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INFERENCE AND SPECIFICATION TESTING IN THRESHOLD REGRESSION WITH ENDOGENEITY

By

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Inference and Specification Testing in Threshold Regression with Endogeneity*

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Abstract

We propose three new methods of inference for the threshold point in endogenous threshold regression and two specification tests designed to assess the presence of endogeneity and threshold effects without necessarily relying on instrumentation of the covariates. The first inferential method is a parametric two-stage least squares method and is suitable when instruments are available. The second and third methods are based on smoothing the objective function of the integrated difference kernel estimator in different ways and these methods do not require instrumentation. All three methods are applicable irrespective of endogeneity of the threshold variable. The two specification tests are constructed using a score-type principle. The threshold effect test extends conventional parametric structural change tests to the nonparametric case. A wild bootstrap procedure is suggested to deliver finite sample critical values for both tests. Simulations show good finite sample performance of these procedures and the methods provide flexibility in testing and inference for practitioners working with threshold models.

Keywords: Threshold regression, Endogeneity, Identification, Confidence interval, 2SLS, IDKE, Specification testing, Bootstrap, U-statistic

JEL-Classification: C21, C24, C26.

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1 Introduction

In recognition of potential shifts in economic relationships, threshold models have become increasingly popular in recent econometric practice. Hansen (2011) provides an overview of the methods and their various applications in economics and finance. One typical application of the threshold model in time series is to illustrate asymmetric effects of shocks over the business cycle (e.g., Potter, 1995). Threshold models are also useful in cross section applications. For example, Hansen (2000) applied a threshold model to show that depending on the starting point, rich countries and poor countries have different growth patterns. All this literature assumes exogenous regressors and an exogenous threshold variable. But in practical work there is often uncertainty about exogeneity and threshold model applications have commonly encountered endogeneity issues. For example, the empirical growth models used in Papageorgiou (2002) and Tan (2010) both suffer from endogenous regressor problems, as argued in Frankel and Romer (1999) and Acemoglu et al. (2001).

The standard model formulation for endogenous threshold regression is

$$y = x' \beta_1 1(q \leq \gamma) + x' \beta_2 1(q > \gamma) + \varepsilon =: x' \beta + x' \delta 1(q \leq \gamma) + \varepsilon,$$

with $E[\varepsilon|x] \neq 0$, where $x = (1, x', q)' \in \mathbb{R}^{d+1} =: \mathbb{F}$, and where $d$ and $\tilde{d}$ are the dimensions of the nonconstant covariates $(x', q)'$ and all covariates including the constant. The parameter of interest is $\theta = (\beta_1', \beta_2', \gamma)'$ or equivalently, $\theta := (\beta', \delta', \gamma)'$ with $\beta = \beta_2$, $\delta = \beta_1 - \beta_2$ and $\gamma \in \Gamma$. This setup is similar to endogenous linear regression except that the regression coefficients depend on whether the threshold variable $q$ crosses the threshold point $\gamma$.

The literature on estimation of this model includes the following three main contributions. First, Caner and Hansen (2004) (CH hereafter) use a two-stage least squares (2SLS) method to estimate $\gamma$ in the small-threshold-effect framework of Hansen (2000), but assuming $q$ is exogenous so that $E[\varepsilon|x] = E[\varepsilon|x]$ holds. Second, working in Hansen’s (2000) framework, Kourtellos, Stengos and Tan (2016) (KST hereafter) use a control function approach to deal with the case where $q$ is also endogenous\(^1\). Their setup is parametric (see Kourtellos et al. (2017) for a semiparametric extension) and the asymptotic theory is flawed. Specifically, Yu, Liao and Phillips (2018) (YLP hereafter) show that the structural threshold regression (STR) estimator of the threshold point in KST is not consistent unless the endogeneity level of the threshold variable is low compared to the threshold effect. Third, Yu and Phillips (2018) (YP hereafter) use an integrated difference kernel estimator (IDKE) to estimate $\gamma$ in the fixed-threshold-effect framework of Chan (1993). Their estimator can be applied irrespective of whether $q$ is endogenous or whether instruments are available (as required in the previous two methods). Even when there are no instruments available and the model reduces to a nonparametric threshold regression, their estimator is still $n$-consistent, just as in the parametric setup\(^2\). The endogeneity problem is also considered in the related structural change literature, where the threshold variable is a simple time index and is always exogenous; see YP for a detailed literature review.

In spite of the theoretical developments on the estimation of $\gamma$ in endogenous threshold regression, infer-

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\(^1\)If $q$ is exogenous, then KST’s estimator is asymptotically equivalent to CH’s estimator.

\(^2\)There are two other estimators of $\gamma$ in nonparametric threshold regression with different motivations and objective functions from the IDKE. The first estimator is the semiparametric M-estimator of Henderson et al. (2017). That estimator can be treated as an extension of the partial linear estimator of Porter (2003) (see also Yu (2016)) in regression discontinuity designs to the case with unknown discontinuity point and extra covariates (beyond $q$), but this estimator can be applied only to the case with constant threshold effects; readers are referred to the supplementary materials of YP to see why the authors avoid using a generalized version of this estimator. The second estimator is the least squares estimator of Chiu et al. (2018). Chiu et al. (2018) can be treated as a nonparametric parallel of Bai and Perron (1998); for example, they allow for multiple regimes, use sequential tests to determine the number of regimes, and $q \notin x$ (because $q$ in structural change models is the time index which is usually not a covariate in $x$); again, readers are referred to the supplementary materials of YP to see why the authors avoid using this estimator.
ence concerning the threshold parameter $\gamma$ still presents practical difficulties especially when $q$ is endogenous. First, the CH method should be applied only if $q$ is exogenous. For as shown in Yu (2013a), the CH estimator is generally inconsistent when $q$ is endogenous. Second, as mentioned above, the KST approach is not generically applicable. Third, the asymptotic distribution of the IDKE in YP is too complicated to be readily applied in empirical work. This paper seeks to alleviate these practical difficulties by proposing three new methods of confidence interval (CI) construction for $\gamma$. All three methods can be applied regardless of whether $q$ is endogenous. To our knowledge, these methods are the only valid and applicable inferential tools that are robust to endogeneity of $q$ in the sense that the procedures need no modification when $q$ is endogenous. The first method is a parametric two stage least squares (2SLS) method and requires instruments, while the second and third methods are based on smoothing the objective function of the IDKE in different ways so that instruments are unnecessary. In discussing the first method of inference, we also discuss the identification issue of $\gamma$ using moment conditions. Of the two remaining methods, the second method assumes fixed threshold effects and uses the data around the threshold point marginally, while the third method assumes shrinking threshold effects and makes full use of data around the threshold point. So the second and third methods are similar in spirit to the smoothed maximum score (SMS) estimator of Horowitz (1992). On the other hand, the two methods differ from the SMS estimator in the sense that they have slower convergence rates than the IDKE in YP while the SMS estimator improves the convergence rate over the original maximum score estimator of Manski (1975, 1985). This feature of the methods is in some sense similar to the smoothed least squares estimator (SLSE) of Seo and Linton (2007) which also has a slower convergence rate than the usual least squares estimator (LSE) in, e.g., Yu and Fan (2019). Like the original IDKE approach, both of these IDKE-smoothing methods are nonparametric and require kernel and bandwidth selection. Practitioners can select an approach to inference from among these three methods based on their suitability to the data and on data availability. For example, if instruments are available, then the first method can be used; otherwise, the second and third methods may be preferable.

This paper also proposes two specification tests. The first tests for the existence of endogeneity, and the second tests for the presence of threshold effects with and without instruments. Both tests are score-type tests in the sense that they are constructed under the null, and their asymptotic properties are therefore easy to develop. More importantly, these tests of structural shifts are easier to implement in practice than the popular Wald test especially when instruments are unavailable. Because the Wald and score tests of structural shifts when instruments are available are standard extensions of existing tests and are well understood in the literature, these tests are relegated to an online supplement and the main text of the paper concentrates on the test without instruments. Both specification tests discussed in the main text take a nonparametric form and have an asymptotic normal (null) distribution. We suggest a wild bootstrap procedure to obtain critical values for these tests. Practitioners are encouraged to give greater attention to the inference methods and specification tests that are developed without instrumentation because good instruments are often hard to find and justify in practical work.

The rest of this paper is organized as follows. Section 2 provides an overview of the three inference methods and the two specification tests. Sections 3 to 5 analyze the three inference methods in turn and derive the corresponding asymptotic theory for constructing CIs. Section 6 presents the limit theory of the two specification tests and shows how to bootstrap the critical values. Section 7 includes some simulation results and Section 8 concludes.\footnote{Proofs of theorems with supporting propositions and lemmas are given in Supplements A, B and C. Additional discussion on parametric tests for the presence of threshold effects.}

\footnote{We will not discuss inference on regular parameters such as $\beta$ and $\delta$ because these cases fall within the standard literature; see, e.g., CH, YP and YLP. The first inference method in this paper also covers $\beta$ and $\delta$.}

\footnote{The dissertation of Qin Liao includes a serious empirical application using the techniques in this paper. To restrain the length of the present paper, we decided to discuss this application in separate work.}
when instruments are available is given in Supplement D. These supplements are collected together for online access in Yu et al. (2019).

A word on notation. The three inference methods in the paper are labeled Methods I, II and III. The symbol \( \ell \) is used to indicate the two regimes in \( \ell = 1, 2 \) or the two specification tests, and is not always written explicitly as ‘\( \ell = 1, 2 \)’. For matrices, \( A > 0 \) means that \( A \) is positive definite, \( \text{span}(A) \) denotes the column space of \( A \), and \( I_n \) is the \( n \times n \) identity matrix. For a random sequence \( Z_n \), \( \text{plim}Z_n \) means the probability limit of \( Z_n \) as the sample size \( n \to \infty \). For any random vector \( x \), \( x_{\leq \gamma} := x1(q \leq \gamma) \) and \( x_{> \gamma} \) is similarly defined. For any two random vectors \( x \) and \( y \), \( x \perp y \) means \( x \) and \( y \) are independent. A parameter with a subscript 0 is the true value of the parameter.

2 Overview of Inferential Methods and Specification Testing

This section briefly overviews the three inferential methods and the two specification tests, introduces notation useful in the subsequent development, and details assumptions employed in the asymptotic theory.

2.1 Methods of Inference for the Threshold Point

If we write the model \( y = G(x, q; \theta) + \varepsilon \), with \( E[\varepsilon | x, q] \neq 0 \), where \( G(x, q; \theta) = x' \beta + x' \delta 1(q \leq \gamma) \) is a nonlinear function of \( x, q \), then estimation of \( \gamma \) can be treated as in a nonlinear regression model with endogeneity. As argued in Section 2.1.6 of Blundell and Powell (2003), the fitted-value method of 2SLS relies heavily on linearity of the regression function, a feature that can explain why the 2SLS estimators in Yu (2013a) are not consistent. To restore consistency of 2SLS, we need to maintain the linear structure of the model. In other words, instead of projecting \( (x, q) \) on instruments \( z \), we first project \( (x, x_{\leq \gamma}) \) for a fixed \( \gamma \) on \( z \) to get the predictors \( \tilde{x} \) and \( \tilde{x}_{\leq \gamma} \); then we can find \( \tilde{\beta}(\gamma) \) and \( \tilde{\delta}(\gamma) \) by regressing \( y \) on \( \tilde{x} \) and \( \tilde{x}_{\leq \gamma} \); finally, \( \hat{\gamma} \) is obtained by minimizing \( \sum_{i=1}^{n}(y_i - \tilde{x}_i' \tilde{\beta}(\gamma) - \tilde{x}_{\leq \gamma, i}' \tilde{\delta}(\gamma))^2 \), from which we obtain \( \tilde{\beta} = \tilde{\beta}(\hat{\gamma}) \) and \( \tilde{\delta} = \tilde{\delta}(\hat{\gamma}) \). It is easy to see that this procedure is equivalent to the instrumental variable (IV) extremum estimation problem

\[
\left( \hat{\beta}, \hat{\delta}, \hat{\gamma} \right) = \arg \min_{\beta, \delta, \gamma} \left( Y - X \beta - X_{\leq \gamma} \delta \right)' P_Z(Y - X \beta - X_{\leq \gamma} \delta),
\]

where \( Y, X, X_{\leq \gamma} \) and \( Z \) are matrices stacking \( y_i, x_i', x_{\leq \gamma, i}' \) and \( z \) respectively, and \( P_Z \) is the projection matrix onto the instrument space \( \text{span}(Z) \). This method, labeled as Method I, treats \( \gamma \) as a regular parameter and is just nonlinear 2SLS, as in Amemiya (1974). We will also show that this 2SLS estimator may be viewed as a version of the GMM estimator considered in Hall, Han and Boldea (2012) (HHB hereafter; see also Andrews, 1993) but one that turns out to have more desirable asymptotic properties, including consistency, in the endogenous threshold regression case.

To better understand the estimator \( \hat{\gamma} \), we consider the case where \( x = 1, \beta_0 = 0 \) and \( \delta_0 = 1 \) are known, and \( z = 1 \). For this simple case, \( y = 1 (q \leq \gamma_0) + \varepsilon \), and the moment condition is \( E[\varepsilon | x] = E[y] - F_q(\gamma_0) = 0 \), where \( F_q(\cdot) \) is the cdf of \( q \). In other words, \( \gamma_0 = F_q^{-1}(E[y]) \) is the \( E[y] \)-th quantile of \( q \), and \( \hat{\gamma} = F_q^{-1}(\hat{\gamma}) \). Following this intuition, we will show that: (i) \( \hat{\gamma} \) is \( \sqrt{n} \)-consistent, asymptotically normal, and the asymptotic variance involves the density of \( q \) at \( \gamma_0 \) (i.e., \( f_q(\gamma_0) \)) as in quantile regression; (ii) different from the usual threshold regression estimators where \( \hat{\gamma} \) is asymptotically independent of \( (\hat{\beta}, \hat{\delta}) \), this new estimator \( \hat{\gamma} \) is correlated with \( (\hat{\beta}, \hat{\delta}) \) asymptotically. Given these two results, a valid CI for \( \gamma \) can be constructed jointly with \( (\beta, \delta) \) by means of the bootstrap, just as in quantile regression to avoid nonparametric estimation of the density \( f_q(\gamma_0) \) in the asymptotic distribution.

Differing from CH estimation, this 2SLS estimator can be applied irrespective of whether \( q \) is endogenous.
Yu (2013a) shows that when q is exogenous, the CH estimator is inconsistent if the first stage predictor is a projection rather than a conditional mean. By contrast our 2SLS estimator requires only a linear projection in the first stage, so it is more robust in this respect.

Moving away from instrument-based estimation, we next introduce instrument-free estimators in Methods II and III by extending the IDKE of YP in different directions. Without instruments, the model reduces to a nonparametric threshold regression that can be written in the form

\[ y = m(x, q) + u = m_-(x, q)1(q \leq \gamma) + m_+(x, q)1(q > \gamma) + u, \]

where \( u = \varepsilon - \mathbb{E}[\varepsilon|x, q] \), \( m_-(x, q) = x'\beta_1 + \mathbb{E}[\varepsilon|x, q] \), \( m_+(x, q) = x'\beta_2 + \mathbb{E}[\varepsilon|x, q] \), \( g(x, q) = m_+(x, q) \) when \( q > \gamma \) and is the smooth extension of \( m_+(x, q) \) when \( q \leq \gamma \), and \( \Delta(x, q) = m_-(x, q) - m_+(x, q) \). This setup allows \( \mathbb{E}[\varepsilon|x, q] \) to be kinked or discontinuous at \( \gamma \). When \( \mathbb{E}[\varepsilon|x, q] \) is smooth, then \( g(x, q) = x'\beta + \mathbb{E}[\varepsilon|x, q] \) and \( \Delta(x, q) = x''\delta \); otherwise, \( g(x, q) \neq x'\beta + \mathbb{E}[\varepsilon|x, q] \) for \( q \leq \gamma \) and \( \Delta(x, q) \neq x''\delta \). When \( \mathbb{E}[\varepsilon|x, q] \) is continuous \( \Delta(x, \gamma) = (1, x', \gamma)\delta \).

To construct the IDKE of \( \gamma \), we start by defining a generalized kernel function, following Müller (1991).

**Definition:** \( k_h(\cdot, \cdot) \) is called a univariate generalized kernel function of order \( p \) if \( k_h(u, t) = 0 \) when \( u > t \) or \( u < t - 1 \) and for all \( t \in [0, 1] \),

\[
\int_{t-1}^{t} u^j k_h(u, t) du = \begin{cases} 
1, & \text{if } j = 0, \\
0, & \text{if } 1 \leq j \leq p - 1.
\end{cases}
\]

A popular example of the generalized kernel function is obtained as follows. Define the space

\[ M_p([a, b]) = \left\{ g \in \text{Lip}([a, b]), \int_a^b x^j g(x) dx = \begin{cases} 
1, & \text{if } j = 0, \\
0, & \text{if } 1 \leq j \leq p - 1
\end{cases} \right\}, \]

where \( \text{Lip}([a, b]) \) denotes the space of Lipschitz continuous functions on \([a, b]\). Define \( k_+(\cdot, \cdot) \) and \( k_-(\cdot, \cdot) \) as follows:

(i) The support of \( k_-(x, r) \) is \([-1, r] \times [0, 1]\) and the support of \( k_+(x, r) \) is \([-r, 1] \times [0, 1]\).

(ii) \( k_-(\cdot, r) \in M_p([-1, 1]) \) and \( k_+(\cdot, r) \in M_p([-r, 1]) \).

(iii) \( k_+(x, r) = k_-(x, r) \).

(iv) \( k_-(1, r) = k_+(1, r) = 0 \).

Condition (iv) implies that \( k_-(\cdot, r) \) is Lipschitz on \(( -\infty, r] \) and \( k_+(\cdot, r) \) is Lipschitz on \([-r, \infty) \). This assumption is important in deriving the asymptotic distribution of the IDKE of \( \gamma \). Readers are referred to Appendix A of Porter and Yu (2015) for related discussion in the DKE case.

To simplify the construction of \( k_h(u, t) \), the following constraints are imposed on the support of \( x \) and on the parameter space.

**Assumption S:** \( (y, x', q') \in \mathbb{R} \times \mathcal{X} \times \mathcal{Q} \subset \mathbb{R}^7, \mathcal{X} = [0, 1]^{d-1}, \mathcal{Q} = [q, \bar{q}], \) and \( \gamma \in \Gamma = [\underline{\gamma}, \bar{\gamma}] \subset \mathcal{Q}, \beta \in B \subset \mathbb{R}^7, \delta \in \Xi \subset \mathbb{R}^2, \) where \( q \) can be \(-\infty \) and \( \bar{q} \) can be \( \infty \), and \( \Gamma, B \) and \( \Xi \) are compact.

We do not restrict \( \delta_0 \) to be fixed or to shrink to zero in all cases. Rather, \( \delta_0 \) is taken as fixed in Method II and shrinks to zero in Method III. We assume \( x \) is continuously distributed, but note that continuous and discrete
components may be accommodated, at least in a conceptually straightforward manner but at the expense of additional notational complexity, by using the continuous covariate estimator within samples homogeneous in the discrete covariates. Requiring the support of \( x \) to be \([0,1]^{d-1}\) is not restrictive as this support can be achieved by use of a suitable monotone transformation such as the empirical percentile transformation. The compactness assumption on \( \mathcal{X} \) simplifies the proof and may be relaxed by imposing restrictions on the tail of the distribution of \( x \).

Define

\[
\begin{align*}
k(+, 1) &= k(-, 1) \in \mathcal{M}_p([-1,1]), \\
k_0(+, 0) &= \mathcal{M}_p([0,1]), \\
k_0(-, 0) &= \mathcal{M}_p([-1,0]),
\end{align*}
\]

and

\[
k_h(u,t) = \begin{cases} 
\frac{1}{h} k \left( \frac{u}{h} \right), & \text{if } h \leq t < 1-h, \\
\frac{1}{h} k_0 \left( \frac{u}{h} \right), & \text{if } 0 \leq t < h, \\
\frac{1}{h} k_0 \left( \frac{u}{h} \right), & \text{if } 1-h \leq t \leq 1.
\end{cases}
\]  

Then \( k_h(u,t) \) is a generalized kernel function of order \( p \). We may construct a corresponding multivariate generalized kernel function of order \( p \) by taking the product of univariate generalized kernel functions of order \( p \). We only require \( k_h(u,t) \) to be a first order kernel function in Method II but may require it to be a higher order kernel function in Method III\(^5\). In particular, we use the following two conditions. For Method II we use:

**Assumption K**: \( k_h(u,t) \) takes the form of (3) with \( p = 1 \), \( k_0(0) = k_0(0) = 0 \), and \( k_0'(0) > 0 \), \( k_0'(0) < 0 \).

For Method III we use:

**Assumption K':** \( k_h(u,t) \) takes the form of (3) with \( p = s \), and \( k_0(0) = k_0(0) > 0 \), where \( s \) is the smoothness index of \( g(x,q) \) and will be defined in Assumption G below.

Assumption K mimics Assumptions B2 and B3 of Delgado and Hidalgo (2000) (DH hereafter) and Assumption K' is Assumption K in YP with the additional requirement that \( p = s \). Higher order kernels are required in Assumption K' only to achieve the optimal convergence rate of \( \gamma \). In practice, \( p = 1 \) is sufficient.

Given \( k_h(u,t) \), the IDKE of \( \gamma \) is constructed as the extremum estimator which satisfies

\[
\hat{\gamma} = \arg \max_{\gamma} \frac{1}{n} \sum_{i=1}^{n} \left[ \frac{1}{n-1} \sum_{j=1, j \neq i}^{n} y_j K_{h,i}^{\gamma_-} - \frac{1}{n-1} \sum_{j=1, j \neq i}^{n} y_j K_{h,i}^{\gamma_+} \right]^2 \quad (4)
\]

\[
= \arg \max_{\gamma} \frac{1}{n} \sum_{i=1}^{n} \Delta_i^2(\gamma) =: \arg \max_{\gamma} \hat{Q}_n(\gamma),
\]

where

\[
\begin{align*}
K_{h,i}^{\gamma_-} &= \prod_{l=1}^{d-1} k_h(x_{ij} - x_{il}, x_{il}) \cdot k_h^+ (q_j - \gamma) =: K_{h,i}^{\gamma_-} \cdot K_h^+ (q_j - \gamma), \\
K_{h,i}^{\gamma_+} &= \prod_{l=1}^{d-1} k_h(x_{ij} - x_{il}, x_{il}) \cdot k_h^+ (q_j - \gamma) =: K_{h,i}^{\gamma_+} \cdot K_h^+ (q_j - \gamma), \\
\Delta_i(\gamma) &= \frac{1}{n-1} \sum_{j=1, j \neq i}^{n} y_j K_{h,i}^{\gamma_-} - \frac{1}{n-1} \sum_{j=1, j \neq i}^{n} y_j K_{h,i}^{\gamma_+}.
\end{align*}
\]

\(^5\)Note here that the usual symmetric kernel is a second order kernel, but the boundary kernel is only a first order kernel because \( \int u k_h(u,t) = 0 \).
For notational convenience, we here use the same bandwidth for each dimension of \((x', q)\), although there may be some finite sample improvement from using different bandwidths in each dimension. As suggested in Yu (2012, 2015b), we need only check the mid-points of the contiguous \(q\)'s in the optimization process of \(f_0\). In other words, the argmax operator is a mid-point operator. The summation in the parenthesis of \(f_0\) excludes \(j = i\), which is a standard strategy in converting a V-statistic to a U-statistic. Also, the normalization factor \(\sum_{j=1,j\neq i}^n K_{h,ij}^{\pm}\) does not appear in the construction of \(\hat{\gamma}\), thereby avoiding random denominator issues in conditional mean estimation and simplifying the derivation of the limit distribution of \(\hat{\gamma}\), a technique that dates back at least to Powell et al. (1989). This form of \(\hat{\gamma}\) has some practical advantages especially when \(d\) is large. Since the conditional mean is estimated at the boundary point \(q = \gamma\), the local linear smoother (LLS) or local polynomial estimator (LPE) may be considered to ameliorate bias. However, when \(d\) is large, there are not many data points in a \(h\) neighborhood of \((x', \gamma)'\). As a result, not only does the LLS lose degrees of freedom (by estimating more parameters) but its denominator matrix can be close to singular, which disrupts finite sample performance. Further, differing from regular parameter estimation (such as conditional mean estimation), use of the LLS in this context does not affect the first-order asymptotic distribution of \(\hat{\gamma}\).

