^{a} / \sigma_{n,J} \rightarrow N(0, 1) \text{ in distribution,}
\]
where
\[
\sigma_{n,J}^2 = \frac{2(h_1 \ldots h_q)}{n^2} \sum_i \sum_{j \neq i} \sum_{w \in S_w} E \left\{ \left( K_{x_i, x_j I_u, I_w} \right)^2 p_f(w)^{-4} + \left( K_{y_i, y_j I_v, I_w} \right)^2 p_g(w)^{-4} \right\} \\
+ \left[ \left( K_{x_i, y_j I_u, I_w} \right)^2 + \left( K_{y_i, x_j I_v, I_w} \right)^2 \right] (p_f(w)p_g(w))^{-2}.
\]

It is straightforward to show that \( \hat{\sigma}_{n,J}^2 = \sigma_{n,J}^2 + o_p(1) \). Hence, we have
\[
(A.23) \quad (n^2 h_1 \ldots h_q)^{-1/2} J_{2n}^{a} / \hat{\sigma}_{n,J} \rightarrow N(0, 1).
\]

Next, we consider \( J_{2n}^{b} \). We will show that \( J_{2n}^{b} \) has an order smaller than that of \( J_{2n}^{a} \). Because we only need to evaluate the order of \( J_{2n}^{b} \), we will omit \( \sum_{w \in S_w} \) for simplify the notation. Alternatively, one can think that, for each \( w \in S_w \), we derive an upper bound for \( J_{2n}^{b}(w) \) for any fixed value of \( w \). Since \( S_w \) is a finite set, the same bound holds true for \( \max_{w \in S_w} |J_{2n}^{b}(w)| \).

Note that \( J_{2n}^{b} \) contains three summations, therefore, it can be written as a third order U-statistic, i.e.,
\[
J_{2n}^{b} = \frac{1}{3n(n-1)(n-2)} \sum_{i=1}^{n} \sum_{j \neq i} \sum_{l \neq i,j} [J_{n,ijl}^b + J_{n,jil}^b + J_{n,lji}^b] \\
= \frac{1}{3n(n-1)(n-2)} \sum_{i=1}^{n} \sum_{j \neq i} \sum_{l \neq i,j,l} H_{n,ijl}^b,
\]
where \( H_{n,ijl}^b = J_{n,ijl}^b + J_{n,jil}^b + J_{n,lji}^b \), and

\[
J_{n,ijl}^b = \left[ \frac{2}{p_f} (p_f - I_{u_i,w}) W_{h,x_i,x_j} I_{u_i,w} I_{u_j,w} + \frac{2}{p_g} (p_g - I_{v_i,w}) W_{h,y_i,y_j} I_{v_i,w} I_{v_j,w} \right. \\
- \left. \frac{(p_g - I_{v_i,w})}{p_f p_g} + \frac{(p_f - I_{v_i,z})}{p_f p_g} \right] [W_{h,x_i,y_j} I_{u_i,w} I_{v_j,w} + W_{h,x_j,y_i} I_{u_j,w} I_{v_i,w}] .
\]

Let \( Z_i = (X_i, Y_i, U_i, V_i) \). Below we show that \( E(H_{n,ijl}^b | Z_i) = 0 \) under \( H_0^b \). Note that

\[
E(H_{n,ijl}^b | Z_i) = E(J_{n,ijl}^b | Z_i) + E(J_{n,jil}^b | Z_i) + E(J_{n,lji}^b | Z_i).
\]

From \( E(p_f - I_{v_i,z}) = 0 \) we immediately have \( E(J_{n,jil}^b | Z_i) = 0 \) and \( E(J_{n,lji}^b | Z_i) = 0 \). Now,

\[
E(J_{n,ijl}^b | Z_i) = (p_f - I_{u_i,w}) p_f^{-2} \left[ 2p_f^{-1} E(K_{\gamma,x_i,x_j} I_{u_i,w} I_{u_j,w}) \right. \\
- \left. p_g^{-1} E(K_{\gamma,y_i,y_j} I_{u_i,w} I_{v_j,w}) - p_g^{-1} E(K_{\gamma,x_j,y_i} I_{u_j,w} I_{v_i,w}) \right] \\
+ (p_g - I_{v_i,w}) p_g^{-2} \left[ 2p_g^{-1} E(K_{\gamma,y_i,y_j} I_{v_i,w} I_{v_j,w}) \right. \\
- \left. p_f^{-1} E(K_{\gamma,x_j,y_i} I_{v_j,w} I_{v_i,w}) - p_f^{-1} E(K_{\gamma,x_i,x_j} I_{v_i,w} I_{v_j,w}) \right] \\
= 0,
\]

by the same arguments as we used in the proof of \( E(H_{n,ij}^a | Z_i) = 0 \) (since \( f(x,w)/p_f(w) = g(x,w)/p_g(w) \)).

Hence, \( E(H_{n,ijl}^a | Z_i) = 0 \), and \( J_{n}^a \) is a degenerate U-statistic. Define \( H_{n,ij}^b = E(H_{n,ijl}^b | Z_i, Z_j) \).

Then by the standard change-of-variable argument, one can show that

\[
(A.24) \quad \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j \neq i}^{n} H_{n,ij}^b = \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j \neq i}^{n} H_{n,ij,0}^b + O_P(|h|^2 n^{-1}),
\]

where \( |h|^2 = \sum_{s=1}^{q} h_s^2 \), and

\[
H_{n,ij,0}^b = E(H_{n,ijl}^b | Z_i, Z_j) = 2p_f^{-3}(p_f - I_{u_i,w}) f(x_j,w) I_{u_j,w} + 2p_g^{-3}(p_g - I_{v_i,w}) g(y_j,w) I_{v_j,w} \\
- \left[ \frac{(p_g - I_{v_i,w})}{p_f p_g} + \frac{(p_f - I_{v_i,z})}{p_f p_g} \right] [f(y_j,w) I_{v_j,w} + g(x_j,w) I_{u_j,w}] .
\]
Therefore, by the U-statistic H-decomposition (see Lee (1990)), we have

\begin{equation}
J_{2n}^b = J_{2n,0}^b + O_p(n^{-1}),
\end{equation}

where

\begin{align*}
J_{2n,0}^b = \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j \neq i} H_{n,ij,0}^b.
\end{align*}

Note that \( E(H_{n,ij,0}^b|Z_i) = 0 \), i.e., \( J_{2n,0}^b \) is a degenerate U-statistic. Also note that \( H_{n,ij,0}^b \) is unrelated to the smoothing parameters \( (h_1, \ldots, h_q) \). Then it is easy to show that \( E[(J_{2n,0}^b)^2] = O(n^{-2}) \), which implies that

\begin{equation}
J_{2n,0}^b = O_p(n^{-1}).
\end{equation}

(A.24) to (A.26) lead to

\begin{equation}
J_{2n}^b = O_p \left( n^{-1} \right) = o_p \left( (n^2 h_1 \ldots h_q)^{-1/2} \right).
\end{equation}

Combining (A.16), (A.20), (A.23), (A.27), and the fact that \( J_{1n} = O_p(n^{-1}) \), we have completed the proof of Theorem 3.1. \( \square \)

REFERENCES