CI construction based on the limit distribution of \(\hat{\gamma}\) under Assumption K’ and fixed threshold effects (i.e., in the framework of YP) is challenging because the asymptotics involve a compound Poisson process, making simulation awkward. Methods II and III use different smoothing schemes to achieve more convenient asymptotic distributions. Method II assumes fixed threshold effects but uses data in the neighborhood of \(\gamma_0\) only marginally. The resulting asymptotic theory is normal and the CI can be constructed by inverting either the \(t\) or likelihood ratio (LR) statistic. Method III fully utilizes data in the neighborhood of \(\gamma_0\) but assumes shrinking threshold effects. The limit distribution then involves a two-sided Brownian motion. As suggested in Hansen (2000) we can invert the LR statistic (rather than the \(t\)-statistic) to improve finite-sample performance. As expected, due to insufficient usage of data information in the neighborhood of \(\gamma_0\), the convergence rates of the IDKE in both these methods are slower than the \(O(n)\) rate in YP. In Method II we also require \(k'_\pm(0) \neq 0\), or else the convergence rate of the IDKE is even slower.

We next provide some intuition that helps to justify the extremum estimator \(\hat{\gamma}\). For this purpose we impose the following Assumption F on the distribution of \((x', q)'\) in Method II and Assumptions G and G’ on \(g(x, q)\) in Method III.

Assumption F: The density \(f(x, q)\) of \((x', q)'\) is second order continuously differentiable and satisfies \(0 < f \leq f(x, q) \leq \overline{f} < \infty\) for \((x', q)' \in \mathcal{X} \times \Gamma_{\varepsilon}\), where \(\Gamma_{\varepsilon} := (\gamma - \varepsilon, \gamma + \varepsilon)\) for some \(\varepsilon > 0\) and \((f, \overline{f})\) are some fixed quantities.

Assumption G: \(g(x, q)\) is second order continuously differentiable on \(\mathcal{X} \times \Gamma_{\varepsilon}\).

Assumption G’: \(g(x, q)\) is \(s\)’th order continuously differentiable on \(\mathcal{X} \times \Gamma_{\varepsilon}\) with \(s \geq d\).

Assumption F implies that \(f_0(\gamma)\) is continuous, and \(0 < f_0 \leq f_0(\gamma) \leq \overline{f}_0 < \infty\) for \(\gamma \in \Gamma_{\varepsilon}\) and fixed \((\underline{f}_0, \overline{f}_0)\), and the conditional density \(f_{x|q}(x|q)\) is bounded below and above for \((x', q)' \in \mathcal{X} \times \Gamma_{\varepsilon}\); see Yu and Zhao (2013) for relaxation of these conditions. The first part of Assumption F implies that there are no discrete covariates in \(x\). As mentioned earlier in the remarks following Assumption S, this assumption is made for simplicity, just as in Robinson (1988), and is not critical to the methodology or the limit theory. The second

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6 Although in the fixed-threshold-effect framework with \(k_\pm(0) > 0\) the asymptotic distribution of \(\hat{\gamma}\) depends on whether the left endpoint or the middle point of the maximizing interval is taken as the maximizer, the asymptotic distributions of the two IDKEs in the present paper are both invariant to such choices of \(\hat{\gamma}\).
part of Assumption F implies that $\gamma_0$ is not on the boundary of $\mathcal{Q}$. Under these two assumptions, we can expect the objective function $\tilde{Q}_n(\gamma)$ to converge to

$$
\mathbb{E} \left[ (\mathbb{E}[y|x, q = \gamma^-] f(x, \gamma) - \mathbb{E}[y|x, q = \gamma^+] f(x, \gamma))^2 \right] = \int (\mathbb{E}[y|x, q = \gamma^-] - \mathbb{E}[y|x, q = \gamma^+])^2 f(x, \gamma)^2 f(x) dx.
$$

Since $f(x)$ and $f(x, \gamma)$ are continuous in $x$ and $\gamma$, there will be a jump in the limit only if $\gamma = \gamma_0$ which provides identifying information. In view of these properties, the threshold point can then be identified and consistently estimated by maximizing $\tilde{Q}_n(\gamma)$ under an additional requirement on the differential

$$
\Delta(x, \gamma_0) := \mathbb{E}[y|x, q = \gamma^-] - \mathbb{E}[y|x, q = \gamma^+],
$$

which enables identification of $\gamma_0$.

**Assumption I:** $\Delta(x, \gamma_0) \neq 0$ for $x$ in some set of positive Lebesgue measure in $\mathcal{X}$.

In Method I, $\Delta(x, \gamma_0) = (1, x', \gamma_0) \delta_0$, so we can replace $\Delta(x, \gamma_0)$ by $(1, x', \gamma_0) \delta_0$ in Assumption I. For comparison, we state the following Assumption I'.

**Assumption I':** $\delta_0 \neq 0$, where $\neq$ here means that at least one element is unequal.

Note that Assumption I is stronger than Assumption I' when $\Delta(x, \gamma_0) = (1, x', \gamma_0) \delta_0$. For example,

$$
\delta_0 = \begin{cases} 
(1, 0, -\frac{1}{\gamma_0})', & \text{if } \gamma_0 \neq 0, \\
(0, 0, 1)', & \text{if } \gamma_0 = 0,
\end{cases}
$$

is nonzero but does not satisfy Assumption I. Assumption I implies that $P((1, x', \gamma_0) \delta_0 \neq 0) > 0$, which excludes the continuous threshold regression (CTR) of Chan and Tsay (1998) (see also Hansen (2017)).

For comparison, we also review the DKE in DH here. Define the DKE

$$
\hat{\gamma} = \arg \max_{\gamma} \left[ \frac{1}{n} \sum_{j=1}^{n} y_j K_{h,j}^\gamma - \frac{1}{n} \sum_{j=1}^{n} y_j K_{h,j}^{\gamma^+} \right]^2 \equiv \arg \max_{\gamma} \Delta_2^\gamma(\gamma) =: \arg \max_{\gamma} \tilde{Q}_n(\gamma),
$$

where

$$
K_{h,j}^\gamma = \prod_{l=1}^{d-1} k_h(x_{lj} - x_{ol}, x_{ol}) \cdot k_h^\gamma(q_j - \gamma), \quad K_{h,j}^{\gamma^+} = \prod_{l=1}^{d-1} k_h(x_{lj} - x_{ol}, x_{ol}) \cdot k_h^{\gamma^+}(q_j - \gamma),
$$

and $x_o$ is some fixed point in the interior of $\mathcal{X}$. As explained in YP, selection of $x_o$ is difficult from both theoretical and practical perspectives. As distinct from the DKE, the IDKE procedure integrates the jump information over all the $x_i$, thereby removing the problem of choosing $x_o$. Further, usage of all the data ensures that the IDKE has greater identifying capability than the DKE in both Methods II and III.

---

7Strictly speaking, DH normalize the first term of $\Delta_2(\gamma)$ by $\tilde{f}_o^- := \frac{1}{n} \sum_{j=1}^{n} K_{h,j}^\gamma$ and the second term by $\tilde{f}_o^{\gamma^+} := \frac{1}{n} \sum_{j=1}^{n} K_{h,j}^{\gamma^+}$. However, their estimator is asymptotically equivalent to $\arg \max_{\gamma} \Delta_2^\gamma(\gamma) / f(x_o, \gamma_0)^2$ and has the same asymptotic distribution as $\hat{\gamma}$.
2.2 Overview of Two Specification Tests

The first specification test addresses potential endogeneity and the corresponding hypotheses $H^{(1)}$ are formulated as follows

$$H^{(1)}_0 : \mathbb{E}[\varepsilon|x,q] = 0,$$

$$H^{(1)}_1 : \mathbb{E}[\varepsilon|x,q] \neq 0.$$

This exogeneity test can be conducted prior to model estimation. When instruments are available, the Hausman test in Kapetanios (2010) can be applied. In the present paper we therefore consider only the case without instruments and apply the techniques developed in Fan and Li (1996) and Zheng (1996) to test the null $H^{(1)}_0$. In the second test, the hypotheses $H^{(2)}$ are

$$H^{(2)}_0 : \beta_1 = \beta_2 \text{ or } \delta = 0,$$

$$H^{(2)}_1 : \beta_1 \neq \beta_2 \text{ or } \delta \neq 0.$$

If $H^{(1)}_0$ is not rejected, i.e., there is no evidence of endogeneity, then $H^{(2)}$ involve a conventional parametric structural change test, such as that considered in Davies (1977, 1987), Andrews (1993), Andrews and Ploberger (1994) and Hansen (1996), among others. If $H^{(1)}_0$ is rejected, the ensuing situation is more complex. When there are instruments, Wald-type test statistics such as the sup-statistic in Section 5 of Caner and Hansen (2004) or score-type statistics such as those in Yu (2013b) can be used. Since the asymptotic distributions of both these types of test statistics are not pivotal, the simulation method of Hansen (1996) and De Jong (1996) can be applied to obtain critical values. Details concerning these tests are given in Supplement D of the paper because techniques for these tests are nowadays standard.

When there are no instruments, the Wald-type statistic is hard to implement since its asymptotic distribution is hard to derive given that $\hat{\delta}$ can only be estimated at a nonparametric rate – see Section 3.3 of Porter and Yu (2015) for discussion. However, the score-type test of Porter and Yu (2015) can be extended to this case with some technical complications. Importantly, the hypotheses $H^{(2)}$ relate to whether $m(x,q)$ is continuous, so $H^{(2)}_0$ encompasses more data generating processes (DGPs) than the null hypothesis in the usual structural change literature where $m(x,q)$ has a simple parametric form. In other words, the usual parametric tests have power against alternatives in which $m(x,q)$ does not take the form $x^\top \beta + x^\top \delta \mathbb{1}(q \leq \gamma)$ (see, e.g., Section 5.4 of Andrews (1993)) but our test has only trivial asymptotic power in such continuous $m(x,q)$ cases. A simple example may clarify the point. Suppose $m(x,q) = \alpha + \beta q$, in contrast to the specifications employed for our tests, which are based on $\mathbb{1}(q \leq \gamma)$, or $y = \alpha + \delta \mathbb{1}(q \leq \gamma) + \varepsilon$. It is easy to see that the usual tests have power against $m(x,q)$, which is very smooth in this case. In summary, the usual tests have power against both misspecification and certain types of structural change, whereas our test has non-trivial power only against threshold structural change, which may be more relevant in practical work.

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8Gao et al. (2008) discuss an average form of such a test in the time series context. But their test is not easy to extend to the case with a nonparametric threshold boundary as in the present framework. See also Hidalgo (1995) for a nonparametric conditional moment test for structural stability in a fully nonparametric environment, which focuses on global stability rather than local stability as here.

9In this framework and assuming $m(x,q) = x^\top \beta(q)$, the structural change tests focus on whether $\beta(q) = \beta$. See, e.g., Chen and Hong (2012), Kristensen (2012) and references therein for related tests in the time series context using nonparametric techniques. Actually, we can test whether $\beta(q)$ is continuous by extending the tests below, e.g., we can construct residuals $\tilde{e}_i$ in $I^{(2)}_0$ by estimating $\beta(q)$ using estimation techniques from the varying coefficient model (VCM) literature - see Robinson (1989, 1991), Cleveland et al. (1992) and Hastie and Tibshirani (1993) for early developments, and Fan and Zhang (2008) for a summary of recent developments.

10In the same way, there are also cases where the parametric test does not have power when there is a nonparametric threshold effect; see Example 1 of Hidalgo (1995).
advantage does not come for free: the usual tests have power against \( n^{-1/2} \) local alternatives, while our test needs a larger (than \( n^{-1/2} \)) local alternative to generate non-trivial power. Understandably so, because our test is essentially nonparametric whereas the usual tests are parametric.

In the following discussion of the two specification tests, \( H_0 \) indicates both \( H_0^{(1)} \) and \( H_0^{(2)} \), and \( H_1 \) indicates both \( H_1^{(1)} \) and \( H_1^{(2)} \). To construct a feasible test statistic, we need sample analogues of \( m(\cdot) \), \( \bar{m}(\cdot) \), and \( \bar{m}(\cdot) \) under both the null and alternative in these tests.

For the second test, we use the statistic
\[
I_n^{(2)} = \frac{n h d / 2}{n (n - 1)} \sum_i \sum_{j=1}^{n / h} K_{b,ij} e_i \hat{e}_j,
\]
and, for the second, we use
\[
I_n^{(2)} = \frac{n h d / 2}{n (n - 1)} \sum_i \sum_{j=1}^{n / h} 1_i^T 1_j^T K_{b,ij} e_i \hat{e}_j.
\]

Note further that \( e = u \) under \( H_0 \), so \( e \) has the same meaning in both \( I_n^{(1)} \) and \( I_n^{(2)} \) under \( H_0 \). Observe that \( \mathbb{E}[e \mathbb{E}[e|x,q] f(x,q)] = \mathbb{E}[\mathbb{E}[e|x,q]^2 f(x,q)] \geq 0 \) in the first test and \( \mathbb{E}[e \mathbb{E}[e|x,q] f(x,q) 1_q^T] = \mathbb{E}\left[\mathbb{E}[e|x,q]^2 f(x,q) 1_q^T\right] \geq 0 \) in the second test where the equalities hold if and only if \( H_0 \) holds. So we can construct the statistic based on the moment \( \mathbb{E}[e \mathbb{E}[e|x,q] f(x,q)] \) in the first test and the moment \( \mathbb{E}[e \mathbb{E}[e|x,q] f(x,q) 1_q^T] \) in the second test. Here, \( f(x,q) \) is added in to avoid the random denominator problem in kernel estimation, and \( 1_q^T \) appears in the second test because threshold effects can occur only on \( q \in \Gamma \).

To construct a feasible test statistic, we need sample analogues of \( e \) and \( \mathbb{E}[e|x,q] f(x,q) \). For the first
In this section, we derive the asymptotic distribution of the 2SLS estimator of $\theta$ and discuss some identifiability results for $\gamma$ when estimation is based on moment conditions. First, note that the 2SLS estimator of $\gamma$ can be written in GMM form as

$$\hat{\gamma} = \arg \min_{\gamma} \hat{Q}_n (\gamma),$$

where

$$\hat{Q}_n (\gamma) = \min_{\beta, \delta} \hat{Q}_n (\beta, \delta) := \min_{\beta, \delta} \hat{g}_n (\beta, \delta)' \hat{W} \hat{g}_n (\beta, \delta)$$

with $\hat{\theta} = (\hat{\beta}', \hat{\delta})'$, $\hat{W} = (n^{-1}Z'Z)^{-1}$ and

$$\hat{g}_n (\beta, \delta) = \frac{1}{n} \sum_{i=1}^{n} g_i (\beta, \delta) = \frac{1}{n} \sum_{i=1}^{n} z_i (y_i - x_i' \hat{\beta} - x_i' \delta, \delta).$$

To develop asymptotic properties of $\hat{\theta}$ we make the following assumption. First, throughout our analysis we use the notation $\delta_n$ for the true value of $\delta$ when we allow $\delta$ to shrink to zero, as in Hansen (2000), and we use $\delta_0$ to denote the true value of $\delta$ when $\delta$ is fixed to signify this difference.

**Assumption IV:** $\mathbb{E} [z e] = 0$, $\dim (z) = l \geq 2d + 1$, $\delta_n / \| \delta_n \| \to c$,

$$G = \left( \mathbb{E} [z x'], \mathbb{E} \left[ z x'_{\leq \gamma_0} \right], \mathbb{E} [z x' | q = \gamma_0] c f_q (\gamma_0) \right)$$
is of full column rank, $G_\gamma = E \left[ z(x', x'_{\leq \gamma}) \right] =: (G_1, G_{2, \gamma})$ is of full column rank for any $\gamma \in \Gamma$, $W := E [zz'] > 0$, and $\Omega := E [zz' \varepsilon^2] > 0$. Also, there does not exist a vector $a \in \mathbb{R}^2$ such that $G_\gamma a = G_{2, \gamma_0} c$ for any $\gamma \neq \gamma_0$.

When $\delta$ is fixed, the parameter $c$ is just the normalized form of $\delta$. When $\|\delta_n\| \to 0$, only the components of $c$ that correspond to the lowest shrink rate of $\delta_n$ are nonzero. Full column rankness of $G$ excludes CTR models where $x'\delta_n|q = \gamma_0$ is always zero so that the third part of $G$ is a zero matrix. But this assumption is nonetheless weaker than full column rankness of $E [zx'|q = \gamma_0]$. This is because if $E [zx'|q = \gamma_0]$ has full column rank, then 1 and $q$ cannot be elements of $x$ simultaneously; otherwise, the first and the last columns of $E [zx'|q = \gamma_0]$ would be collinear.

All other conditions in Assumption IV are standard except the last condition. This condition is required for the identification of $\gamma_0$. Take the fixed-$\delta$ case as an example where $c$ can be taken as $\delta_0$ and no normalization on $\delta$ is required. Note that if $E [zx] = 0$, then $E [zy] = G_{\gamma_0} \theta_0$. If there exists an $a := (a_1', a_2')'$ such that $G_\gamma a = G_{2, \gamma_0} \delta_0$ for some $\gamma \neq \gamma_0$, then under this $\gamma$, we can still let $\theta = (\beta_0 + a_1, a_2)$ satisfy the moment conditions, in which case the model is not identified by the moment conditions. This condition requires that the $l$-dimensional vector $G_{2, \gamma_0} \delta_0 \notin \bigcup_{\gamma \neq \gamma_0} \text{span} \left( G_1, G_{2, \gamma} \right)$, where $\text{span} \left( G_1, G_{2, \gamma} \right)$ is a $2l < l$ dimensional space. There is an important case where this condition is violated. If $q$ is independent of $(x', x')'$ in the structural change model where $q$ is the time index, then $G_{2, \gamma_0} = G_1 F_q (\gamma_0)$ and $a = (F_q (\gamma_0) \delta_0, 0)'$ satisfy $G_\gamma a = G_{2, \gamma_0} \delta_0$. In the TR context, if $q$ is independent of the rest of the system, then $q$ should be included in $z$, and cannot be independent of $(x', x')'$. This condition also implies the usual assumption that $x$ cannot be independent of the endogenous variables $(x', q)'$.

If this were the case, then $G_{2, \gamma_0} = E [z] E \left[ x'_{\leq \gamma_0} \right]$ would span a one-dimensional space, which can obviously be spanned by $G_1 = E [z] E \left[ x' \right]$ and $G_{2, \gamma} = E [z] E \left[ x'_{\gamma_0} \right]$.

**Theorem 1** Under Assumptions F, I, IV and S, $\hat{\theta}$ and especially $\hat{\gamma}$ are consistent and have limit distribution given by

$$
\left( \begin{array}{cc}
\sqrt{n} I_2 & 0 \\
0 & \sqrt{n} \|\delta_n\|
\end{array} \right) \left( \begin{array}{c}
\hat{\theta} - \theta_0 \\
\hat{\gamma} - \gamma_0
\end{array} \right) \xrightarrow{d} N (0, V)
$$

where $V = (G'WG)^{-1}G'W\Omega WG(G'WG)^{-1}$.

Note that we need only Assumption I′ to show the consistency of $\hat{\gamma}$. But to derive the asymptotic distribution we need Assumption I. Otherwise, $G$ need not be of full column rank. Also, as predicted in Section 2.1, $f_q (\gamma_0)$ appears in $V$ and $\hat{\gamma}$ is not asymptotically independent of $(\beta, \delta)$. $W$ can be any positive definite matrix besides $(n^{-1} Z'Z)^{-1}$. We still use $\tilde{W}$ to denote such a general weight matrix and use $W$ to denote its limit.

To provide some intuition on the asymptotic variance of $\hat{\gamma}$, consider the simple example in Section 2.1 again. In this example, $G = f_q (\gamma_0)$, $\Omega = \text{Var} (\varepsilon)$ and $W$ is irrelevant, so $V = \text{Var} (\varepsilon) / f_q (\gamma_0)^2$. In fact, $\hat{\gamma} = \hat{F}_q^{-1} (\hat{\gamma})$, so $\sqrt{n} (\hat{\gamma} - \gamma_0) = \sqrt{n} \left[ F_q^{-1} (\gamma) - F_q^{-1} (\hat{\gamma}) \right] + \sqrt{n} \left( F_q^{-1} (\mathbb{E} [ y ]) - F_q^{-1} (\mathbb{E} [ y ]) \right) + \sqrt{n} \left( F_q^{-1} (\mathbb{E} [ y ]) - F_q^{-1} (\mathbb{E} [ y ]) \right)$. The first term is, roughly speaking, $\sqrt{n} (\sum_{i=1}^n \psi_i (\hat{\gamma}) - \psi_i (\mathbb{E} [ y ]))$ with $\psi_i (\tau) = \psi_i (\gamma - \gamma_0)$ does not imply $E [zx'_{\leq \gamma_0}] \delta_n = 0$ or $E [zx'_{\leq \gamma_0}] \delta_n = 0$.

\footnote{Note that zero $x'\delta_n|q = \gamma_0$ does not imply $E [zx'_{\leq \gamma_0}] \delta_n = 0$ or $E [zx'_{\leq \gamma_0}] \delta_n = 0$.}

\footnote{In the nonlinear scenario, uncorrelatedness in the linear scenario should be strengthened to independence. Also, all elements of $(x', q)'$ should be endogenous; otherwise, $z$ should include the exogenous elements of $(x', q)'$ and cannot be independent of $(x', q)'$.}

\footnote{In Remark 2 of Seo and Shin (2016) where their FD-GMM estimator, which is similar to our estimator, is used to estimate the dynamic panel threshold regression, they claim that the asymptotic distribution of $\hat{\theta}$ is invariant to whether the model is CTR. That statement is not correct as can be seen by noting that their $G_\gamma (\gamma_0) = 0$ in CTR so that the asymptotic variance matrix of $\hat{\theta}$ is undefined. The zerosness of $G_\gamma (\gamma_0)$ is due to some redundancy in the parameter $\theta$ when the model is CTR. If we rewrite the CTR as $y = x'\beta + (q - \gamma) \delta_1 (q \leq \gamma) + \varepsilon$, then the corresponding $G$ under the moment conditions $E [zx] = 0$ is \( \left( \mathbb{E} [zx'], \mathbb{E} [x_{\leq \gamma_0} (q - \gamma_0)] \right) \), $E [z (q - \gamma_0) | q = \gamma_0] F_q (\gamma_0) - E [z_{\leq \gamma_0}] \delta_0$), which is of full column rank.}
and when the model is homoskedastic, i.e., $\mathbb{E} [e^2 | x] = \sigma^2$, our 2SLS estimator is efficient. For inference concerning $\gamma$ we suggest use of bootstrap methods such as in Hall and Horowitz (1996), Brown and Newey (2002) or Lee (2014) to avoid estimating $\mathbb{E} [z' x | q = \gamma_0]$, and by the Delta method, $\sqrt{n} (\hat{\beta} - \beta_0) \stackrel{d}{\rightarrow} N (0, \Omega^{-1})$, where $\Omega$ is estimated using numerical derivatives, as in Section 7.3 of Newey and McFadden (1994), or by some kernel or series method.

### 3.1 Comparison with the GMM of HHB and the 2SLS of CH

In the structural change context, HHB show that the GMM estimator based on the following criterion is generally inconsistent:

$$\hat{\gamma} = \arg\min_{\gamma} \tilde{Q}_n (\gamma),$$

where

$$\tilde{Q}_n (\gamma) = \min_{\beta_1, \beta_2} \tilde{Q}_n (\theta) := \min_{\beta_1, \beta_2} \tilde{m}_n (\theta)' \tilde{W} \tilde{m}_n (\theta)$$

with

$$\tilde{m}_n (\theta) = \frac{1}{n} \sum_{i=1}^{n} m_i (\theta) = \frac{1}{n} \sum_{i=1}^{n} \left( \begin{array}{c} m_{1i} (\theta) \\ m_{2i} (\theta) \end{array} \right) = \frac{1}{n} \sum_{i=1}^{n} z_i \left( \begin{array}{c} (y_i - x_i' \beta_1) 1(q_i \leq \gamma) \\ (y_i - x_i' \beta_2) 1(q_i > \gamma) \end{array} \right).$$

As commented by HHB, inconsistency of $\hat{\gamma}$ stems from the fact that the minimand is a quadratic form in the sample moment, thereby taking the form of a square of sums. This "square of sums" structure provides an opportunity for the effects of misspecification associated with the selection of the wrong threshold point to be an offsetting balancing factor in the minimand, leading to inconsistency. In contrast, the objective function of the 2SLS estimator in CH takes a "sum of squares" form, which generates a consistent estimator of $\gamma$. Specifically, the objective function of the 2SLS of CH is

$$\tilde{S}_n (\theta) = \frac{1}{n} \sum_{i=1}^{n} \left( y_i - \beta_1 (\tilde{W} z_i) 1(q_i \leq \gamma) - \beta_2 (\tilde{W} z_i) 1(q_i > \gamma) \right)^2,$$  \(\text{(13)}\)
where \(\hat{\Pi}'z_i\) delivers a first-stage prediction of \(x_i\). Given the comments by HHB, it may seem surprising that our GMM estimator is consistent even if the mininum is also a square of sums. The key point, however, is not the distinction between the "square of sums" and "sum of squares" criteria in this case, but rather the fact that the threshold variable \(q\) in the structural change model is a time index which is independent of the other components of the system (so that offsetting is possible, resulting in inconsistency).

Before a formal discussion on these points, note first that our 2SLS estimator is a special GMM estimator of HHB. Specifically, it is easy to check that when

\[
\hat{W} = \begin{pmatrix} I_t & \hat{W} \end{pmatrix},
\]

\(\hat{Q}_n(\theta) = \hat{Q}_n(\theta)\). This \(\hat{W}\) is only positive semidefinite, not positive definite. In other words, our 2SLS estimator does not fully explore the information in \(\tilde{m}_n(\theta)\). This is why we need \(l > 2\tilde{d}\) instruments, whereas HHB’s GMM estimator needs only \(l > \tilde{d}\) instruments. This is also why our 2SLS estimator is not consistent when \(q\) has properties like a time index (because the general GMM estimator is not consistent). The moment conditions in \(\tilde{m}_n(\theta)\) explore the special structure of threshold regression - \(\beta_1\) and \(\beta_2\) are involved only in one regime of the system, while the moment conditions in \(\tilde{g}_n(\theta)\) are designed for any nonlinear system \(y = G(x, q; \theta) + \varepsilon\) with \(E[\varepsilon|x, q] \neq 0\). Essentially, the moment conditions \(\tilde{m}_n(\theta)\) explore the validity of the moments \(E[z_{\leq \gamma_0}] = 0\) and \(E[z_{> \gamma_0}] = 0\), whereas \(\tilde{g}_n(\theta)\) explores only \(E[z] = E[z_{\leq \gamma_0}] + E[z_{> \gamma_0}] = 0\).

We can now formally state the consistency of \(\hat{\gamma}\) when \(q\) is not independent of \((z', \vec{x}', \varepsilon)\). First, we impose the following assumption. Because we concentrate on the identification issue below, we here assume that \(\delta\) is fixed for notational simplicity.

**Assumption IV’**: \(\dim(z) = l \geq \tilde{d} + 1, \hat{W} \overset{p}{\to} W > 0\) and \(E[zx_{\leq \gamma}]\) and \(E[zx_{> \gamma}]\) are of full column rank for any \(\gamma \in \Gamma\). (i) If \(q\) is exogenous (i.e., \(q\) is included in \(z\), and \(E[\varepsilon|z] = 0\)), then there does not exist \(a = (a_1', a_2')' \in \mathbb{R}^{2\tilde{d}}\) such that \(E[zx_{> \gamma}] a_2 = E[zx_{> \gamma}] \delta_0\) for any \(\gamma < \gamma_0\) or \(E[zx_{\leq \gamma}] a_1 = E[zx_{\leq \gamma}] \delta_0\) any \(\gamma > \gamma_0\). (ii) If \(q\) is endogenous and only \(E[z_{\leq \gamma_0}] = 0\) and \(E[z_{> \gamma_0}] = 0\) hold, then there does not exist \(a = (a_1', a_2')' \in \mathbb{R}^{2\tilde{d}}\) such that \(E[zx_{\leq \gamma}] a_1 = E[z_{\leq \gamma}], E[zx_{> \gamma}] \delta_0 + E[zx_{> \gamma}] a_2 = E[z_{> \gamma}]\) for any \(\gamma < \gamma_0\) or \(E[zx_{\leq \gamma}] a_1 - E[zx_{\leq \gamma}] \delta_0 = E[z_{\leq \gamma}]\) and \(E[zx_{> \gamma}] a_2 = E[z_{> \gamma}]\) for any \(\gamma > \gamma_0\).

**Theorem 2** Under Assumptions F, I’, IV’ and S, \(\hat{\gamma}\) is consistent.

HHB assume \(\hat{W} = \text{diag}(\hat{W}_1, \hat{W}_2)\) with \(\hat{W}_1 \overset{p}{\to} W_1 > 0\) and \(\hat{W}_2 \overset{p}{\to} W_2 > 0\), but we do not need such a restriction to show the consistency of \(\hat{\gamma}\) or inconsistency of \(\hat{\gamma}\) in the HHB setup. Similar to \(\hat{\gamma}\), we only require Assumption I’ rather than the stronger Assumption I to prove the consistency of \(\hat{\gamma}\). In contrast to Assumption IV, we need different assumptions here for the identification of \(\gamma_0\) depending on whether \(q\) is exogenous or not. In this sense, reducing \(\tilde{m}_n(\theta)\) to \(\tilde{g}_n(\theta)\) makes the treatment of identification more uniform although there is some loss of information in doing so. When \(q\) is exogenous, Assumption IV’(i) requires some extra variation in \(E[zx|q = \gamma]\) when \(\gamma\) moves away from \(\gamma_0\). This condition implicitly precludes the possibility that \(q\) is independent of \((z', \vec{x}')\) because if this is the case, then \(E[zx_{> \gamma}] = E[zx'] (1 - F_q(\gamma))\) and \(E[zx_{\leq \gamma}] = E[zx'] (1 - F_q(\gamma_0))\), so \(a_2\) can be chosen as \(\frac{F_q(\gamma_0)}{1 - F_q(\gamma)} \delta_0\) and, similarly, \(a_1\) can be chosen as \(\frac{F_q(\gamma_0)}{1 - F_q(\gamma)} \delta_0\). When \(q\) is endogenous, we need also to take account of the variation in \(E[z_{\leq \gamma}]\) and \(E[z_{> \gamma}]\) as \(\gamma\) moves away from \(\gamma_0\). We provide more intuition on such identifying information in the following discussion.

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14 Following the general setup of this paper we do not assume a threshold effect in the first stage.

15 To save notation, we still use \(W\) to denote the limit of \(\hat{W}\). This should not introduce any confusion.
The proof of the theorem also establishes the following results. First, if \( q \) is exogenous and also independent of \( (z', x')' \), then \( \gamma_0 \) cannot be identified by \( \tilde{Q}_n(\gamma) \). This is essentially the case considered by HHB, and we label it case (o). Second, in case (i), both groups of moment conditions in \( m_i(\theta) \) are required to identify \( \gamma_0 \); using only \( m_{1i}(\theta) \) or \( m_{2i}(\theta) \) is not enough. Third, in case (ii), either group of moment conditions in \( m_i(\theta) \) can identify \( \gamma_0 \). For example, if there does not exist \( a_1 \in \mathbb{R}^d \) such that \( E[z_{x<\gamma}] a_1 = E[z_{x<\gamma}] \) for any \( \gamma < \gamma_0 \) and \( E[z_{x<\gamma}] a_1 = E[z_{x<\gamma}] \delta_0 = E[z_{x<\gamma}] \) for any \( \gamma > \gamma_0 \), then \( \gamma_0 \) can be identified by only \( m_{1i}(\theta) \). In comparison with case (i), we can see that the identifying power in either group of moment conditions in \( m_i(\theta) \) comes solely from the correlation between \( q \) and \( \varepsilon \). In other words, endogeneity is helpful in identifying \( \gamma_0 \) by moment conditions.

It seems that the correlation of \( q \) with the rest of the system is critical for the identification of \( \gamma_0 \). When \( q \) is independent of \( (z', x', \varepsilon)' \), then even the combination of \( m_1 \) and \( m_2 \) cannot identify \( \gamma_0 \); if \( q \) is independent of \( \varepsilon \) but not \( (z', x')' \), then combination of \( m_1 \) and \( m_2 \) can (but \( m_1 \) or \( m_2 \) individually cannot) identify \( \gamma_0 \); if \( q \) is correlated with all of \( (z', x', \varepsilon)' \), then either \( m_1 \) or \( m_2 \) can identify \( \gamma_0 \). What is the intuition here? We can understand these results by using Lemma 2.3 of Newey and McFadden (1994) which states that as \( \gamma \) is independent of \( (z', x', \varepsilon)' \), then \( m_1 \) and \( m_2 \) cannot identify \( \gamma_0 \). When \( q \) is correlated with all of \( (z', x', \varepsilon)' \), then either \( m_1 \) or \( m_2 \) can identify \( \gamma_0 \).

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when $\gamma < \gamma_0$, and equal to

$$
\left( \begin{array}{c}
E[zx'_{\leq \gamma_0}] \\
E[zx'_{\gamma_0 < \leq \gamma}] \\
E[zx'_{\gamma < \gamma}] \\
E[z\varepsilon_{\leq \gamma}]
\end{array} \right) \beta_{10} + \left( \begin{array}{c}
E[zx'_{\leq \gamma_0}] \\
E[zx'_{\gamma_0 < \leq \gamma}] \\
E[zx'_{\gamma < \gamma}] \\
E[z\varepsilon_{\leq \gamma}]
\end{array} \right) \beta_{20} + E[z\varepsilon_{\gamma}]
$$

when $\gamma > \gamma_0$. In case (o), $E[z\varepsilon_{\leq \gamma}] = E[z\varepsilon_{>\gamma}] = 0$, $E[zx'_{\leq \gamma_0}] = E[zx']F_\gamma(\gamma)$, $E[zx'_{\gamma < \gamma}] = E[zx'](1 - F_\gamma(\gamma))$ and $E[zx'_{\gamma_0 < \leq \gamma_2}] = E[zx'](F_\gamma(\gamma_2) - F_\gamma(\gamma_1))$, where $\gamma_1 < \gamma_2$. So we can choose $\beta_1$ and $\beta_2$ such that

$$
\beta_1 = \beta_{10} + (1 - F_\gamma(\gamma))\beta_2 = (F_\gamma(\gamma_0) - F_\gamma(\gamma))\beta_{10} + (1 - F_\gamma(\gamma_0))\beta_{20}
$$

(15)

when $\gamma < \gamma_0$ and

$$
\beta_2 = \beta_{20} \text{ and } F_\gamma(\gamma)\beta_1 = F_\gamma(\gamma_0)\beta_{10} + (F_\gamma(\gamma) - F_\gamma(\gamma_0))\beta_{20}
$$

(16)

when $\gamma > \gamma_0$ to make the equalities hold.

In other words, $\text{plim} \tilde{Q}_n(\gamma) = 0$ for any $\gamma$. In case (i), $E[z\varepsilon_{\leq \gamma}] = E[z\varepsilon_{>\gamma}] = 0$. So when $\gamma < \gamma_0$, we can choose $\beta_1 = \beta_{10}$ to make $E[zx'_{\leq \gamma}]\beta_1 = E[z\varepsilon_{\leq \gamma}]$ but cannot choose $\beta_2$ such that $E[zx'_{\gamma < \gamma}]\beta_2 = E[z\varepsilon_{>\gamma}]$, and when $\gamma > \gamma_0$, we can choose $\beta_2 = \beta_{20}$ to make $E[zx'_{\gamma_0 < \leq \gamma_2}]\beta_2 = E[z\varepsilon_{\leq \gamma}]$ but cannot choose $\beta_1$ such that $E[zx'_{\gamma < \gamma}]\beta_1 = E[z\varepsilon_{>\gamma}]$. In other words, if we use only $m_1$, then $\text{plim} \tilde{Q}_n(\gamma) = 0$ for $\gamma \in [\gamma_0, \gamma_1]$ and if we use only $m_2$, then $\text{plim} \tilde{Q}_n(\gamma) = 0$ on $[\gamma_0, \gamma_1]$, while if we use both $m_1$ and $m_2$, then $\text{plim} \tilde{Q}_n(\gamma) = 0$ only if $\gamma = \gamma_0$. In case (ii), $E[z\varepsilon_{\leq \gamma}] \neq 0$ and $E[z\varepsilon_{>\gamma}] \neq 0$. So even if we use only $m_1$ or $m_2$, the equalities can hold only at $\gamma = \gamma_0$.

The above arguments also show a key difference between the identification sources of the HHB GMM estimator and the CH 2SLS estimator. In CH,

$$
\text{plim} \tilde{S}_n(\theta) = E \left[ (y - \beta_1^a (\Pi' z) 1(q \leq \gamma) - \beta_2^a (\Pi' z) 1(q > \gamma))^2 \right],
$$

which assumes that $E[y|z, q] = \beta_1^a (\Pi' z) 1(q \leq \gamma) + \beta_2^a (\Pi' z) 1(q > \gamma)$ and uses the conditional mean difference of $y$ below $\gamma_0$ and above $\gamma_0$ to identify $\gamma_0$ (just as in standard least squares estimation where $E[z|x] = 0$ and we can calibrate $y$ against its conditional mean to identify the parameters in the conditional mean). Since $y = \beta_1^a (\Pi' z + u) 1(q \leq \gamma) + \beta_2^a (\Pi' z + u) 1(q > \gamma) + \varepsilon$, where the first stage regression is assumed to be $x = \Pi' z + u$, we must assume $E[u|z, q] = 0$ and $E[\varepsilon|z, q] = 0$ and then the conditional mean of $y$ is $\beta_1^a (\Pi' z) 1(q \leq \gamma) + \beta_2^a (\Pi' z) 1(q > \gamma)$. To achieve such conditions, we must assume that $q$ is exogenous so that it can be included in $z$. Also, as argued in Yu (2013a), the first stage must be a regression rather than a projection, i.e., $E[u|z] = 0$ rather than only $E[zu'] = 0$. In a nonlinear environment, such a requirement does not seem too stringent. On the contrary, the identification of $\gamma_0$ by HHB’s GMM is based on the matching of covariances just as in the usual linear GMM estimation. If $\gamma_0$ were known, we can identify $\beta_1$ by matching $E[z\varepsilon_{\leq \gamma_0}]$ with $E[zx'_{\leq \gamma_0}]$, and $\beta_2$ by matching $E[z\varepsilon_{>\gamma_0}]$ with $E[zx'_{>\gamma_0}]$. It is the nonlinear structure introduced by the unknown $\gamma$ that necessitates the division of identification into three different cases; in such a nonlinear system, endogeneity of $q$ is helpful rather than harmful to identification as in CH’s 2SLS. As far as inference is concerned, the bootstrap is questionable for CH’s 2SLS given the negative findings in Yu (2014) where it is shown that bootstrap inference for $\gamma$ is invalid when the objective function takes the "sum of squares" form.

As for the requirement on the number of instruments, CH’s assumption that $E[\Pi' z z' \Pi|q = \gamma_0] = \Pi' E[z z'|q = \gamma_0] \Pi > 0$ implies that $l \geq 7$ instruments are needed. As mentioned in Assumptions IV and IV', $l \geq 2d + 1$ instruments are required in our 2SLS approach because $g_\theta(\theta)$ contains $2d + 1$ unknown parameters, and $l \geq 7 + 1$ instruments are required in HHB’s GMM approach because each of $m_{11}(\theta)$ and

\footnote{Note that we can choose $\beta_1$ and $\beta_2$ freely, so the choice of $\beta_1$ and $\beta_2$ depends on $\gamma$ here.}

18
$m_{2i}(\theta)$ contains $d + 1$ unknown parameters. On the other hand, more instruments typically imply greater identification power: CH’s 2SLS cannot handle the endogenous $q$ case, whereas the other two estimators can; and, as discussed before Assumption IV’, HBB’s GMM relies on the special structure of (1) whereas our 2SLS approach can handle any nonlinear system $y = G(x, q; \theta) + \varepsilon$ with $\mathbb{E}[\varepsilon|x, q] \neq 0$.

It is well known that more moment conditions generally imply higher asymptotic efficiency. Why then is our 2SLS the suggested approach rather than HBB’s GMM? The reason is that the derivative $d\mathbb{E}[m_i(\theta_0)]/d\theta'$ does not exist as is normally required in usual GMM asymptotic derivations. Specifically,

$$
\left. \frac{\partial \mathbb{E}[m_i(\beta_0, \delta_0, \gamma)]}{\partial \gamma} \right|_{\gamma=\gamma_0^+} = \left( \mathbb{E}[xz|q = \gamma_0]f_q(\gamma_0) - \mathbb{E}[zx'|q = \gamma_0]\delta_0f_q(\gamma_0) \right),
$$

whereas

$$
\left. \frac{\partial \mathbb{E}[m_i(\beta_0, \delta_0, \gamma)]}{\partial \gamma} \right|_{\gamma=\gamma_0^-} = \left( -\mathbb{E}[xz|q = \gamma_0]f_q(\gamma_0) - \mathbb{E}[zx'|q = \gamma_0]\delta_0f_q(\gamma_0) \right).
$$

Here, note that $\mathbb{E}[zx1(q \leq \gamma_0)] = 0$ and $\mathbb{E}[zx1(q > \gamma_0)] = 0$ do not imply $\mathbb{E}[zx|q = \gamma_0] = 0$. Even if $\mathbb{E}[zx|q = \gamma_0] = 0$ as in case (i), if $\mathbb{E}[zx'|q = \gamma_0]$ is of full column rank, the derivative $\partial \mathbb{E}[m_i(\beta_0, \delta_0, \gamma)]/\partial \gamma$ does not exist. This makes the asymptotic distribution of $\gamma$ a nonnormal mixture that depends on the one-sided derivatives, rendering inference based on $\gamma$ difficult. On the contrary, in our 2SLS approach we have

$$
\left. \frac{\partial \mathbb{E}[g_i(\beta_0, \delta_0, \gamma)]}{\partial \gamma} \right|_{\gamma=\gamma_0^+} = (I_1, I_1) \left. \frac{\partial \mathbb{E}[m_i(\beta_0, \delta_0, \gamma)]}{\partial \gamma} \right|_{\gamma=\gamma_0^+} - \mathbb{E}[zx'|q = \gamma_0] \delta_0f_q(\gamma_0),
$$

which makes bootstrapping inference valid.

We close this subsection with a further comment on identification based on moment conditions. As mentioned above, more moment conditions typically imply higher asymptotic efficiency, but such results rely in the first place on identification. With 2SLS, even if $\mathbb{E}[zx] = 0$ is replaced by $\mathbb{E}[\varepsilon|z] = 0$ which implies more moment conditions, the identification results are unaltered. For example, when $q$ is independent of $(z', x')'$, $\gamma_0$ is not identified even by $\mathbb{E}[\varepsilon|z] = 0$. Similar identification results apply to HBB’s GMM.

### 3.2 A Simple Illustration

We illustrate the identification results above based on the example in Section 2.1. In this example $y = 1(q \leq \gamma_0) + \varepsilon$, where $q \sim U[0, 1]$, $\gamma_0 = 1/2$, $\beta_0 = 0$ and $\delta_0 = 1$ are known, $x = 1$, and $\text{Var}(\varepsilon) = 1$.

First assume $\mathbb{E}[\varepsilon|q] = 0$, so there is no endogeneity. Let $z = 1$, giving case (o) with $\mathbb{E}[\varepsilon|q, z] = 0$ and $q \perp (z, x)$. The moment conditions used for identifying $\gamma_0$ are

$$
\mathbb{E}
\left[
\begin{pmatrix}
(y - 1)1(q \leq \gamma) \\
y1(q > \gamma)
\end{pmatrix}
\right] = 0.
$$

Suppose $\tilde{W} = I_2$. Then

$$
\text{plim} Q_n(\gamma) = \mathbb{E}[(y - 1)(q \leq \gamma)]^2 + \mathbb{E}[y1(q > \gamma)]^2.
$$

---

We will discuss the asymptotic properties of $\gamma$ and bootstrap inference based on $\gamma$ in a separate paper.
After some algebra,
\[
\text{plim} \tilde{Q}_n(\gamma) = \left[ -\left( \gamma - \frac{1}{2} \right)_+ \right]^2 + \left( \frac{1}{2} - \gamma \right)_+^2 = \left( \frac{1}{2} - \gamma \right)_+^2,
\]
where for \( a \in \mathbb{R}, a_+ = \max(a, 0) \), and \( \frac{1}{2} - \gamma = -(\gamma - \frac{1}{2})_+ + (\frac{1}{2} - \gamma)_+ \). Obviously, \( \arg \min_{\gamma \in \mathbb{R}} \text{plim} \tilde{Q}_n(\gamma) = 1/2 \).

This seems to contradict the nonidentification result of HHB in case (o). In fact, this outcome is because \( \beta_0 \) and \( \delta_0 \) are known. In \([15]\) and \([16]\), \( \beta_1 \) and \( \beta_2 \) are fixed at \( \beta_{10} = 1 \) and \( \beta_{20} = 0 \). So when \( \gamma < \gamma_0 \),
\[
(1 - F_q(\gamma)) \beta_{20} = (F_q(\gamma_0) - F_q(\gamma)) \beta_{10} + (1 - F_q(\gamma_0)) \beta_{20} \text{ or } (F_q(\gamma_0) - F_q(\gamma)) \delta_0 = 0
\]
cannot hold as long as \( \delta_0 \neq 0 \) and \( f_q(\gamma) > 0 \) on \([\gamma_0, \gamma]\). Similarly, when \( \gamma > \gamma_0 \), \( F_q(\gamma) \beta_{1} = F_q(\gamma_0) \beta_{10} + (F_q(\gamma) - F_q(\gamma_0)) \beta_{20} \text{ or } (F_q(\gamma) - F_q(\gamma_0)) \delta_0 = 0 \) cannot hold as long as \( \delta_0 \neq 0 \) and \( f_q(\gamma) > 0 \) on \([\gamma_0, \gamma]\). If they can be chosen freely, then it is obvious that the system cannot be identified - there are two equations and three unknowns.

Next, let \( z = (1, q)' \) as in case (i), for which the moment conditions used for identifying \( \gamma_0 \) are
\[
\mathbb{E} \left[ \begin{pmatrix} z(y - 1) \mathbb{I}(q \leq \gamma) \\ zy \mathbb{I}(q > \gamma) \end{pmatrix} \right] = 0.
\]

Suppose \( \tilde{W} = \begin{pmatrix} \tilde{W}_1 & \tilde{W}_2 \\ \tilde{W}_{12} & \tilde{W}_2 \end{pmatrix} = \begin{pmatrix} c_1 I_2 & \tilde{W}_{12} \\ \tilde{W}_{12} & c_2 I_2 \end{pmatrix}, \)

and then
\[
\text{plim} \tilde{Q}_n(\gamma) = c_1 \left[ -\left( \gamma - \frac{1}{2} \right)_+ \right]^2 + c_1 \left[ -\left( \frac{1}{2} - \gamma \right)_+ \right]^2 + c_2 \left( \frac{1}{2} - \gamma \right)_+^2 + c_2 \left( \frac{1}{2} - \gamma \right)_+^2,
\]
\[
+ 2 \left( -\left( \gamma - \frac{1}{2} \right)_+ , -\left( \frac{1}{2} - \gamma \right)_+ \right) \tilde{W}_{12} \left( -\left( \gamma - \frac{1}{2} \right)_+ , -\left( \frac{1}{2} - \gamma \right)_+ \right)'.
\]

If \( \tilde{W}_{12} = 0 \), \( c_1 = 1 \) and \( c_2 = 0 \), then we use only \( m_1(\gamma) \) and Figure 1 shows that \( \arg \min_{\gamma \in \mathbb{R}} \text{plim} \tilde{Q}_n(\gamma) = [0, 1/2] \).

If \( \tilde{W}_{12} = 0 \), \( c_1 = 0 \) and \( c_2 = 1 \), then we use only \( m_2(\gamma) \) and Figure 1 shows that \( \arg \min_{\gamma \in \mathbb{R}} \text{plim} \tilde{Q}_n(\gamma) = [1/2, 1] \).\(^{20}\) If \( c_1 \neq 0 \) and \( c_2 \neq 0 \), then we use both \( m_1(\gamma) \) and \( m_2(\gamma) \) and Figure 1 shows that when either \( c_1 = 1, c_2 = 2 \) and \( \tilde{W}_{12} = 0 \) or \( c_1 = 1, c_2 = 1.5 \) and \( \tilde{W}_{12} = \begin{pmatrix} 0.5 & 0.5 \\ 0.5 & 0.5 \end{pmatrix} \), \( \arg \min_{\gamma \in \mathbb{R}} \text{plim} \tilde{Q}_n(\gamma) = 1/2 \).\(^{21}\)

Section 3.3 of Yu (2015b) considers the following joint distribution of \((q, \varepsilon)\):
\[
f_{q,\varepsilon}(q, \varepsilon) = \begin{cases} \phi(\varepsilon - 1/2), & \text{if } 0 \leq q < \frac{1}{4}, \\ \phi(\varepsilon + 1/2), & \text{if } \frac{1}{4} \leq q \leq \frac{1}{2}, \\ \phi(\varepsilon - 1/2), & \text{if } \frac{1}{2} < q \leq \frac{3}{4}, \\ \phi(\varepsilon + 1/2), & \text{if } \frac{3}{4} < q \leq 1, \\ 0, & \text{otherwise}, \end{cases}
\]
where \( \phi(\cdot) \) is the standard normal density. Obviously, \( \mathbb{E}[\varepsilon | q] \neq 0 \), and \( \mathbb{E}[\mathbb{I}(q \leq \gamma_0)] = \mathbb{E}[\mathbb{I}(q > \gamma_0)] = 0 \), so this is case (ii). Suppose \( z = 1 \). Then the moment conditions used for identifying \( \gamma_0 \) are \([19]\). Suppose
\[^{20}\]This implies that in case (o) (i.e., \( z = 1 \)), using only one moment condition cannot identify \( \gamma_0 \).
\[^{21}\]In this example, \( \tilde{W}_{12} \) does not play any role, i.e., \( \text{plim} \tilde{Q}_n(\gamma) \) depends only on \( \tilde{W}_1 \) and \( \tilde{W}_2 \) because the last term of \( \text{plim} \tilde{Q}_n(\gamma) \) is zero.
\[ \mathbf{\tilde{W}} = \begin{pmatrix} A & C \\ C & B \end{pmatrix}. \]

Then
\[ \text{plim} \tilde{Q}_n(\gamma) = A E [(y - 1)1(q \leq \gamma)]^2 + B E [y1(q > \gamma)]^2 + 2CE [(y - 1)1(q \leq \gamma)] E [y1(q > \gamma)], \]

and with some algebra this reduces to
\[ \text{plim} \tilde{Q}_n(\gamma) = \begin{cases} A \left( \frac{3\gamma}{2} \right)^2 + B \left( \frac{1-3\gamma}{2} \right)^2 + 2C \left( \frac{3\gamma - 1}{2} \right), & \text{if } 0 \leq \gamma < \frac{1}{4}, \\ (A + B + 2C) \left( \frac{1}{2} \right)^2, & \text{if } \frac{1}{4} \leq \gamma \leq \frac{3}{4}, \\ A \left( 1 - \frac{3\gamma}{2} \right)^2 + B \left( \frac{3\gamma - 1}{2} \right)^2 + 2C \left( 1 - \frac{3\gamma}{2} \right), & \text{if } \frac{3}{4} < \gamma \leq 1, \end{cases} \]

Figure 2 shows that when \( A = 1, B = 2 \) and \( C = 0 \), \( \arg \min_{\gamma \in \Gamma} \text{plim} \tilde{Q}_n(\gamma) = 1/2 \). Actually, if only \( m_1 \) is used (i.e., \( A = 1, B = 0 \) and \( C = 0 \)) or only \( m_2 \) is used (i.e., \( A = 0, B = 1 \) and \( C = 0 \)), \( \arg \min_{\gamma \in \Gamma} \text{plim} \tilde{Q}_n(\gamma) = 1/2 \), where the parameter space \( \Gamma \) excludes the neighborhoods of 0 and 1. If \( C \neq 0 \), \( \arg \min_{\gamma \in \Gamma} \text{plim} \tilde{Q}_n(\gamma) = \gamma_0 \) as long as \( \mathbf{\tilde{W}} > 0 \). For instance, Figure 2 shows that when \( A = 1, B = 2 \) and \( C = 1 \), \( \arg \min_{\gamma \in \Gamma} \text{plim} \tilde{Q}_n(\gamma) = 1/2 \).

We now check the behavior of the 2SLS estimators of this paper and CH for these three cases: (o) \( E[x;q] = 0 \) and \( z = 1 \); (i) \( E[x;q] = 0 \) and \( z = (1,q)^T \); and (ii) \( E[x;q] \neq 0 \) and \( z = 1 \). For 2SLS, the moment conditions are \( E[y - 1(q \leq \gamma)] = 0 \) in cases (o) and (ii) and \( E[y - 1(q \leq \gamma)] = 0 \) in case (i). In case (o),
\[ \text{plim} \tilde{Q}_n(\gamma) = \left( \frac{1}{2} - \gamma \right)^2, \]

the same as \( \text{plim} \tilde{Q}_n(\gamma) \). In case (i) with \( \mathbf{\tilde{W}} = I_2 \),
\[ \text{plim} \tilde{Q}_n(\gamma) = \left( \frac{1}{2} - \gamma \right)^2 + \left( \frac{1}{2} - \gamma \right)^2, \]

where for comparison with \( \text{plim} \tilde{Q}_n(\gamma) \), note that \( \frac{1}{8} - \frac{1}{2} \gamma^2 = - \left( \frac{1}{2} \gamma^2 - \frac{1}{8} \right) \downarrow \left( \frac{1}{8} - \frac{1}{2} \gamma^2 \right) \). In case (ii),
\[ \text{plim} \tilde{Q}_n(\gamma) = \left( \frac{1}{2} - \gamma \right)^2, \]

the same as in case (o). For CH’s 2SLS,
\[ \text{plim} \tilde{S}_n(\gamma) = E \left[ (y - 1(q \leq \gamma))^2 \right] \]

in all three cases. In cases (o) and (i),
\[ \text{plim} \tilde{S}_n(\gamma) = 1 + |\gamma - 1/2| ; \]

in case (ii),
\[ \text{plim} \tilde{S}_n(\gamma) = \begin{cases} \frac{7}{4} - 2\gamma, & \text{if } 0 \leq \gamma < \frac{1}{4}, \\ \frac{5}{4}, & \text{if } \frac{1}{4} \leq \gamma \leq \frac{3}{4}, \\ 2\gamma - \frac{1}{4}, & \text{if } \frac{3}{4} < \gamma \leq 1. \end{cases} \]

Figure 3 shows \( \text{plim} \tilde{Q}_n(\gamma) \) and \( \text{plim} \tilde{S}_n(\gamma) \) in these three cases. For our 2SLS estimator, \( \arg \min_{\gamma \in \Gamma} \text{plim} \tilde{Q}_n(\gamma) = 1/2 \) in all cases, giving the same identifying results as HHB’s GMM using both \( m_1 \) and \( m_2 \). For CH’s 2SLS,
Figure 1: $\text{plim} Q_n(\gamma)$ when $E[\varepsilon_1 q] = 0$ and $z = (1, q)'$

arg min $\text{plim} \tilde{S}_n(\gamma) = 1/2$ in cases (o) and (i), whereas arg min $\text{plim} \tilde{S}_n(\gamma) = [1/4, 3/4]$ in case (ii), which is unidentified.

We next check whether the expected moment conditions are differentiable. In case (i), for HBB’s GMM, $E[m_i(\beta_0, \delta_0, \gamma)] = 
\begin{pmatrix}
- (\gamma - \frac{1}{2})_+
- (\frac{1}{2} \gamma^2 - \frac{1}{8})_+
(\frac{1}{2} - \gamma)_+
(\frac{1}{8} - \frac{1}{2} \gamma^2)_+
\end{pmatrix}
$ is not differentiable at $\gamma_0 = 1/2$, whereas for our 2SLS $E[g_i(\beta_0, \delta_0, \gamma)] = \left( \frac{1}{2} - \gamma \right)$, which is differentiable at $\gamma_0 = 1/2$. In case (ii), $E[m_i(\beta_0, \delta_0, \gamma)]$ is differentiable at $\gamma_0$. This is due to the special design of the DGP in this simple example. Specifically, using the formulae in (17) and (18), it turns out that $\frac{\partial E[m_i(\beta_0, \delta_0, \gamma)]}{\partial \gamma}$ is $-1/2$ at $\gamma_0 = 1/2$. But, in general, $E[m_i(\beta_0, \delta_0, \gamma)]$ is not differentiable at $\gamma_0$ in case (ii). On the other hand, for our 2SLS, $E[g_i(\beta_0, \delta_0, \gamma)] = \frac{1}{2} - \gamma$, which is always differentiable at $\gamma_0$.

3.3 Summary of Identification Results

Before summarizing the identification results for the existing estimators of $\gamma_0$, we provide a further comment on the distinction between "sum of squares" and "square of sums" criteria. Note that "sum of squares" criteria need not have more identification power than "square of sums" criteria. When $q$ is endogenous, the example in Section 2.1 of Yu (2013a) shows that the 2SLS estimator of CH is not consistent, and the example in the previous subsection shows that the limit objective function of CH’s 2SLS need not even have a unique minimizer. On the contrary, either the 2SLS estimator of this paper or the GMM estimator of HBB can
Figure 2: \( \text{plim} \hat{Q}_n (\gamma) \) when \( E [\varepsilon | q] \neq 0 \)

Figure 3: \( \text{plim} \hat{Q}_n (\gamma) \) and \( \text{plim} \hat{S}_n (\gamma) \)
generate a consistent estimator of $\gamma_0$.  

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Table 1: Identification of $\gamma_0$ by Various (Possibly Valid) Estimators in Different Scenarios

Table 1 summarizes the identification results for all possibly consistent estimators of $\gamma_0$ in various scenarios. The first four estimators require instruments and the last two do not. Among the last two, Perron and Yamamoto (2015) (PY in Table 1) use the LSE to estimate $\gamma$ in a structural change model even when there is endogeneity. However, as shown in Yu (2015a), this strategy is valid only in the structural change context. From Table 1, it seems that the IDKE of YP has the most extensive identification power even though the method makes no use of instruments. Nevertheless, the IDKE cannot identify $\gamma_0$ in CTR models while HHB’s GMM estimator and our 2SLS estimator can identify $\gamma_0$ even in such models (although inference needs further investigation). Table 1 also lists KST’s STR estimator. As mentioned in the Introduction, when $q$ is exogenous, their estimator is equivalent to CH’s 2SLS estimator so it is consistent; but when $q$ is endogenous, their estimator is not generally consistent unless the endogeneity is relatively small compared to the threshold effect. Taking Table 1 as a whole, we can see some interesting differences between structural change models and TR models – specifically case (o) vs. cases (i) and (ii). It is, however, becoming folklore in the literature that these two kinds of models are considered similar to each other (at least in terms of their asymptotic properties). The present findings reveal that such folklore is misleading when there is endogeneity.

---

22 Of course, we can claim that CH’s 2SLS cannot be applied when $q$ is endogenous; see YLP for modifications of CH’s 2SLS to generate consistent estimators of $\gamma_0$.

23 Here, we implicitly assume $z$ and $x$ do not include $q$.

24 Both $m_1(\theta)$ and $m_2(\theta)$ are required to prove consistency.

25 Either $m_1(\theta)$ or $m_2(\theta)$ is enough to prove consistency.

26 Yu (2015a) strengthens this result a little. Specifically, let $z = (g', q)'$ and $x = (x, q)'$. If $q \perp z$ and $E[x|x, q] = g(z) + q\lambda$, i.e., $q$ need not be independent of $x$, then projecting $x$ only on $z$ in the first stage would generate a consistent estimator of $\gamma_0$.

27 $z$ is not necessary here.

28 Yu (2015a)’s result is a little stronger. Specifically, if $q \perp x$ and $E[x|x, q] = g(x) + q\lambda$, i.e., $q$ need not be exogenous, then the LSE is consistent.

29 In such models, the IDKE can be extended also to take into account slope differences at each $\gamma \in \Gamma$ beyond level differences to identify $\gamma_0$.

30 In structural change models, case (i) corresponds to the circumstance that the moments $E[(z_t', x_t')]$ are not equal for all $t$, i.e., that there is some nonstationarity in the mean of $E[(z_t', x_t')]$; case (ii) corresponds to $E[z_t|x_t] \neq 0$ for all $t$ but with $\frac{1}{T} \sum_{t=1}^{T_0} E[z_t x_t] = \frac{1}{T-T_0} \sum_{t=T_0+1}^{T} E[z_t x_t] = 0$, where $T_0$ is the break point, i.e., $z_t$ is not a valid instrument for all $t$ but is valid when the information in each regime is integrated. These results echo the finding in YP that nonstationarity is often helpful in establishing identification.
4 Inference Based on the IDKE with \( k_{\pm}(0) = 0 \)

This section presents limit theory for the IDKE \( \hat{\gamma} \) in Method II where \( k_{\pm}(0) = 0 \). To facilitate formulation of the limit distribution of \( \hat{\gamma} \), we define the following quantities,

\[
\Delta_i = \mathbb{E}[y_i|x_i, q_i = \gamma_0] - \mathbb{E}[y_i|x_i, q_i = \gamma_0^+] =: m_-(x_i) - m_+(x_i),
\]

\[
\Delta f(x_i) = \Delta_i \cdot f(x_i, \gamma_0),
\]

where \( \Delta_i = \Delta (x_i, \gamma_0) \) and \( \Delta f(x_i) \) is the limit of \( \hat{\Delta}_i (\gamma_0) \) with \( \hat{\Delta}_i (\gamma) \) and \( \Delta (x, \gamma_0) \) defined in (5) and (6) respectively. To derive the asymptotic distribution of \( \hat{\gamma} \) we use the following assumptions on \( f(u|x, q) \), which is allowed to be discontinuous at \( q = \gamma_0 \).

**Assumption U:**

(a) \( f(u|x, q) \) is continuous in \( u \) for \( (x', q) \in X \times \Gamma_\epsilon^- \) and \( (x', q') \in X \times \Gamma_\epsilon^+ \), where \( \Gamma_\epsilon^- = (\gamma - \epsilon, \gamma_0] \) and \( \Gamma_\epsilon^+ = (\gamma_0, \gamma + \epsilon) \) for some \( \epsilon > 0 \).

(b) \( f(u|x, q) \) is Lipschitz in \( (x', q') \) for \( (x', q) \in X \times \Gamma_\epsilon^- \) and \( (x', q') \in X \times \Gamma_\epsilon^+ \).

(c) \( \mathbb{E}[u^2|x, q] \) is uniformly bounded on \( (x', q') \in X \times \Gamma_\epsilon \), where \( \Gamma_\epsilon = \Gamma_\epsilon^- \cup \Gamma_\epsilon^+ \).

Given Assumption U, we impose the following conditions on the bandwidth \( h \).

**Assumption H:** \( h \to 0 \), and \( \sqrt{n}h^d / \ln n \to \infty \).

Observe that \( nh^d = \sqrt{n} \ln n \frac{\sqrt{nh^d}}{\ln n} \to \infty \) when \( \sqrt{n}h^d / \ln n \to \infty \). The limit theory for \( \hat{\gamma} \) is given in the next result.

**Theorem 3** Under Assumptions F, G, H, I, K, S and U,

\[
\sqrt{n/h} (\hat{\gamma} - \gamma_0) \xrightarrow{d} N(0, \Sigma)
\]

where

\[
\Sigma = \frac{\mathbb{E}[\Delta_i^2 f(x_i) f^2(x_i) (\sigma_x^2(x_i) + \sigma_u^2(x_i)) | q_i = \gamma_0] \xi(1)}{f_q(\gamma_0) (\mathbb{E}[\Delta f(x_i) \Delta_i f(x_i) | q_i = \gamma_0])^2 k_+(0)^2}
\]

with \( \xi(1) = \int_0^1 k_+(t)^2 dt \) and \( \sigma_x^2(x) = \mathbb{E}[u^2|x, q = \gamma_0 \pm] \).

This result shows that \( \hat{\gamma} \) converges to \( \gamma_0 \) at the rate \( \sqrt{n/h} \), a much faster rate than that of the DKE \( \tilde{\gamma} \) of DH because \( \tilde{\gamma} \) utilizes more data information in estimation. Specifically, the convergence rate of DKE is \( \sqrt{n}h^{d-2} \) and the relative rate \( \sqrt{n}h^{d-2} / \sqrt{n/h} = \sqrt{h^{d-1}} \to 0 \). Based on Theorem 2 of DH, the asymptotic variance of their estimator \( \tilde{\gamma} \) is

\[
\Sigma_o = \frac{(\sigma_x^2(x_o) + \sigma_u^2(x_o) \kappa^2 \xi(1))}{f(x_o, \gamma_0) \Delta_o^2 k_+(0)^2}, \tag{21}
\]

where \( \kappa^2 = \int K(u_x)^2 du_x \) with \( K(u_x) = \prod_{i=1}^{d-1} k(u_{x_i}) \), and \( \Delta_o = m_-(x_o) - m_+(x_o) \) which is equal to \( 1, x_o, \gamma_0 \delta_0 \) when \( \mathbb{E}[c|x, q] \) is continuous. This asymptotic variance is comparable to \( \Sigma \), but critically relies on the choice of \( x_o \). If \( \Delta_i = \Delta \) and \( \sigma_x^2(x) = \sigma_u^2 \), then \( \Sigma = O \left( \frac{\sigma_x^2(x_o) \Delta^2 k_+(0)^2}{f_q(\gamma_0) \Delta_o^2 k_+(0)^2} \right) \). As expected, \( \Sigma \) is decreasing in \( f_q(\gamma_0) \), \( |\Delta| \) and \( k_+(0) \) and increasing in \( \sigma^2 \).

The convergence rate \( \sqrt{n/h} \) of \( \tilde{\gamma} \) exceeds the usual parametric rate \( \sqrt{n} \). To understand this increase over the parametric rate, some heuristic analysis is helpful. For this purpose, we use the simple case where
\(d = 1\), so that \(q\) is the only covariate. The convergence rate is then determined by the balance between an empirical process and a deterministic centering process. Recall that

\[
\hat{\gamma} = \arg \max_{\gamma \in \Gamma} \hat{Q}_n(\gamma) = \arg \max_{\gamma \in \Gamma} \left\{ \hat{Q}_n(\gamma) - \hat{Q}_n(\gamma_0) \right\}.
\]

Because \(\hat{\gamma}\) maximizes \(\hat{Q}_n(\gamma) - \hat{Q}_n(\gamma_0)\) on \(\Gamma\) and \(\gamma_0 \in \Gamma\), we have the decomposition

\[
0 \leq \hat{Q}_n(\hat{\gamma}) - \hat{Q}_n(\gamma_0) = \left[ Q_0(\hat{\gamma}) - Q_0(\gamma_0) \right] + \left[ \left( \hat{Q}_n(\hat{\gamma}) - Q_0(\hat{\gamma}) \right) - \left( \hat{Q}_n(\gamma_0) - Q_0(\gamma_0) \right) \right].
\]

where the first term on the extreme right side is the limit centering process and is less than zero because \(\gamma_0 = \arg \max_{\gamma \in \Gamma} Q_0(\gamma)\), whereas the second term is the modulus of continuity of the empirical process, which exceeds zero. Hence, \(Q_0(\hat{\gamma}) - Q_0(\gamma_0)\) and

\[
\sup_{|\gamma - \gamma_0| \leq \delta} \left[ \left( \hat{Q}_n(\hat{\gamma}) - Q_0(\hat{\gamma}) \right) - \left( \hat{Q}_n(\gamma_0) - Q_0(\gamma_0) \right) \right] = \frac{\delta^2 n^{(\gamma_0+h)}}{\sqrt{n}}
\]

must balance out so their sum is greater than zero.

In the \(h\) neighborhood of \(\gamma_0\), we can treat the model as a parametric one, so without loss of generality, assume

\(y_i = \Delta 1(q_i \leq \gamma_0) + u_i\),

where \(q_i = i/n, i = 1, \ldots, n\) \[u_i \sim \mathcal{N}(0, 1)\] and \(\Delta > 0\). Now, \(\hat{\gamma}\) tries to maximize \(\hat{\Delta}(\gamma) - \hat{\Delta}(\gamma_0)\), where

\[
\hat{\Delta}(\gamma) = \frac{1}{nh} \sum_{i=1}^{n\gamma} k_-(i - n\gamma \frac{1}{nh}) y_i - \frac{1}{nh} \sum_{i=1}^{n(\gamma + h)} k_+(i - n\gamma \frac{1}{nh}) y_i.
\]

For \(|\gamma - \gamma_0| \leq \delta\), and \(\gamma < \gamma_0\), we have

\[
\begin{align*}
\hat{\Delta}(\gamma) - \hat{\Delta}(\gamma_0) & = \left[ \frac{1}{nh} \sum_{i=n(\gamma - h)}^{n\gamma} k_-(\frac{i-n\gamma}{nh}) y_i - \frac{1}{nh} \sum_{i=n(\gamma - h)+1}^{n(\gamma + h)} k_+(\frac{i-n\gamma}{nh}) y_i \right] - \left[ \frac{1}{nh} \sum_{i=n(\gamma - h)}^{n\gamma} k_-(\frac{i-n\gamma}{nh}) y_i - \frac{1}{nh} \sum_{i=n\gamma+1}^{n(\gamma + h)} k_+(\frac{i-n\gamma}{nh}) y_i \right] \\
& = \Delta + \frac{1}{nh} \sum_{i=n(\gamma - h)}^{n\gamma} k_-(\frac{i-n\gamma}{nh}) y_i - \frac{1}{nh} \sum_{i=n\gamma+1}^{n(\gamma + h)} k_+(\frac{i-n\gamma}{nh}) y_i \\
& = -\frac{1}{nh} \sum_{i=n(\gamma - h)+1}^{n\gamma} k_+(\frac{i-n\gamma}{nh}) u_i + \frac{1}{nh} \sum_{i=n\gamma}^{n(\gamma + h)-1} k_-(\frac{i-n\gamma}{nh}) u_i \\
& = -\frac{1}{nh} \sum_{i=n(\gamma - h)+1}^{n\gamma} k_+(\frac{i-n\gamma}{nh}) u_i + \frac{1}{nh} \sum_{i=n\gamma}^{n(\gamma + h)-1} k_-(\frac{i-n\gamma}{nh}) u_i \\
& \approx O \left( \frac{\Delta}{nh} k_+(\gamma) dv + O_p \left( \frac{1}{nh} \int_{-\frac{1}{2}}^{-\frac{1}{2}} \left( \frac{\gamma_0 - \gamma}{n} \right)^2 dv + k_+(0) \int_{-\frac{1}{2}}^{\frac{1}{2}} \left( \frac{\gamma_0 - \gamma}{n} \right)^2 dv \right) \right)
\]

If \(\Delta\) is fixed, \(k_+(0) = 0\) and \(k_+'(0) > 0\), then \(\int_{-\frac{1}{2}}^{-\frac{1}{2}} k_+(\gamma) dv = O \left( \left( \frac{\gamma_0 - \gamma}{n} \right)^2 \right), \int_{-\frac{1}{2}}^{\frac{1}{2}} \left( \frac{\gamma_0 - \gamma}{n} \right)^2 dv = \int_{-\frac{1}{2}}^{\frac{1}{2}} \left( \frac{\gamma_0 - \gamma}{n} \right)^2 dv = \)

\[31\]Locally, \(q_i\) follows a uniform distribution on \([\gamma_0 - h, \gamma_0 + h]\), so we assume it follows the discrete form of \(U[0, 1]\) here.
The LSE is changed to \((2007)\) where a parametric TR model is considered. In the SLSE of Seo and Linton (2007), nonsmoothness is less severe (the left and right derivatives of the distribution. The extra \(h\) (in the neighborhood of \(Q\)) is still nonsmooth but the nonsmoothness is less severe (the left and right derivatives of \(Q_0(\gamma) - Q_0(\gamma_0)\) at \(\gamma_0\) are \(O(1/h)\) order in YP), so the convergence rate is slower than that in YP. For comparison, consider the IDKE in YP and in Method III. In the former, \(\int_0^{\frac{2n-\gamma}{h}} k_+ (v)dv = \frac{2n-\gamma}{h}\), \(\int_{-1}^{0} \left( \frac{2n-\gamma}{h} \right)^2 dv = \frac{2n-\gamma}{h}\), and \(\int_0^{\frac{2n-\gamma}{h}} k_+ (v) + k_- \left( \frac{2n-\gamma}{h} - v \right) \right)^2 dv = \frac{2n-\gamma}{h}\) since \(k_\pm (0) > 0\). As a result, \(Q_0(\gamma) - Q_0(\gamma_0) = O(\delta^2/h^2)\) and \(\phi_n (\delta) = \sqrt{\delta^2/h^2}\) since \(\frac{2n-\gamma}{h}\) is quadratic in \(r_n\), so solving \(\frac{\Delta}{r_n h} \approx \sqrt{\frac{\delta^2}{nh^2}}\), we get \(r_n = n\Delta^2\).

Figure 4 illustrates these heuristics. For example, in Method II, because \(Q_0(\gamma) - Q_0(\gamma_0)\) is quadratic (in the neighborhood of \(\gamma_0\)) as in the regular parameter case, we expect \(\hat{\gamma}\) to have an asymptotic normal distribution. The extra \(h\) in the convergence rate \(\sqrt{n/h}\) arises because the variations in \(Q_0(\gamma) - Q_0(\gamma_0)\) and \(\phi_n (\delta)\) are both in the scale of \(h\) in this nonparametric setup. In YP, \(Q_0(\gamma) - Q_0(\gamma_0)\) is a nonsmooth function of \(\gamma\) (in the neighborhood of \(\gamma_0\)) such that \(\gamma\) can be more easily identified than in Method II, so the convergence rate \(n\) is faster than \(\sqrt{n/h}\). In Method III, \(Q_0(\gamma) - Q_0(\gamma_0)\) is still nonsmooth but the nonsmoothness is less severe (the left and right derivatives of \(Q_0(\gamma) - Q_0(\gamma_0)\) at \(\gamma_0\) are \(O(1/h)\) order in YP), so the convergence rate is slower than that in YP. In YP and Method III, the convergence rates of \(\hat{\gamma}\) are actually the same as in the parametric cases.

For inference of \(\gamma\) based on inverting the \(t\) statistic, we need to estimate \(\Sigma\) in Theorem 3. A straightforward
approach is to use the sample analog. Specifically, we can estimate $\Sigma$ by

$$\hat{\Sigma} = \frac{\frac{1}{n} \sum_{i=1}^{n} k_h(q_i - \gamma) \hat{\Delta}_i^2(\gamma) \hat{f}_i^2(x_i) 2 \hat{u}_i^2 \xi_1}{\left( \frac{1}{n} \sum_{i=1}^{n} k_h(q_i - \gamma) \hat{\Delta}_i^2(\gamma) \hat{f}_i^{-1}(x_i, \gamma) \hat{f}(x_i) \right)^2},$$

where

$$\hat{\Delta}_i^2(\gamma) = \left( \hat{m}_-(x_i, \gamma) \hat{f}_-(x_i, \gamma) - \hat{m}_+(x_i, \gamma) \hat{f}_+(x_i, \gamma) \right),$$

$$\hat{f}(x_i) = \frac{1}{n - 1} \sum_{j=1, j \neq i}^{n} K^x_h, \hat{f}(x_i, \gamma) = \frac{1}{n - 1} \sum_{j=1, j \neq i}^{n} K^x_h k_h(q_j - \gamma),$$

$$\hat{u}_i(\gamma) = y_i - \hat{m}_-(x_i, \gamma) 1(q_i \leq \gamma) - \hat{m}_+(x_i, \gamma) 1(q_i > \gamma),$$

with

$$\hat{m}_\pm(x_i, \gamma) = \frac{1}{n - 1} \sum_{j=1, j \neq i}^{n} K^x_h, k_h^\pm(q_j - \gamma) y_j,$$

$$\hat{f}_\pm(x_i, \gamma) = \frac{1}{n - 1} \sum_{j=1, j \neq i}^{n} K^x_h k_h^\pm(q_j - \gamma).$$

The next result establishes that $\hat{\Sigma}$ is consistent.

**Theorem 4** Under the assumptions of Theorem 3, $\hat{\Sigma} \overset{p}{\rightarrow} \Sigma$.

Another method of inference is based on inverting the LR statistic. Although this method has been proposed in the small-threshold-effect framework by Hansen (2000), it seems new in the current setting. Our LR statistic can be used to test whether $\gamma = \gamma_0$ and is defined as

$$LR_n^{(1)}(\gamma) = n \hat{k}_1(0) \xi_1 \frac{\mathbb{E}[\Delta_f(x_i) \Delta_f(x_i) | q_i = \gamma_0]}{\mathbb{E}[\hat{\Delta}_f^2(x_i) f^2(x_i)(\sigma^2_f(x_i) + \sigma^2_f(x_i)|q_i = \gamma_0)]} \left( \hat{Q}_n(\gamma) - \hat{Q}_n(\gamma_0) \right).$$

**Corollary 2** Under the assumptions of Theorem 3,

$$LR_n^{(1)}(\gamma_0) \overset{d}{\rightarrow} \chi^2_1.$$ 

To construct a CI for $\gamma$ based on $LR_n^{(1)}$, we need to estimate $\frac{\mathbb{E}[\Delta_f(x_i) \Delta_f(x_i) | q_i = \gamma_0]}{\mathbb{E}[\hat{\Delta}_f^2(x_i) f^2(x_i)(\sigma^2_f(x_i) + \sigma^2_f(x_i)|q_i = \gamma_0)]}$. The natural estimator is $\frac{1}{n} \sum_{i=1}^{n} k_h(q_i - \gamma) \hat{\Delta}_i^2(\gamma) \hat{f}_i^{-1}(x_i, \gamma) \hat{f}(x_i)$, which is consistent from Theorem 4. Hence, the $(1 - \alpha)100\%$ LR-CI for $\gamma$ is

$$\left\{ \gamma : \frac{1}{n} \sum_{i=1}^{n} k_h(q_i - \gamma) \hat{\Delta}_i^2(\gamma) \hat{f}_i^{-1}(x_i, \gamma) \hat{f}(x_i) \leq cv_\alpha \right\},$$

where $\frac{1}{n} \sum_{i=1}^{n} k_h(q_i - \gamma) \hat{\Delta}_i^2(\gamma) \hat{f}_i^{-1}(x_i, \gamma) \hat{f}(x_i)$ replaces $\frac{\mathbb{E}[\Delta_f(x_i) \Delta_f(x_i) | q_i = \gamma_0]}{\mathbb{E}[\hat{\Delta}_f^2(x_i) f^2(x_i)(\sigma^2_f(x_i) + \sigma^2_f(x_i)|q_i = \gamma_0)]}$ in $LR_n^{(1)}(\gamma)$ by its estimate, and $cv_\alpha$ is the $(1 - \alpha)100\%$ quantile of $\chi^2_1$.

## 5 Inference Based on the IDKE with Shrinking Threshold Effects

In the previous section, the IDKE was adjusted by letting $k_1(0) = 0$ to construct a CI for $\gamma$ and the threshold effect was taken as fixed. In this section, the IDKE is adjusted from a different perspective by allowing for
the threshold effect to shrink to zero with the sample size but requiring that \( k_\delta(0) > 0 \).

### 5.1 Optimal Rate of Convergence for \( \gamma \)

First, we discuss the interpretation of a shrinking threshold effect. As argued in Section 2.4 of YP, the local shifter \( 1(q > \gamma) \) plays the role of an instrument. When \( q \) shifts from the left side of \( \gamma \) to its right side, the shift in the mean of \( y \) shrinks to zero. This behavior can be interpreted as the manifestation of a weak IV problem in the threshold regression context. A natural question that then arises is the identifiability of \( \gamma \) as \( \delta \) shrinks to zero. To put this question a different way, we can ask what is the minimum magnitude of \( \delta \) that ensures identification of \( \gamma \). For this purpose, we cast the model in the following framework.

Suppose \( \mathcal{P} \) is a family of probability models on some fixed measurable space \((\Omega, A)\). Let \( \gamma \) be a functional defined on \( \mathcal{P} \). Given an estimator \( \hat{\gamma} \) of \( \gamma \) and a loss function \( L(\hat{\gamma}, \gamma) \), the maximum expected loss over \( P \in \mathcal{P} \) is defined to be

\[
R(\hat{\gamma}, \mathcal{P}) = \sup_{P \in \mathcal{P}} \mathbb{E}_P [L(\hat{\gamma}, \gamma(P))],
\]

where \( \mathbb{E}_P \) is the expectation operator under the probability measure \( P \). A popular loss function (e.g., Stone (1980)) is the 0-1 loss

\[
L(\hat{\gamma}, \gamma) = 1\{ |\hat{\gamma} - \gamma| > \frac{\epsilon}{2} \}
\]

for some fixed \( \epsilon > 0 \), which will be used in this paper. Under this loss, \( R(\hat{\gamma}, \mathcal{P}) \) is the maximum probability that \( \hat{\gamma} \) is not in the \( \epsilon/2 \) neighborhood of \( \gamma \). The goal is to find an achievable lower bound for the minimax risk defined by

\[
\inf_{\hat{\gamma}} R(\hat{\gamma}, \mathcal{P}) = \inf_{\hat{\gamma}} \sup_{P \in \mathcal{P}} \mathbb{E}_P [L(\hat{\gamma}, \gamma(P))]. \tag{23}
\]

Only if \( \delta \) is large enough, will the right side converge to zero. The best rate of convergence of \( R(\hat{\gamma}, \mathcal{P}) \) to zero is then called the **optimal rate of convergence** or the **minimax rate of convergence**. Now \( P \in \mathcal{P} \) in our model is characterized by \( m_\pm(x, q) \) and \( \gamma \) as follows:

\[
\mathcal{P}(s, B) = \left\{ P_{m_\pm, \gamma} : \frac{dP_{m_\pm, \gamma}}{d\mu} = f(x, q) \varphi_{x, q} (y - m_-(x, q)1(q \leq \gamma) - m_+(x, q)1(q > \gamma)), \right. \\
m_-(x, q) \in C_s (B, \mathcal{X} \times \Gamma^-_s), m_+(x, q) \in C_s (B, \mathcal{X} \times \Gamma^+_s), \left. \int w \varphi_{x, q}(u) du = 0, \gamma \in \Gamma \right\},
\]

where \( \mu \) is Lebesgue measure on \( \mathbb{R}^d \), \( \varphi_{x, q}(u) \) is the conditional density of \( u \) given \((x', q)', \) and \( C_s (\cdot, \cdot) \) is defined in Section 2.2.

To formulate a precise statement of our next result, let

\[
\rho_n = \sqrt{\int_X (m_-(x, \gamma) - m_+(x, \gamma))^2 f(x|\gamma) dx},
\]

where \( \gamma = \gamma(P) \), and \( f(x|\gamma) \) can be replaced by any weight function \( w(x) \) with \( 0 < c \leq w(x) \leq C < \infty \) and \( \int_X w(x) = 1 \). If \( \mathbb{E}[\varepsilon | x, q] \) is continuous, then \( \rho_n = \left[ \mathbb{E}[X | q = \gamma] \right] \delta_n = O(\|\delta_n\|) \), similar to the \( \|\delta_n\| \) in Section 3.

**Theorem 5** Suppose Assumptions F, S and U hold, and \( P \in \mathcal{P}(s, B) \) with \( s \geq 1 \). If \( n^{\frac{\mu_s}{\delta}} \rho_n \to \infty \), then

\[
\lim_{n \to \infty} \inf_{\hat{\gamma}} \sup_{P \in \mathcal{P}(s, B)} P \left( n \rho_n^2 |\hat{\gamma} - \gamma(P)| > \frac{\epsilon}{2} \right) \geq C;
\]

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and if \( n^{2/z_{00}} \rho_n = O(1) \), then

\[
\lim_{n \to \infty} \inf_{\hat{\gamma}} \sup_{P \in P(x, B)} P \left( \left| \hat{\gamma} - \gamma(P) \right| > \frac{\epsilon}{2} \right) \geq C,
\]

for some positive constant \( C \) and small \( \epsilon > 0 \).

We begin our discussion of this result by clarifying a key difference between the parametric and nonparametric threshold models with shrinking threshold effects. In the former, as long as the jump size is \( n^{-\alpha} \) with \( 0 < \alpha < 1/2 \) (i.e., larger than \( n^{-1/2} \)), \( \gamma \) can be identified; in the latter, however, we require a jump size larger than \( n^{-2/z_{00}} \) to identify \( \gamma \). In other words, the minimum rate of convergence for \( \gamma \) in the nonparametric model must be larger than \( n^{-2/z_{00}} \) rather than any rate diverging to infinity as in the parametric model. In the parametric model, \( s = \infty \), so \( n^{-2/z_{00}} = n^{-1/2} \) and \( n^{-2/z_{00}} = 0 \), i.e., the parametric result is a limiting special case of Theorem 5 as \( s \to \infty \). Such a difference between the parametric model and the nonparametric model does not seem to have been explicitly recognized in the literature. For example, Müller and Song (1997) show that the convergence rate of the DKE is \( n\rho_n^2 \) when \( q \) is the only regressor by implicitly assuming a trade off in rates under which \( \rho_n \) is taken to be larger than \( n^{-2/z_{00}} \). In fact, when \( \gamma \) can be identified, the optimal rate of convergence for \( \gamma \) is the same as in the parametric case. This rate is achieved by the IDKE, as shown in the next section.

5.2 Asymptotics for \( \hat{\gamma} \)

To facilitate finding an expression for the limit distribution of \( \hat{\gamma} \), we define the following quantities

\[
D_n = \mathbb{E}[\Delta_f(x_i)\Delta_f(x_i)|q_i = \gamma_0]/\rho_n^2,
\]

\[
V_{1n} = \mathbb{E}[\Delta_f^2(x_i)f^2(x_i)\sigma_+^2(x_i)|q_i = \gamma_0]/\rho_n^2,
\]

\[
V_{2n} = \mathbb{E}[\Delta_f^2(x_i)f^2(x_i)\sigma_+^2(x_i)|q_i = \gamma_0]/\rho_n^2,
\]

where \( \rho_n \) is evaluated at \( \gamma_0 \). We also impose the following conditions on the bandwidth \( h \).

**Assumption H’:** \( h \to 0, \sqrt{n}h^d/\ln n \to \infty, \rho_n \to 0, \rho_n/h^s \to \infty, \sqrt{n}\rho_n^2 \to \infty \).

If we employ the optimal bandwidth \( h = O \left( n^{-2/z_{00}} \right) \), then \( \sqrt{n}h^d/\ln n = n^{-2/z_{00}d+1/2} \ln n \to \infty \) under Assumption G’ \((s \geq d)\). Also, \( \rho_n/h^s \to \infty \) and \( n\rho_n^2 \to \infty \) hold when \( n^{-2/z_{00}} \rho_n \to \infty \). The limit distribution of \( \hat{\gamma} \) is given in the following theorem.

**Theorem 6** Under Assumptions F, G’, H’, I, K’, S and U, if \( D_n \to D \) and \( V_{1n} \to V_1 \) as \( n \to \infty \), then

\[
n\rho_n^2(\hat{\gamma} - \gamma_0) \overset{d}{\to} \omega \cdot \Lambda(\lambda),
\]

where \( \omega = \frac{1}{\mathbb{E}^{1/V(\gamma_0)} D_1} \), and

\[
\Lambda(\lambda) = \arg\max_r \begin{cases} W_1(-r) - |r|, & \text{if } r \leq 0, \\ \sqrt{\lambda}W_2(r) - |r|, & \text{if } r > 0, \end{cases}
\]

with \( \lambda = V_2/V_1 \), and \( W_\ell(r), \ell = 1, 2 \), being two independent standard Wiener processes on \([0, \infty)\).

In some special cases, the asymptotic distribution of \( \hat{\gamma} \) can be simplified. For example, if \( \sigma_+^2(x_i) = \sigma_{10}^2 \) and \( \sigma_+^2(x_i) = \sigma_{20}^2 \), so that the model is locally homoskedastic within each regime, then \( \lambda = \sigma_{20}^2/\sigma_{10}^2 \); if
Figure 5: Comparison Between the PDFs of $N(0, 1)$ and $\Lambda(\lambda)$: $\lambda = 0.5, 1, 2$

$\sigma^2(x_i) = \sigma^2_+(x_i) = \sigma^2_0$, so that the model is homoskedastic locally around $\gamma_0$, then $\lambda = 1$. To compare with the asymptotic distribution of $\hat{\gamma}$ in Method II, note that by Slutsky’s theorem,

$$n\rho^2_n f_q(\gamma_0)D^2_n \left(\hat{\gamma} - \gamma_0\right) \xrightarrow{d} \Lambda(\lambda)$$

in Method III, whereas

$$\sqrt{n\rho^2_n f_q(\gamma_0)D^2_n \frac{k'_+(0)^2}{V_{1n}} (\hat{\gamma} - \gamma_0)} \xrightarrow{d} N(0, 1)$$

(24)

in Method II, so we use a different normalization on $(\hat{\gamma} - \gamma_0)$ to achieve a nondegenerate limit distribution. We can show that when $\rho_n \to 0$, the result in (24) still holds and the CI based on $LR^{(1)}(\gamma)$ remains valid. The convergence rate of $\hat{\gamma}$ in Method II is $\sqrt{n\rho^2_n f_q(\gamma_0)D^2_n \frac{k'_+(0)^2}{V_{1n}} (\hat{\gamma} - \gamma_0)} \xrightarrow{d} N(0, 1)$

with $(\hat{\gamma} - \gamma_0)$ being the cdf of $N(0, 1)$. Intuitively, when the heteroskedasticity measure $\lambda > 1$, it is more likely that $\Lambda(\lambda)$ achieves the maximum at $r > 0$. This intuition explains why the right tail of $\Lambda(\lambda)$ is heavier than the left tail. Interestingly, the effects of $\lambda$ on the limit distributions of the two $\gamma$ estimators are different: its effect on the estimator $\hat{\gamma}$ in Method II is to increase variance (but maintain symmetry), whereas its effect on the estimator $\hat{\gamma}$ in Method III is to introduce skewness.

For comparison, we state the limit distribution of the DKE in the following corollary. For this purpose,
we adjust Assumption H as follows.

**Assumption H**: \( h \to 0, \sqrt{n}h^d / \ln n \to \infty, \Delta_\alpha \to 0, \Delta_\alpha / h^s \to \infty, nh^d \Delta_\alpha^2 \to \infty, \) where \( \Delta_\alpha \) is defined in (21) and is equal to \((1, x'_\alpha, \gamma_0) \delta_n\) when \( E[\varepsilon|x, q] \) is continuous.

The optimal bandwidth \( h = O \left( n^{-\frac{2}{3d+2}} \right) \) satisfies Assumption H".

**Corollary 3** Under Assumptions F, G', H', I, K', S and U,

\[
nh^{d-1} \Delta_\alpha^2 (\hat{\gamma} - \gamma_0) \xrightarrow{d} \omega_\alpha \cdot \Lambda (\lambda_\alpha, \kappa),
\]

where \( \omega_\alpha = \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} \), and

\[
\Lambda (\lambda_\alpha, \kappa) = \arg\max_r \left\{ \frac{W_1(-r) - |r|}{\kappa}, \quad \text{if } r \leq 0, \right.
\]

\[
\left. \frac{\sqrt{\lambda_\alpha} W_2(r) - |r|}{\kappa}, \quad \text{if } r > 0, \right\}
\]

with \( \lambda_\alpha = \frac{\sigma^2(x_o)}{\sigma^2(x_o)} \), \( \kappa^2 \) being defined in (21), and standard Brownian motions \( W_\ell(r), \ell = 1, 2, \) as in Theorem 0.

The distribution of \( \Lambda (\lambda_\alpha, \kappa) \) is derived in Proposition 1 of Stryhn (1996). Since it will not be used for inference, it is omitted here. Compared with the convergence rate of \( \hat{\gamma} \) (viz., \( n \rho_n^2 \)), the convergence rate of \( \hat{\gamma} \) (viz., \( nh^{d-1} \Delta_\alpha^2 \)) is much slower especially when \( d \) is large. But it is still faster than the convergence rate of the DKE in Method II because the ratio \( nh^{d-1} \Delta_\alpha^2 / \sqrt{n}h^{d-2} \Delta_\alpha^2 = \sqrt{n}h^d \Delta_\alpha^2 \to \infty \). To compare with the limit distribution of \( \hat{\gamma} \) in Method II, note that by Slutsky’s theorem,

\[
\frac{nh^{d-1} \Delta_\alpha^2 f(x_o, \gamma_0)}{\sigma^2(x_o)} (\hat{\gamma} - \gamma_0) \xrightarrow{d} \Lambda (\lambda_\alpha, \kappa)
\]

in Method III, whereas

\[
\sqrt{\frac{nh^{d-1} \Delta_\alpha^2}{h \sigma^2(x_o)} \frac{k'_\alpha(0)^2}{(1 + \lambda_\alpha) \kappa^2 \xi_{(1)}} (\hat{\gamma} - \gamma_0)} \xrightarrow{d} N(0, 1)
\]
in Method II. The limit distributions of \( \hat{\gamma} \) in both methods involve only information local to \( x_o \). Similar to \( \lambda \) in the limit distributions of \( \hat{\gamma} \), \( \lambda_\alpha \) affects only the variance of \( \hat{\gamma} \) in Method II, but affects symmetry in Method III. A new factor \( \kappa^2 \) also appears in the limit distributions of \( \hat{\gamma} \); different from \( \lambda_\alpha \), the factor \( \kappa^2 \) increases variance but does not affect symmetry in either case.

It is also interesting to notice that the limit distribution of \( \hat{\gamma} \) in Method III does not depend on the kernel choice whereas the limit distribution of \( \hat{\gamma} \) in Method III does depend on the kernel choice on \( x \) (although not on \( q \)). These results echo Theorem 1 and Corollary 1 of YP where \( \rho_n \) is fixed and \( k_{+}(0) > 0 \). But when \( k_{+}(0) = 0 \), from Theorem 3 the asymptotic distribution of \( \hat{\gamma} \) depends on the kernel choice on \( q \), and from (21), the asymptotic distribution of \( \hat{\gamma} \) depends on the kernel choice on both \( x \) and \( q \). In other words, whether \( k_{+}(0) = 0 \) or not does indeed affect the role of the kernel on \( q \) with respect to data usage (and hence efficiency) of the estimators.

We next discuss inference concerning the threshold parameter \( \gamma \) based on our IDKE approach. Although we can construct a CI for \( \gamma \) by inverting the asymptotic distribution of \( \hat{\gamma} \) in Theorem 6, Hansen (2000) shows that such CIs perform poorly due to the identification failure when \( \rho_n = 0 \). He suggests constructing CIs for
γ by inverting the LR statistic instead, which in our case is defined as

\[ LR_n^{(2)}(γ) = nh \frac{1}{4k_+(0)} \frac{D_n}{V_{1n}} \left( \hat{Q}_n(\hat{γ}) - \hat{Q}_n(γ) \right). \]

To do so, we make use of the following result.

**Corollary 4** Under the assumptions of Theorem 6,

\[ LR_n^{(2)}(γ_0) \xrightarrow{d} M(λ), \]

where \( M(λ) \) follows the distribution \( P(λ(λ) \leq z) = (1 - e^{-z})(1 - e^{-z/λ}) \) with \( λ \) defined in Theorem 6.

To compare the LR statistic \( LR_n^{(2)}(γ) \) with \( LR_n^{(1)}(γ) \) in the last section, note that \( LR_n^{(1)}(γ) \) can be expressed as

\[ LR_n^{(1)}(γ) = nh \frac{k_+(0)}{ξ(1)} \frac{D_n}{V_{1n} + V_{2n}} \left( \hat{Q}_n(\hat{γ}) - \hat{Q}_n(γ) \right). \]

If \( \hat{Q}_n(\hat{γ}) - \hat{Q}_n(γ) \) are the same in these two LR statistics, then

\[ \frac{LR_n^{(1)}(γ)}{LR_n^{(2)}(γ)} = 4 \frac{k_+(0)k_+(0)}{ξ(1)} \frac{V_{1n}}{V_{1n} + V_{2n}} =: R_n. \]

If the model is locally homoskedastic in each regime, then \( R_n = 4 \frac{σ_{10}^2}{σ_{20}^2 + σ_{20}^2} \frac{k_+(0)k_+(0)}{ξ(1)} \), which is further simplified to \( 2k_+(0)k_+(0) \) when the model is locally homoskedastic in both regimes. However, \( \hat{Q}_n(\hat{γ}) - \hat{Q}_n(γ) \) are not the same in \( LR_n^{(1)}(γ) \) and \( LR_n^{(2)}(γ) \) because the employed kernels are different and the \( γ \)'s are different. To compare the asymptotic distributions of these two LR statistics, we plot the asymptotic pdfs in Figure 6 and report their 95% critical values in Table 2.
Figure 6: Comparison Between the PDFs of $\chi^2_1$ and $M$ for $\lambda = 0.5, 1, 2$

<table>
<thead>
<tr>
<th>Test Stat.</th>
<th>$LR_n^{(1)}$</th>
<th>$LR_n^{(2)}(\lambda = 0.5)$</th>
<th>$LR_n^{(2)}(\lambda = 1)$</th>
<th>$LR_n^{(2)}(\lambda = 2)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>95% crit</td>
<td>3.841</td>
<td>3.040</td>
<td>3.676</td>
<td>6.081</td>
</tr>
</tbody>
</table>

Table 2: 95% Critical Values of $LR_n^{(1)}$ and $LR_n^{(2)}$ for $\lambda = 0.5, 1, 2$

### 5.3 Comparison With the Parametric LSE

We close this section by comparing the IDKE with the LSE in the parametric case (see, e.g., Hansen (2000)). From YP, the LSE $\hat{\gamma}_{LSE}$ obtained by minimizing

$$
\min_{\beta, \delta} (Y - X\beta - X_{\leq \gamma}\delta)' (Y - X\beta - X_{\leq \gamma}\delta) = Y'Y - Y'P_{\gamma}Y
$$

is equivalent to estimation by maximizing

$$
\left( \hat{\delta}'X' \right) \left[ X (X'X)^{-1} X_{>\gamma} (X'X)^{-1} X_{\leq \gamma} X_{\leq \gamma} (X'X)^{-1} X' \right] \left( X \hat{\delta} \right),
$$

where $P_{\gamma}$ is the projection matrix onto $span (X, X_{\leq \gamma})$, $\hat{\delta} (\gamma)$ is the LSE of $\delta$ based on splitting according to the threshold $\gamma$, and $X$ and $X_{\leq \gamma}$ are defined in (2) with a similar definition for $X_{>\gamma}$. Here, $X \hat{\delta} (\gamma) = \left\{ x_i^T \hat{\delta} (\gamma) \right\}_{i=1}^n$, and $x_i^T \hat{\delta} (\gamma)$ is an estimator of the conditional mean differential $\Delta (x_i, \gamma)$. Since $\hat{\Delta}_i (\gamma)$ is estimating $\Delta (x_i, \gamma) f(x_i, \gamma)$, $X \hat{\delta} (\gamma)$ is mimicking $\left\{ \hat{\Delta}_i (\gamma) / f (x_i, \gamma) \right\}_{i=1}^n$. In the parametric case, $f(x_i, \gamma_0)$ and $f(x_i)$ do not appear in $D_n$, $V_{1n}$ and $V_{2n}$, so $\rho_n^2 D_n^2 / V_{1n}$ reduces to $\frac{\mathbb{E}[\Delta^2 | q = \gamma_0]}{\mathbb{E}[\Delta^2 | q = \gamma_0]}$ and $\lambda$ reduces to $\frac{\mathbb{E}[\Delta^2 | q = \gamma_0]}{\mathbb{E}[\Delta^2 | q = \gamma_0]}$; if $\sigma^2 (x_i) = \sigma^2_{10}$ and $\sigma^2 (x_i) = \sigma^2_{20}$, then $\rho_n^2 D_n^2 / V_{1n}$ further reduces to $\frac{\mathbb{E}[\Delta^2 | q = \gamma_0]}{\sigma^2_{10}}$ and $\lambda$ further reduces to $\frac{\sigma^2_{20}}{\sigma^2_{10}}$. A natural question is how to generate the same asymptotic distribution as in the parametric case when $\mathbb{E}[\varepsilon | x, q]$ is continuous. By careful inspection of the derivations in the proofs we can
show that, if the objective function of the IDKE changes to
\[
\frac{1}{n} \sum_{i=1}^{n} \hat{f}(x_i, \gamma) \left[ \frac{1}{n-1} \sum_{j=1, j \neq i}^{n} y_j K_{h,ij}^\gamma - \frac{1}{n-1} \sum_{j=1, j \neq i}^{n} K_{h,ij}^\gamma \right] f(x_i)^2,
\]
then the asymptotic distribution of the IDKE is the same as that of the parametric LSE, where \( \hat{f}(x_i, \gamma) \) and \( \hat{f}(x_i) \) are consistent estimators of \( f(x_i, \gamma) \) and \( f(x_i) \), respectively. Asymptotically, we impose a weight \( f_{q|x}(\gamma_{0}|x_{i}) \) on \( \Delta_2^2 \). This weight is intuitive in the sense that when there are more data points in the neighborhood of \( q = \gamma_0 \) at \( x_i \), we impose a larger weight on \( \Delta_2^2 \). In fact, we can also show that the IDKE using this objective function has the same asymptotic distribution as the LSE even in the framework of YP.\(^{33}\) In other words, the complex-looking weights in the square bracket of the objective function (27) are asymptotically equivalent to \( \{ f_{q|x}(\gamma_{0}|x_{i}) \}^{n}_{i=1} \). Such an equivalence result is not at all obvious from the original least squares objective function [26].

6 Two Specification Tests

In this section, we study limit theory of the two specification tests in Section 2.2. We first specify some regularity conditions which are modifications of Assumptions F, G and U given earlier.

Assumption F': \( f(x, q) \in C_1(B, \mathcal{X} \times \mathcal{Q}) \).

Assumption F'': \( f(x, q) \in C_\lambda(B, \mathcal{X} \times \Gamma_\varepsilon) \) with \( \lambda \geq 1 \), and \( 0 < f \leq f(x, q) \leq \overline{f} < \infty \) for \( (x', q') \in \mathcal{X} \times \Gamma_\varepsilon \).

Assumption G': \( i) \ g(x, q) \in C_\delta(B, \mathcal{X} \times \mathcal{Q}) \) with \( \delta \geq 2 \); \( ii) \ g(x, q) \in C_{\delta}(B, \mathcal{X} \times \Gamma_\varepsilon) \) with \( \delta \geq 2 \).

Assumption U':

(a) \( f(u|x, q) \) is continuous in \( u \) for \( (x', q') \in \mathcal{X} \times \mathcal{Q}^- \) and \( (x', q') \in \mathcal{X} \times \mathcal{Q}^+ \), where \( \mathcal{Q}^- = [q, \gamma_0] \) and \( \mathcal{Q}^+ = (\gamma_0, q] \).

(b) \( f(u|x, q) \) is Lipschitz in \( (x', q') \) for \( (x', q') \in \mathcal{X} \times \mathcal{Q}^- \) and \( (x', q') \in \mathcal{X} \times \mathcal{Q}^+ \).

(c) \( \mathbb{E}[u^4|x, q] \) is uniformly bounded on \( (x', q') \in \mathcal{X} \times \mathcal{Q} \).

The following Assumptions B1 and B2 are made on the bandwidths used in the first and second tests.

Assumption B1: \( nh^d \rightarrow \infty, h \rightarrow 0 \).

Assumption B2: \( nh^d \rightarrow \infty, b \rightarrow 0, nh^{d/2}b^{2\eta} \rightarrow 0 \), where \( \eta = \min(\lambda + 1, s) \).

Given \( d > 1 \), \( h/b \rightarrow 0 \) implies \( h^{d/2}/b \rightarrow 0 \), so \( nh^d \rightarrow \infty \) implies that \( nh^{d/2}b \rightarrow \infty \), where \( nh^{d/2}b \) is the magnitude of \( I_n^{(2)} \) under \( H_0^{(2)} \). The quantity \( nh^{d/2}b^{2\eta} \) is the bias of \( I_n^{(2)} \) under \( H_0^{(2)} \), so the assumption \( nh^{d/2}b^{2\eta} \rightarrow 0 \) guarantees that \( I_n^{(2)} \) is centered at the origin. Under \( H_0^{(1)} \), the bias of \( I_n^{(1)} \) is \( h^{d/2} \), so \( h \rightarrow 0 \) ensures that \( I_n^{(1)} \) is also centered at the origin. The condition \( h/b \rightarrow 0 \) requires that \( h \) is smaller than \( b \), which helps to generate power under \( H_1^{(2)} \) and shrink the bias under \( H_0^{(2)} \) to zero. Intuitively, if \( h/b \rightarrow 0 \), then the term \( K_{h,ij} \) in \( I_n^{(2)} \) makes the product \( \tilde{e}_i \tilde{e}_j \) behave like a squared term, producing an effect that generates power. In the first test, \( m(x, q) \) under \( H_0^{(1)} \) is parametric, so the corresponding bandwidth of \( b \) is a constant so that \( h \rightarrow 0 \) necessarily implies \( h/b \rightarrow 0 \). In testing \( H_0^{(2)} \) versus \( H_1^{(2)} \), our test statistic \( I_n^{(2)} \)

---

\(^{33}\) Using such an objective function, the asymptotic distribution of the IDKE with \( k_{\pm}(0) = 0 \) in Section 4 would also change. For example, \( \Sigma \) would change to \( \frac{\mathbb{E}[\Delta_2^2|\tilde{y}_i=\gamma_{0}]+\mathbb{E}[\Delta_2^2|\tilde{y}_i=\gamma_{0}^+]\xi(i)}{\mathbb{E}[\Delta_2^2|\tilde{y}_i=\gamma_{0}]+\mathbb{E}[\Delta_2^2|\tilde{y}_i=\gamma_{0}]}\left( \hat{Q}_n(\gamma) - \hat{Q}_n(\gamma) \right) \).
still applies when \( \mathbb{E}[\varepsilon|x,q] \) is not smooth at \( q = \gamma_0 \). But in this case the null hypothesis is better modified to the equivalence \( m_-(x) = m_+(x) \) for all \( x \in \mathcal{X} \) and \( y \) in Assumption \( G'' \) does not need to be smooth at \( q = \gamma_0 \). Also, we need to add the requirement \( nh^{d/2}b^3 \rightarrow 0 \) to Assumption \( B2 \), where \( nh^{d/2}b^3 \) is the bias of \( I_n^{(2)} \) attributed to the cusp of \( m(x,q) \) at \( q = \gamma_0 \).

In the second test, we impose the following assumption on the kernel \( l_b(\cdot,t) \).

**Assumption L**: \( l_b(\cdot,t) \) takes the form of (3) with order \( p = s + \lambda - 1 \).

Thus, \( l_b(\cdot,t) \) may be a higher order kernel to reduce the bias in \( \hat{f}_t \).

### 6.1 Limit Theory for the Two Tests

The following two theorems give the limit distribution of \( I_n^{(1)} \) under the null \( H_0^{(1)} \) and local power under \( H_1^{(1)} \). Note that the main component of \( I_n^{(1)} \) under \( H_1^{(1)} \) is a degenerate \( U \)-statistic, so the asymptotic distribution is normal rather than a functional of a chi-square process, as in the usual structural change literature.

**Theorem 7** Under Assumptions \( B1, F', G''(i), K, S, \) and \( U' \), the following hold:

(i) \[
I_n^{(1)} \overset{d}{\rightarrow} \mathcal{N}
\]

uniformly over \( \mathcal{H}_0^{(1)} \), where

\[
\Sigma^{(1)} = 2 \int k^{2d}(u)du\mathbb{E} \left[ f(x,q)\sigma^4(x,q) \right], \text{ with } \sigma^2(x,q) = \mathbb{E}[u^2|x,q],
\]

which can be consistently estimated by

\[
v_n^{(1)} = \frac{2h^d}{n(n-1)} \sum_i \sum_{j \neq i} K_h^{ij} \epsilon_i \epsilon_j^2.
\]

Hence, a test based on the studentized test statistic \( T_n^{(1)} = I_n^{(1)}/v_n^{(1)} \)

\[
T_n^{(1)} = 1 \left( T_n^{(1)} > z_{\alpha} \right),
\]

has significance level \( \alpha \), where \( z_{\alpha} \) is the \( 1-\alpha \) quantile of \( \mathcal{N}(0,1) \[34\]

(ii) If under \( H_1^{(1)} \), \( m(x,q) - \overline{m}(x,q) = n^{-1/2}h^{-d/4} \Delta_n(x,q) \) such that \( \int \Delta_n(x,q)^2 f(x,q)^2 dx dq \rightarrow \Delta \), then

\[
I_n^{(1)} \overset{d}{\rightarrow} \mathcal{N} \left( \Delta, \Sigma^{(1)} \right) \text{ and } T_n^{(1)} \overset{d}{\rightarrow} \mathcal{N} \left( \Delta/\sqrt{\Sigma^{(1)}}, 1 \right),
\]

so that the test \( T_n^{(1)} \) is consistent and \( P_m \left( T_n^{(1)} > z_{\alpha} \right) \rightarrow 1 \) for any \( m(\cdot) \) such that \( \int (m(x,q) - \overline{m}(x,q))^2 f(x,q)^2 dx dq \neq 0 \). Furthermore, the result continues to hold when \( z_{\alpha} \) is replaced by any nonstochastic constant \( C_n = o(nh^{d/2}) \).

According to this result, \( I_n^{(1)} \) has only trivial power if \( \mathbb{E} \left[ (m(x,q) - \overline{m}(x,q))^2 f(x,q) \right] = 0 \). Consider the following special example to illustrate. Suppose \( m(x,q) \) under \( H_0^{(1)} \) is \( x'\beta + x'\delta 1 \) \( \left( q \leq \gamma \right) \), and the alternative

\[34\] The test is one-sided because \( I_n^{(1)} \) is based on the \( L^2 \)-distance between \( m(\cdot) \) and \( \overline{m}(\cdot) \).
Theorem 8 Under Assumptions B2, F', G'' (ii), K, L, S, and U, the following hold:

(i) 
\[ I_n^{(2)} \overset{d}{\to} \mathcal{N}(0, \Sigma^{(2)}) \]
uniformly over \( \mathcal{H}_0^{(2)} \), where
\[ \Sigma^{(2)} = 2 \int k^2(u)du \mathbb{E} \left[ T_{\text{x}}^{\text{f}}(x, q) \sigma^4(x, q) \right], \]
which can be consistently estimated by
\[ v_n^{(2)} = \frac{2h^d}{n(n-1)} \sum_i \sum_{j \neq i} 1_i^1 1_j^1 K_{h,j}^2 e_i^2 e_j^2. \]
As a result, the test based on the studentized test statistic \( T_n^{(2)} = I_n^{(2)}/v_n^{(2)} \)
\[ t_n^{(2)} = 1 \left( T_n^{(2)} > z_\alpha \right), \]
has significance level \( \alpha \), where \( z_\alpha \) is the \( 1 - \alpha \) quantile of \( \mathcal{N}(0, 1) \).

(ii) If under \( H_0^{(2)} \), \( m_- - m_+ = n^{-1/2} h^{-d/4} b^{-1/2} \Delta_n(x) \) such that \( \int \Delta_n(x)^2 f(x, \gamma_0)^2 dx \to \Delta \), then
\[ I_n^{(2)} \overset{d}{\to} \mathcal{N} \left( \zeta \Delta, \Sigma^{(2)} \right) \text{ and } T_n^{(2)} \overset{d}{\to} \mathcal{N} \left( \zeta \Delta / \sqrt{\Sigma^{(2)}} , 1 \right), \]
where \( \zeta = 2 \int_0^1 \left( \int_0^1 t(u)du \right)^2 dv \), and the test \( t_n^{(2)} \) is consistent with \( P_m \left( T_n^{(2)} > z_\alpha \right) \to 1 \) for any \( m \) such that \( \int (m_-(x) - m_+(x))^2 f(x, \gamma_0)^2 dx \neq 0 \). The result continues to hold when \( z_\alpha \) is replaced by any nonstochastic constant \( C_n = o \left( nh^{d/2}b \right) \).

These two theorems show that \( I_n^{(1)} \) and \( I_n^{(2)} \) have power against different deviations of \( m(x, q) \) from \( H_0 \). For \( I_n^{(1)} \), power is generated from global deviations of \( m(x, q) \) from \( H_0 \), just as in classical specification testing (see, e.g., Theorem 3 of Zheng (1996) and Theorem 3.1 of Fan and Li (2000)). For \( I_n^{(2)} \), power is generated only from local deviations in the neighborhood of \( q = \gamma_0 \). In consequence, we need a larger deviation for \( I_n^{(2)} \) than for \( I_n^{(1)} \) to generate non-trivial power – specifically, \( n^{-1/2} h^{-d/4} b^{-1/2}/n^{-1/2} h^{-d/4} = b^{-1/2} \to \infty \).

6.2 Bootstrapping Critical Values

As is evident from the proofs of Theorems 7 and 8, the convergence rates of \( T_n^{(1)} \) and \( T_n^{(2)} \) to the standard normal is slow. The bias under \( H_0^{(1)} \) is \( h^{d/2} \) and under \( H_0^{(2)} \) is \( nh^{d/2} h^{d/2} \). Both these rates are low for some standard choices of bandwidth. As argued in the literature of classical specification testing (see, e.g., Härdle and Mammen (1993), Li and Wang (1998), Stute et al. (1998), Delgado and Manteiga (2001), and Gu et al.
In Step 1, a popular way to simulate is to use the two-point wild bootstrap residual $u_i^* = \tilde{e}_i (1 - \sqrt{5}) / 2$ with probability $(1 + \sqrt{5}) / (2\sqrt{5})$, and $u_i^* = \tilde{e}_i (1 + \sqrt{5}) / 2$ with probability $(\sqrt{5} - 1) / (2\sqrt{5})$, then $\mathbb{E}^*[u_i^*] = 0$, $\mathbb{E}^*[u_i^*2] = \tilde{c}_i^2$ and $\mathbb{E}^*[u_i^*3] = \tilde{c}_i^3$, where $\mathbb{E}^*[\cdot] = \mathbb{E}[\cdot | \mathcal{F}_n]$ and $\mathcal{F}_n = \{(x_i, q_i, y_i)\}_{i=1}^n$.

Step 2: Generate the bootstrap resample $\{y_i^*, x_i, q_i\}_{i=1}^n$ by

$$y_i^* = \tilde{g}_i + u_i^*.$$  

Then obtain the bootstrap residuals $\tilde{e}_i^* = y_i^* - \tilde{g}_i^*$, where $\tilde{g}_i^*$ is defined similarly to $\tilde{g}_i$ except that $y_i$ in the construction of $\tilde{g}_i$ is replaced by $y_i^*$.

Step 3: Use the bootstrap samples to compute the statistics

$$I_n^{(1)*} = \frac{nh^{d/2}}{n(n-1)} \sum_i \sum_{j \neq i} K_{h,i} \tilde{e}_i \tilde{e}_j,$$

$$I_n^{(2)*} = \frac{nh^{d/2}}{n(n-1)} \sum_i \sum_{j \neq i} \Gamma_i^{(1)} \Gamma_j^{(1)} K_{h,i,j} \tilde{e}_i \tilde{e}_j,$$

and the estimated asymptotic variances

$$v_n^{(1)*2} = \frac{2h^d}{n(n-1)} \sum_i \sum_{j \neq i} K_{h,i,j}^2 \tilde{e}_i^2 \tilde{e}_j^2,$$

$$v_n^{(2)*2} = \frac{2h^d}{n(n-1)} \sum_i \sum_{j \neq i} \Gamma_i^{(1)} \Gamma_j^{(1)} K_{h,i,j}^2 \tilde{e}_i \tilde{e}_j^2.$$

The studentized bootstrap statistics are

$$T_n^{(i)*} = I_n^{(i)*} / v_n^{(i)*}.$$  

Here, the same $b$ and $h$ are used as in $I_n^{(i)}$ and $v_n^{(i)}$ in Theorems 7 and 8.

Step 4: Repeat steps 1-3 $B$ times, and use the empirical distribution of $\{T_n^{(i)*}_{n,k}\}_{k=1}^B$ to approximate the null distribution of $T_n^{(i)}$. We reject $H_0^{(i)}$ if $T_n^{(i)} > T_n^{(i)*}_{n(\alpha B)}$, where $T_n^{(i)*}_{n(\alpha B)}$ is the upper $\alpha$-percentile of $\{T_n^{(i)*}_{n,k}\}_{k=1}^B$.

In Step 1, a popular way to simulate $u_i^*$ in the second test is based on $\tilde{e}_i$’s centralized counterpart $\bar{e}_i = \tilde{e}_i - \bar{e}$ rather than $\tilde{e}_i$ itself, where $\bar{e} = \sum_{i=1}^n \bar{e}_i \Gamma_i^{(1)} / \sum_{i=1}^n \Gamma_i^{(1)}$, $\Gamma_b = (\gamma - b, \gamma + b)$; see, e.g., Gijbels and Genestiaux (2004) and Su and Xiao (2008). Such a formulation can lead to $\sum_{i=1}^n \bar{e}_i \Gamma_i^{(1)} / \sum_{i=1}^n \Gamma_i^{(1)} = 0$, which will not

\footnote{To construct $I_n^{(2)*}$, we need only the data with $q_i \in [\gamma - b, \gamma + b]$.
\footnote{If we use a data-adaptive bandwidth such as cross-validation based on each bootstrap sample, then the algorithm is extremely time-consuming. See Chapter 3 of Mammen (1992) for related discussions.
\footnote{In the first test, $\frac{1}{n} \sum_{i=1}^n \tilde{e}_i = \frac{1}{n} \sum_{i=1}^n \tilde{e}_i 1(q_i \leq \tilde{\gamma}) + \frac{1}{n} \sum_{i=1}^n \tilde{e}_i 1(q_i > \tilde{\gamma}) = 0$ since the covariates include a constant term.}}
Theorem 9 Under the assumptions of Theorem 7 and 8, so these methods have appeal in applied work. As mentioned in the Introduction, good instruments are hard to find and justify in practice. Another reason for focusing on these two CI constructions and two specification tests is that neither involves the wild bootstrap procedure is not valid if the test statistic $m$ is discontinuous. Also, although Porter and Yu (2015) investigate the finite sample performance of a similar structural change test as $I_n^{(4)}$, no covariates $x$ are included in that work.

Theorem 9 Under the assumptions of Theorem 7 and 8

$$\sup_{z \in \mathbb{R}} \left| P \left( T_n^{(t)} \leq z \mid \mathcal{F}_n \right) - \Phi(z) \right| = o_p(1),$$

where $\Phi(\cdot)$ is the cdf of $\mathcal{N}(0, 1)$.

7 Simulations

We report simulations designed to assess the performance of our CIs and tests. We will concentrate on procedures whose performance is unclear in the literature. For inference, we will check only the CIs that invert the two LR statistics in Sections 4 and 5. It is unnecessary to compare the performance of the IDKE procedures whose performance is unclear in the literature. For inference, we will check only the CIs that invert the two LR statistics, we investigate only the two nonparametric tests developed in the main text because the performances of parametric tests developed in Supplement D are widely available in the literature.

We use a similar DGP as in YP for the simulation designs. Specifically, $y = \delta_1 1(q \leq \gamma) + \varepsilon,$ i.e., the threshold effect does not depend on $x$, where $\gamma = 0$ and $\Gamma = [-0.1, 0.1]$. $x$ and $q$ are independent and each is uniformly distributed over $[-0.5, 0.5]$, and $\varepsilon \mid (x, q) \sim N(-\delta_2 q^3, 0.1^2).$ In CI construction, we let $\delta_1 = 0.1$ and 0.2, indicating small and large threshold effects, respectively, and $\delta_2 = 1$, indicating severe endogeneity. In testing endogeneity, we let $\delta_1 = 0.2$ and $\delta_2 = 0, 0.2, 0.5$ and 1, where $\delta_2 = 0$ corresponds to the null. In testing threshold effects, $\delta_1 = 0, 0.1, 0.2$ and 0.5, and $\delta_2 = 1$, where $\delta_1 = 0$ corresponds to the null. For the IDKE with $k_\pm(0) = 0$,

$$k_-(x, r) = -x(1 + x)(-1 \leq x \leq r) \left\{ \frac{1}{6} - \frac{1}{2} r^2 - \frac{1}{3} r^3 \right\}, 0 \leq r \leq 1,$$
and for the IDKE with shrinking threshold effects,

\[ k_-(x, r) = \frac{3}{4} (1 - x^2) 1(-1 \leq x \leq r) \left( \frac{1}{2} + \frac{3}{4} r - \frac{1}{4} r^2 \right), 0 \leq r \leq 1, \tag{28} \]

which degenerates to the Epanechnikov kernel when \( r = 1 \); \( k_+(x, r) = k_-(x, r) \)\(^{40}\). In both tests, the kernel in \( K_{b,ij} \) is specified in (28), and in the second test, \( \hat{y}_i \) is estimated by the local linear smoother which implies a second-order boundary kernel in \( L_{b,ij} \) as required in Assumption L. Following DH, three bandwidths \( h \) are used based on the formula \( Cn^{-1/2} \) with proportionality constants \( C = 2, 3 \) and \( 4 \); in the second test, \( b = \frac{1}{2} h^{1/2} \) to guarantee \( h/b \to 0 \) and \( nh^{d/2}b^{2n} \to 0 \) with \( \eta = 2 \)\(^{41}\). The simulation study in Müller (1991) shows that a bandwidth without boundary adjustment works well, and we therefore use the same bandwidth for both interior and boundary points. \( N = 500 \) replications with sample size 500 and 1000 are used. In Algorithm WB, \( B = 399 \) when \( n = 500 \) and \( B = 199 \) when \( n = 1000 \). For CI construction the confidence level used is 95%, and for testing the level of significance is 5%.

### 7.1 Two LR-Based CIs

The coverage and average length of the two LR-based CIs are reported in Table 3. From Table 3, both methods perform well in coverage. Reductions in bandwidth and expansion of sample size both marginally improve coverage. On the other hand, different bandwidths and sample sizes have a big impact on CI lengths. Specifically, under our DGP, the medium bandwidth seems to perform satisfactorily for CI length among various scenarios and a larger sample size shrinks the length significantly. Another phenomenon deserving of mention is that CI length decreases sharply when the jump size doubles in both methods. This outcome is expected because larger threshold effects make the threshold point easier to identify. Comparing Method II to Method III, the latter behaves a little better in coverage. This may stem from the fact that the latter makes full use of the data information around \( \gamma_0 \) (\( k_+(0) > 0 \)) while the former makes only marginal use of such information (\( k_+(0) = 0 \)). This improvement comes at the cost that the CIs in Method III are generally longer than those in Method II.

### 7.2 Two Nonparametric Specification Tests

The size and power of the two nonparametric specification tests are reported in Tables 4 and 5, respectively. From these two tables, all tests have size close to the nominal 5% except in the second test with a large \( h \) where the test is undersized. A large \( h \) implies a large bias in \( \widehat{I}_n^{(2)} \) so the rejection probability is adversely affected. The power of the endogeneity test is very good - even when \( \delta_2 = 1 \) and \( n = 500 \), the power is 100%. The power of the second test is also very good - even when \( \delta_1 = 0.5 \) and \( n = 500 \), the power is close to 100%. As a benchmark, \( \delta_1 = 0.2 \) corresponds to two standard deviations of the error term \( u \) which follows \( \mathcal{N}(0, 0.1^2) \).

\(^{40}\)These kernel functions imply \( \xi_{(1)} = 12 \) and \( k_+(0) = 6 \) in Corollary 2 and \( k_+(0) = 1.5 \) in Corollary 4. So \( R_n = 2 \frac{k_+(0)k_+(0)}{\xi_{(1)}} = 1.5 \) in our DGP.

\(^{41}\)Notice that the range of bandwidths chosen is quite large, since the ratio of the proportionality constants between the first and the last is 2. In the estimation of \( \gamma \), the bandwidth is smaller than the usual optimal bandwidth (which is of rate \( n^{-\frac{1}{5+\varepsilon}} = n^{-\frac{1}{6}} \)), just as suggested in Porter and Yu (2015). In the second test, \( N = n \times (2 \times \frac{1}{4}C^{1/2}n^{-1/4})^2 = Cn^{1/2} \) data points are used to obtain \( \hat{y}_i \). When \( C = 2 \) and \( n = 500 \), \( N \approx 45 \). When \( C = 4 \) and \( n = 1000 \), \( N \approx 126 \).
Table 3: Comparison of Inferential Methods: Coverage and average length for nominal 95% confidence for $\gamma$ with $\delta_2 = 1$ and bandwidth proportionality constant $C$.

<table>
<thead>
<tr>
<th>$n$</th>
<th>500</th>
<th>1000</th>
<th>500</th>
<th>1000</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\delta_1$</td>
<td>0.1</td>
<td>0.2</td>
<td>0.1</td>
<td>0.2</td>
</tr>
<tr>
<td>$C = 2$</td>
<td>0.998</td>
<td>0.962</td>
<td>0.998</td>
<td>0.964</td>
</tr>
<tr>
<td>$C = 3$</td>
<td>0.994</td>
<td>0.958</td>
<td>1</td>
<td>0.964</td>
</tr>
<tr>
<td>$C = 4$</td>
<td>0.984</td>
<td>0.954</td>
<td>0.992</td>
<td>0.970</td>
</tr>
</tbody>
</table>

Table 4: Size and Power of $T_{n}^{(1)}$ (%): Nominal significance level 5%, $\delta_1 = 0.2$

<table>
<thead>
<tr>
<th>$n$</th>
<th>500</th>
<th>1000</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\delta_2$</td>
<td>0</td>
<td>0.2</td>
</tr>
<tr>
<td>$C = 2$</td>
<td>5.2</td>
<td>8</td>
</tr>
<tr>
<td>$C = 3$</td>
<td>5.8</td>
<td>10.2</td>
</tr>
<tr>
<td>$C = 4$</td>
<td>5.2</td>
<td>10.8</td>
</tr>
</tbody>
</table>

Table 5: Size and Power of $T_{n}^{(2)}$ (%): Nominal significance level 5%, $\delta_2 = 1$

<table>
<thead>
<tr>
<th>$n$</th>
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</tr>
</thead>
<tbody>
<tr>
<td>$\delta_1$</td>
<td>0</td>
<td>0.1</td>
</tr>
<tr>
<td>$C = 2$</td>
<td>4.4</td>
<td>17.4</td>
</tr>
<tr>
<td>$C = 3$</td>
<td>4.0</td>
<td>22.2</td>
</tr>
<tr>
<td>$C = 4$</td>
<td>3.8</td>
<td>17.2</td>
</tr>
</tbody>
</table>

8 Conclusion

All three methods of estimation presented here for threshold point regression remain valid and invariant to endogeneity of the threshold variable. To the best of our knowledge these methods are the only ones in the literature offering such robustness. The first method is a nonlinear 2SLS method and requires instruments, while the other two methods are based on smoothing the objective function of the IDKE and do not require any instrumentation in their implementation. These are important advantages in empirical work where valid instruments are often scarce.

Our development and discussion of the 2SLS method clarifies some puzzles in the current literature about the properties of threshold regression estimation. We draw attention in particular to the following matters that are resolved in the paper: (i) why the usual GMM method cannot identify the threshold point in structural change models; (ii) why two groups of moments are required to identify the threshold point when the threshold variable is exogenous and correlated with the covariates and instruments, whereas only one
group of moments is sufficient when the threshold variable is endogenous; and (iii) why the bootstrap is valid for our 2SLS method while it is generally invalid for the usual GMM approach.

In discussing the two IDKE-smoothing methods, we show that these IDKEs use different normalizations to obtain operable asymptotic distributions under different assumptions and we explain why their convergence rates are different. We further show how to construct confidence intervals by inverting the LR statistics in both methods. Our three inferential methods provide considerable flexibility to practitioners. When instruments are available, the 2SLS method of estimation can be used, coupled with use of the bootstrap for inference. When instruments are absent, the other two methods can be used.

Two specification tests are suggested, one designed to check for the presence of endogeneity and the other to check for threshold effects. Our results show that it is possible to test for threshold effects in the absence of instrumentation even if endogeneity is present. An important implication of the test for endogeneity in empirical work is that it helps to assess whether instruments are required to achieve consistent estimation of the structural coefficients. Both tests are similar to score tests and have convenient asymptotic normal distributions, although a wild bootstrap procedure is suggested to determine critical values for improved finite sample performance.
References


Hansen, B.E., 1996, Inference when a Nuisance Parameter is not Identified under the Null Hypothesis, *Econometrica*, 413-430.


Yu, P., 2013b, Integrated Quantile Threshold Regression and Distributional Threshold Effects, mimeo.


Yu, P. and X. Fan, 2019, Threshold Regression With A Threshold Boundary, mimeo.


We first collect notation for future reference. The $n \times 1$ vectors $Y$ and $\varepsilon$ stack the variables $y_i$ and $\varepsilon_i$, the $n \times d_x$ matrices $X$, $X_{\leq \gamma}$ and $X_{> \gamma}$ stack the vectors $x_i'$, $x_i'1(q_i \leq \gamma)$ and $x_i'1(q_i > \gamma)$, and the $n \times d_z$ matrices $Z$, $Z_{\leq \gamma}$ and $Z_{> \gamma}$ are similarly defined. The symbol $\approx$ means asymptotic equivalence in the sense that higher order terms are neglected, $\equiv$ signifies equality in distribution, $C$ means a positive constant that may change at each occurrence, and $\Rightarrow$ signifies weak convergence of the respective probability measures over an associated compact metric space.

To aid intuition in the development of Methods II and III, we let $(x;q) = x_0$ throughout the proofs for these two methods.

**Supplement A: Proofs**

**Proof of Theorem 1** We first show the consistency of $\hat{\beta}$. If $\hat{\beta}$ is consistent, then a standard argument can be applied to show that

$$\left( \hat{\beta}, \hat{\delta} \right) = \left( \hat{\beta}(\hat{\gamma}), \hat{\delta}(\hat{\gamma}) \right)$$

with

$$\left( \hat{\beta}(\gamma), \hat{\delta}(\gamma) \right) = \left[ (X'_{\leq \gamma}Z) \hat{W}(Z'X_{\leq \gamma}) \right]^{-1} \left[ (X'_{\leq \gamma}Z) \hat{W}(Z'Y) \right]$$

is consistent, where $X_{\leq \gamma} = (X, X_{< \gamma})$. So we now concentrate on the consistency of $\hat{\gamma}$. First note that the concentrated objective function of (12) after plugging in $\hat{\beta}(\gamma), \hat{\delta}(\gamma)$ is

$$\hat{Q}_n(\gamma) = \hat{g}_n(\gamma)' \hat{W} \hat{g}_n(\gamma).$$

Here,

$$\hat{g}_n(\gamma) = \frac{1}{n} Z' \left( Y - X_{\leq \gamma} \left[ \left( \frac{1}{n} X'_{\leq \gamma} Z \right) \hat{W} \left( \frac{1}{n} Z'X_{\leq \gamma} \right) \right]^{-1} \left[ \left( \frac{1}{n} X'_{\leq \gamma} Z \right) \hat{W} \left( \frac{1}{n} Z'Y \right) \right] \right)$$

$$= \left\{ I - \left( \frac{1}{n} Z'X_{\leq \gamma} \right) \left[ \left( \frac{1}{n} X'_{\leq \gamma} Z \right) \hat{W} \left( \frac{1}{n} Z'X_{\leq \gamma} \right) \right]^{-1} \left( \frac{1}{n} Z'X_{\leq \gamma} \theta_0 + \frac{1}{n} Z'\varepsilon \right) \right\}$$

$$= \left\{ I - \left( \frac{1}{n} Z'X_{\leq \gamma} \right) \left[ \left( \frac{1}{n} X'_{\leq \gamma} Z \right) \hat{W} \left( \frac{1}{n} Z'X_{\leq \gamma} \right) \right]^{-1} \left( \frac{1}{n} Z'X_{\leq \gamma} \delta_n + \frac{1}{n} Z'\varepsilon \right) \right\},$$

and the second equality holds because the first $d$ columns of $Z'X_{\leq \gamma}$ are the same as those of $Z'X_{\leq \gamma}$. We apply Theorem 2.1 of Newey and McFadden (1994) to prove the result. First, we can show

$$\sup_{\gamma \in \Gamma} \left\| \frac{1}{n} \hat{g}_n(\gamma) - g_0(\gamma) \right\| \rightarrow 0,$$

Footnotes:

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where
\[ g_0(\gamma) = \left( I - G_2 \left( G_p'WG_p \right)^{-1} G_p'W \right) G_{2,\gamma_0} c. \]

To see why, note that by a Glivenko-Cantelli argument,
\[ I_i - \left( \frac{1}{n} Z'X_{\leq \gamma} \right) \left( \frac{1}{n} X'X_{\leq \gamma} \right)^{-1} \left( \frac{1}{n} X'Z \right) \overset{p}{\to} \left( I - G_2 \right) \left( G_p'WG_p \right)^{-1} G_p'W \]
uniformly in \( \gamma \). Also, \( \frac{1}{n} Z'X_{\leq \gamma} \frac{4n}{\| \delta_n \|} \overset{p}{\to} G_{2,\gamma_0} c \) and \( \frac{1}{n} Z' \varepsilon = O_p \left( \frac{1}{\sqrt{n} \| \delta_n \|} \right) = o_p (1) \), the result follows.

Second, by the CMT,
\[
\hat{Q}_n (\gamma) / \| \delta_n \|^2 \overset{p}{\to} c' G_{2,\gamma_0} \left( I - W G_2 \left( G_p'WG_p \right)^{-1} G_p' \right) W \left( I - G_2 \left( G_p'WG_p \right)^{-1} G_p' \right) G_{2,\gamma_0} c = c' R_{2,\gamma_0} \left( I - R_c \left( R_c^{-1} R_c^{-1} \right) R_c \right) R_{2,\gamma_0} c = \| R_{2,\gamma_0} c \|^2 =: Q_0 (\gamma)
\]
uniformly in \( \gamma \), where \( R_c = W_{1/2}G_2 \) and \( R_{2,\gamma} = W_{1/2}G_{2,\gamma} \). Obviously, \( Q_0 (\gamma_0) = 0 \). Also, \( \| R_{\gamma} \| \) is a projection on a \( 2d \)-dimensional space, while \( R_{2,\gamma_0} \) is a \( t(2d) \)-dimensional vector, so as long as \( R_{n_0} c \) does not fall in \( \text{span} \left( R_{\gamma} \right) \) when \( \gamma \neq \gamma_0 \), \( Q_0 (\gamma) > 0 \). This requirement is satisfied by virtue of Assumption IV.

We can now adjust Theorem 7.2 of Newey and McFadden (1994) to derive the asymptotic distribution of \( \hat{\theta} \). We only point out the difference in the proof. Replace \( G \) by \( G_n r_n \) and \( \theta - \theta_0 \) by \( r_n^{-1} (\theta - \theta_0) \), where \( G_n = - \left( E \left[ z' x \right], E \left[ z' x' \left| \gamma = \gamma_0 \right. \right] f_q (\gamma_0) \right) \) and \( r_n = \left( \begin{array}{c} I_{2d} \ 0 \\ 0 \ 1/ \| \delta_n \| \end{array} \right) \). Then \( H = -r_n G_n' W G_n r_n \rightarrow -G' W G \) and \( \bar{D} = -r_n G_n' \hat{W} \sqrt{n} \hat{g}_n (\theta_0) \overset{d}{\to} N (0, G' W G I) \). What remains is to show that for any \( h_n \rightarrow 0 \),
\[
\sup_{\| \theta - \theta_0 \| \leq h_n} \sqrt{n} \| \hat{g}_n (\theta) - \bar{g}_n (\theta_0) - g_0 (\theta_0) \| \overset{p}{\to} 0.
\]
This stochastic equicontinuity result is obvious because \( \hat{g}_n (\theta) \) is generated by a VC subgraph class of functions. We mention that this part of proof is similar to the ‘convergence rate and asymptotic normality’ part in the proof of Theorem 1 of Seo and Shin (2016). Their consistency proof is marred by a typo, which has a material effect. Specifically, at the end of page 181, the probability limit should be \( I - A(\gamma) \) • • • rather than \( I + A(\gamma) \) • • •. This is why they did not specify an identification assumption such as that in Assumption IV.

\[ \text{Proof of Theorem 2} \]
First assume \( \hat{W} = \text{diag} \left\{ \hat{W}_1, \hat{W}_2 \right\} \) with \( \hat{W}_1 \overset{p}{\to} W_1 > 0 \) and \( \hat{W}_2 \overset{p}{\to} W_2 > 0 \). Then \( \hat{Q}_n (\theta) \) can be expressed as a sum of two quadratic forms:
\[
\hat{Q}_n (\theta) = \hat{Q}_{1n} (\theta_1) + \hat{Q}_{2n} (\theta_2),
\]
where
\[
\hat{Q}_{1n} (\theta_1) = \hat{m}_{1n} (\theta_1)' \hat{W}_1 \hat{m}_{1n} (\theta_1) \quad \text{and} \quad \hat{Q}_{1n} (\theta_1) = \hat{m}_{2n} (\theta_2)' \hat{W}_2 \hat{m}_{2n} (\theta_2),
\]
\[
\hat{m}_{1n} (\theta_1) = \frac{1}{n} \sum_{i=1}^{n} z_i (y_i - x_i' \beta_1) 1 (q_i \leq \gamma),
\]
\[
\hat{m}_{2n} (\theta_2) = \frac{1}{n} \sum_{i=1}^{n} z_i (y_i - x_i' \beta_2) 1 (q_i > \gamma),
\]
with $\theta_1 = (\beta_1', \gamma')'$ and $\theta_2 = (\beta_2', \gamma')'$. It is not hard to see that
\[
\left( \begin{array}{c}
\bar{\beta}_1(\gamma) \\
\bar{\beta}_2(\gamma)
\end{array} \right) = \left( \begin{array}{c}
[X_{\leq, Z} W_1 (Z^* X_{\leq})]^{-1} \left( (X_{\leq, Z} W_1 (Z Y_{\leq})) 
\right) \\
[X_{\leq, Z} W_2 (Z^* X_{>})]^{-1} \left( (X_{\leq, Z} W_2 (Z Y_{>})) 
\right)
\end{array} \right),
\]
so if $\bar{\gamma}$ is consistent, then both $\bar{\beta}_1$ and $\bar{\beta}_2$ are consistent. So we concentrate on the consistency of $\bar{\gamma}$.

By similar analysis to the proof of Theorem $\Box$
\[
\frac{Q_n(\gamma)}{P} = \| P_{\bar{r}_1, \gamma}(\bar{r}_1, \gamma, \theta_0) \|^2 - \| P_{\bar{r}_2, \gamma}(\bar{r}_2, \gamma, \theta_0) \|^2 - \| P_{\bar{r}_1, \gamma}(\bar{r}_1, \gamma, \theta_0) \|^2,
\]
where $\bar{r}_1, \gamma, \theta_0 = W_1^{1/2} \mathbb{E}[z_{Y_{\leq}}]$, $\bar{r}_2, \gamma, \theta_0 = W_2^{1/2} \mathbb{E}[z_{Y_{>}}]$, $\bar{R}_1, \gamma = W_1^{1/2} \mathbb{E}[z_{X_{\leq}}]$ and $\bar{R}_2, \gamma = W_2^{1/2} \mathbb{E}[z_{X_{>}}]$. Note that
\[
\left( \begin{array}{c}
\mathbb{E}[z_{Y_{\leq}}] \\
\mathbb{E}[z_{Y_{>}}]
\end{array} \right) = \left( \begin{array}{c}
\mathbb{E}[z_{X_{\leq}}] \\
\mathbb{E}[z_{X_{>}}]
\end{array} \right) \left( \begin{array}{cc}
\beta_1 & 0 \\
0 & \beta_2
\end{array} \right) + \left( \begin{array}{c}
\mathbb{E}[z_{\leq}] \\
\mathbb{E}[z_{>}] \end{array} \right),
\]
if $\gamma \leq \gamma_0$. Then we can find $a_1$ and $a_2$ such that $\mathbb{E}[z_{X_{\leq}}] a_1 = \mathbb{E}[z_{X_{>}}] \beta_1 + \mathbb{E}[z_{X_{>}}] \beta_2$.

We consider the following two cases.

(i) $q$ is exogenous (i.e., $q$ is included in $z$, and $\mathbb{E}[z|z] = 0$). In this case, $\mathbb{E}[z_{\leq}] = \mathbb{E}[z_{>}] = 0$. First suppose $\gamma \leq \gamma_0$. Then $\bar{r}_1, \gamma, \theta_0 = W_1^{1/2} \mathbb{E}[z_{X_{\leq}}] \beta_1$ and $\bar{r}_2, \gamma, \theta_0 = W_2^{1/2} \mathbb{E}[z_{X_{>}}] \beta_2$. The question is whether we can find $a_1$ and $a_2$ such that $\mathbb{E}[z_{X_{\leq}}] a_1 = \mathbb{E}[z_{X_{>}}] \beta_1 + \mathbb{E}[z_{X_{>}}] \beta_2$.

(ii) $q$ is endogenous and satisfies only $\mathbb{E}[z_{1}(q < \gamma_0)] = 0$ and $\mathbb{E}[z|z_{1}(q > \gamma_0)] = 0$. Again, first suppose $\gamma \leq \gamma_0$. Then we try to select $a_1$ and $a_2$ such that $\mathbb{E}[z_{X_{\leq}}] a_1 = \mathbb{E}[z_{X_{>}}] \beta_1 + \mathbb{E}[z_{X_{>}}] \beta_2$, and $\mathbb{E}[z_{X_{\leq}}] a_2 = \mathbb{E}[z_{X_{>}}] \beta_1 + \mathbb{E}[z_{X_{>}}] \beta_2 + \mathbb{E}[z_{X_{>}}] \beta_2 + \mathbb{E}[z_{X_{>}}] \beta_2 + \mathbb{E}[z_{X_{>}}] \beta_1$, and $\mathbb{E}[z_{X_{>}}] (a_1 - a_2) + \mathbb{E} [z_{X_{>}}] = 0$. Such an $a_1$ and $a_2$ are impossible by Assumption IV(iii). Next suppose $\gamma > \gamma_0$. Then we try to select $a_1$ and $a_2$ such that $\mathbb{E}[z_{X_{\leq}}] a_1 = \mathbb{E}[z_{X_{>}}] \beta_1 + \mathbb{E}[z_{X_{>}}] \beta_1 + \mathbb{E}[z_{X_{>}}] \beta_1 + \mathbb{E}[z_{X_{>}}] \beta_1$, and $\mathbb{E}[z_{X_{>}}] (a_1 - a_2) + \mathbb{E} [z_{X_{>}}] = 0$. Such an $a_1$ and $a_2$ are impossible by Assumption IV(iii).

Now consider the case where $q$ is exogenous and is independent of $(z', x')$. Suppose $\gamma \leq \gamma_0$. We can set $a_2 = (1 - F_k(\gamma_0))(\beta_2 - \beta_1) + \beta_1 = (1 - F_k(\gamma_0))(\beta_2 + F_k(\gamma_0) - F_k(\gamma_0))\beta_1$, which is $\theta_2^{(3)}(\gamma)$ in Proposition 1(ii) of HBB where $q \sim U[0, 1]$. Now suppose $\gamma > \gamma_0$. We can set $a_1 = F_k(\gamma) \beta_1 + (F_k(\gamma) - F_k(\gamma_0)) \beta_2$, which is $\theta_1^{(1)}(\gamma)$ in Proposition 1(ii) of HBB where $q \sim U[0, 1]$.3}
If $W$ is a general positive definite matrix, then
\[ \bar{Q}_n(\gamma) \xrightarrow{p} \| \tilde{R}_{\gamma,\theta_0} \|^2 - \| P_{\tilde{R}_{\gamma,\theta_0}} \|^2, \]
where $\tilde{R}_{\gamma} = W^{1/2} \bar{G}_{\gamma}$ with $\bar{G}_{\gamma} = \begin{pmatrix} E [z \bar{x}_{\leq 1}] & 0 \\ E [z \bar{x}_{> 1}] & 0 \end{pmatrix}$, and $\tilde{R}_{\gamma,\theta_0} = W^{1/2} \begin{pmatrix} E [z y_{\leq \gamma}] \\ E [z y_{> \gamma}] \end{pmatrix}$. In case (i), if there does not exist $a_1$ and $a_2$ such that
\[ \left( \begin{array}{c} E [z \bar{x}_{\leq 1}] \\ E [z \bar{x}_{> 1}] \end{array} \right) a_1 + \left( \begin{array}{c} 0 \\ E [z \bar{x}_{\leq 1}] \\ E [z \bar{x}_{> 1}] \end{array} \right) a_2 = \left( \begin{array}{c} E [z \bar{x}_{\leq 1}] \\ E [z \bar{x}_{> 1}] \end{array} \right), \]
or $E [z \bar{x}_{\leq 1}] a_1 = E [z \bar{x}_{\leq 1}] \beta_{10}$ and $E [z \bar{x}_{> 1}] a_2 = E [z \bar{x}_{> 1}] \beta_{10} + E [z \bar{x}_{> 1}] \beta_{20}$ for any $\gamma < \gamma_0$, and
\[ a_1 + \left( \begin{array}{c} 0 \\ E [z \bar{x}_{\leq 1}] \\ E [z \bar{x}_{> 1}] \end{array} \right) a_2 = \left( \begin{array}{c} E [z \bar{x}_{\leq 1}] \\ E [z \bar{x}_{> 1}] \end{array} \right), \]
or $E [z \bar{x}_{\leq 1}] a_1 = E [z \bar{x}_{\leq 1}] \beta_{10} + E [z \bar{x}_{> 1}] \beta_{20}$ and $E [z \bar{x}_{> 1}] a_2 = E [z \bar{x}_{> 1}] \beta_{20}$ for any $\gamma > \gamma_0$, then $\gamma_0$ is identified. These conditions are exactly the same as in the diagonal $W$ case. Similarly, in case (ii), the same conditions as in the diagonal $W$ case are required to identify $\gamma_0$.

**Proof of Theorem 3.** Proposition 1 proves the consistency of $\hat{\gamma}$, and Proposition 2 proves $\hat{\gamma} - \gamma_0 = O_p((n/h)^{-1/2})$, so we can apply the argmax continuous mapping theorem (see, e.g., Theorem 3.2.2 of Van der Vaart and Wellner (1996)) to establish the asymptotic distribution of $\sqrt{n/h}(\hat{\gamma} - \gamma_0)$. From Proposition 3, the finite-dimensional limit distributions of $nh \left( \bar{Q}_n(\gamma_0) - \bar{Q}_n(\gamma_0) \right)$ are the same as those of
\[ -v^T E [\Delta_f(x,\gamma) f(x)] q_i = \gamma_0 \gamma_0 f_\theta(\gamma_0) k_4(0) + 2Uv, \]
where $U \sim N(0, E \Delta^2_f(x) f^2(x) (\sigma^2_f(x) + \sigma^2(x))) q_i = \gamma_0 f_\theta(\gamma_0) k_4(0)$. Combining this with the stochastic equicontinuity result in Lemma 4, we have
\[ nh \left( \bar{Q}_n(\gamma_0) - \bar{Q}_n(\gamma_0) \right) \rightsquigarrow -v^T E [\Delta_f(x) \Delta_f(x)] q_i = \gamma_0 f_\theta(\gamma_0) k_4(0) + 2Uv. \]

Thus,
\[ \sqrt{n/h}(\hat{\gamma} - \gamma_0) \xrightarrow{d} v^* = \arg \max_v \left\{ -v^T E [\Delta_f(x) \Delta_f(x)] q_i = \gamma_0 f_\theta(\gamma_0) k_4(0) + 2Uv \right\} = \frac{U}{E \Delta_f(x) \Delta_f(x) [q_i = \gamma_0 f_\theta(\gamma_0) k_4(0)]} \sim N(0, \Sigma). \]

**Proof of Theorem 4.** We prove the theorem in two steps. First, we show that $\tilde{\Sigma}_0 \rightarrow \Sigma$, where $\tilde{\Sigma}_0$ is replacing $\hat{\Sigma}$ in $\tilde{\Sigma}$ by $\gamma_0$. Second, we show that $\tilde{\Sigma} - \tilde{\Sigma}_0 \xrightarrow{p} 0$. The first result is shown in Proposition 4 and the second is shown in Proposition 5.

**Proof of Corollary 2.** From the proof of Theorem 3 and the CMT,
\[ nh \left( \begin{array}{c} E [\Delta_f(x) \Delta_f(x)] q_i = \gamma_0 k_4(0) \\ E [\Delta_f(x) \Delta_f(x)] q_i = \gamma_0 k_4(0) \end{array} \right) \xrightarrow{d} \left( \begin{array}{c} E [\Delta_f(x) \Delta_f(x)] q_i = \gamma_0 k_4(0) \\ E [\Delta_f(x) \Delta_f(x)] q_i = \gamma_0 k_4(0) \end{array} \right) \]
\[ = \frac{U^2}{E [\Delta_f(x) \Delta_f(x)] q_i = \gamma_0 k_4(0)} \xrightarrow{d} \chi^2. \]

By the proof of Theorem 4 and Slutsky’s theorem, $\tilde{\Sigma}_0 \xrightarrow{d} \chi^2$. ■

**Proof of Theorem 5.** Assume the densities of $(x', q)'$ and $e$ are known. Since the minimax risk for a
larger class of probability models is no smaller than that for a smaller class of probability models, the lower bound for a particular distributional assumption also holds for a wider class of distributions. To simplify the calculation, assume \(e_i\) is iid \(N(0, 1)\) and \((x'_i, q_i)'\) is iid uniform on \(X \times [0, 1]\). Such a specification also appears in Fan (1993) where it is called the assumption of richness of joint densities. We will use the technique in Sun (2005) to develop our results. This technique is also implicitly used in Stone (1980) and the essential part of the technique can be cast in the language of Neyman-Pearson testing.

Let \(P, Q\) be probability measures defined on the same measurable space \((\Omega, \mathcal{A})\) with the affinity between the two measures defined as usual to be

\[
\pi(P, Q) = \inf \left( \mathbb{E}_P[\phi] + \mathbb{E}_Q[1 - \phi] \right),
\]

where the infimum is taken over the measurable function \(\phi\) such that \(0 \leq \phi \leq 1\). In other words, \(\pi(P, Q)\) is the smallest sum of type I and type II errors of any test between \(P\) and \(Q\). It is a natural measure of the difficulty of distinguishing \(P\) and \(Q\). Suppose \(\mu\) is a measure dominating both \(P\) and \(Q\) with corresponding densities \(p\) and \(q\). It follows from the Neyman-Pearson lemma that the infimum is achieved by setting \(\phi = 1(p \leq q)\) and then

\[
\pi(P, Q) = \int 1(p \leq q) p d\mu + \int 1(p > q) q d\mu = 1 - \frac{1}{2} \int |p - q| d\mu \geq 1 - \frac{1}{2} \|P - Q\|_1
\]

where \(\|\cdot\|_1\) is the \(L_1\) distance between two probability measures. Now consider a pair of probability models \(P, Q \in \mathcal{P}(s, B)\) such that \(|\gamma(P) - \gamma(Q)| \geq \epsilon\).

For any estimator \(\hat{\gamma}\), we have

\[
1 \left( \|\hat{\gamma} - \gamma(P)\| > \epsilon/2 \right) + 1 \left( \|\hat{\gamma} - \gamma(Q)\| > \epsilon/2 \right) \geq 1.
\]

Let

\[
\phi = \frac{1(|\hat{\gamma} - \gamma(P)| > \epsilon/2)}{1(|\hat{\gamma} - \gamma(P)| > \epsilon/2) + 1(|\hat{\gamma} - \gamma(Q)| > \epsilon/2)}.
\]

Then \(0 \leq \phi \leq 1\) and

\[
\sup_{P \in \mathcal{P}(s, B)} P \left( |\hat{\gamma} - \gamma(P)| > \epsilon/2 \right) \geq \frac{1}{2} \{ P \left( |\hat{\gamma} - \gamma(P)| > \epsilon/2 \right) + Q \left( |\hat{\gamma} - \gamma(Q)| > \epsilon/2 \right) \} \geq \frac{1}{2} \mathbb{E}_P[\phi] + \frac{1}{2} \mathbb{E}_Q[1 - \phi].
\]

Therefore

\[
\inf \sup_{\hat{\gamma}} P \left( |\hat{\gamma} - \gamma(P)| > \epsilon/2 \right) \geq \frac{1}{2} \pi(P, Q)
\]

for any \(P\) and \(Q\) such that \(|\gamma(P) - \gamma(Q)| \geq \epsilon\). So we need only search for the pair \((P, Q)\) which minimize \(\pi(P, Q)\) subject to the constraint \(|\gamma(P) - \gamma(Q)| \geq \epsilon\). To obtain a lower bound with a sequence of independent observations, let \((\Omega, \mathcal{A})\) be the product space and \(\mathcal{P}(s, B)\) be the family of product probabilities on such a space. Then for any pair of finite-product measures \(P = \prod_{i=1}^n P_i\) and \(Q = \prod_{i=1}^n Q_i\), the minimax risk satisfies

\[
\inf \sup_{P \in \mathcal{P}(s, B)} P \left( |\hat{\gamma} - \gamma(P)| > \epsilon/2 \right) \geq \frac{1}{2} \left( 1 - \frac{1}{2} \|\prod_{i=1}^n P_i - \prod_{i=1}^n Q_i\|_1 \right)
\]

provided that \(|\gamma(P) - \gamma(Q)| \geq \epsilon\). From Pollard (1993), if \(dQ_i/dP_i = 1 + \Delta_i(\cdot)\), then

\[
\|\prod_{i=1}^n P_i - \prod_{i=1}^n Q_i\|_1 \leq \exp \left( \sum_{i=1}^n \nu_i^2 \right) - 1,
\]
where \( \nu_i^2 = \mathbb{E}_P[\Delta_i^2(\cdot)] \) is finite. So

\[
\inf \sup_{\gamma} P(\gamma - \gamma(P) > \epsilon/2) \geq \frac{1}{2} \left( \frac{3}{5} - \exp \left( \frac{n}{\sqrt{2}} \right) \right)
\]

(29)

provided that \( |\gamma(P) - \gamma(Q)| \geq \epsilon \).

It remains to find probabilities \( P \) and \( Q \) that are difficult to distinguish by the data set \( \{ (x_i, q_i, y_i) \}_{i=1}^n \).

Under \( P \), the data is generated according to

\[
y_i = g_P(x_i, q_i) + \Delta m_P (x_i, q_i) 1(q_i \leq \gamma_P) + \epsilon_i,
\]

where \( \Delta m_P (x, q_i) = \delta_{\alpha_P} + x_P \delta_{P_P} + q_i \delta_{q_P} \), and under \( Q \), \( g_P \), \( \Delta m_P \) and \( \gamma_P \) are changed to \( g_Q \), \( \Delta m_Q \) and \( \gamma_Q \), respectively. The point here is that only \( \Delta m_P \) instead of \( \delta_P \) matters for our purpose. We now specify \( g \), \( \Delta m \) and \( \gamma \) for each model. First suppose \( n \frac{\sup}{\inf} \rho_n \to \infty \). For \( P \), let \( g_P = 0 \), \( \Delta m_P = 0 \), and \( \gamma_P = 0 \) without loss of generality; for \( Q \), let

\[
g_Q(x, q) = 0, \ \Delta m_Q (x, q) = \xi \rho_n, \ \gamma_Q = \left( \xi n \rho_n^2 \right)^{-1},
\]

where \( \xi \) is a positive constant. Obviously, \( g_Q(x, q) \in \mathcal{C}_0(B, \mathcal{X} \times [0, 1]) \) for some \( B > 0 \), so it remains to compute the \( L_1 \) distance between the two measures. Let the density of \( Q \), with respect to \( P_i \) be \( 1 + \Delta_i(\cdot) \), then

\[
\Delta_i(x_i, q_i, y_i) = \begin{cases}
\phi(y_i - \Delta m_Q (x_i, q_i))/\phi(y_i) - 1, & \text{if } q_i \in [0, \gamma_Q], \\
0, & \text{otherwise}
\end{cases}
\]

where \( \phi(\cdot) \) is the standard normal pdf. Therefore,

\[
\mathbb{E}_{P_i} [\Delta_i^2] = \int_0^\gamma \int_0^1 \int_0^1 \int_{-\infty}^{\infty} \phi(y - \Delta m_Q (x, q))^2/\phi(y)^2 dy dx dq
\]

\[
= \int_0^\gamma \int_0^1 \int_{-\infty}^{\infty} \phi(y - \Delta m_Q (x, q))^2 f_{y|x|} dy dx dq
\]

Plugging in the standard normal pdf yields

\[
\mathbb{E}_{P_i} [\Delta_i^2] = \int_0^\gamma \int_0^1 \int_{-\infty}^{\infty} \frac{1}{\sqrt{2\pi}} \exp \left\{ -\frac{(y - \Delta m_Q (x, q))^2}{2} \right\} dy dx dq - \gamma_Q
\]

\[
= \int_0^\gamma \int_0^1 \int_{-\infty}^{\infty} \phi(\Delta m_Q (x, q))^2 dx dq - \gamma_Q
\]

\[
= \int_0^\gamma \exp (\xi^2 \rho_n^2) dq - \gamma_Q
\]

\[
= \gamma_Q [\exp (\xi^2 \rho_n^2) - 1] = \gamma_Q \xi^2 \rho_n^2 (1 + o(1)) \leq \frac{\xi}{n},
\]

when \( n \) is large enough.

When \( \xi \) is small enough, say \( \xi \leq \log(5/4) \), we have

\[
\exp \left( \frac{n}{\sqrt{2}} \right) \leq \exp (\xi) < \frac{5}{4}.
\]

It follows from (29) that

\[
\inf \sup_{\gamma} P(\gamma - \gamma(P) > \frac{\xi}{n} (n \rho_n^2)^{-1}) \geq \frac{1}{2} \left( \frac{3}{2} - \frac{5}{4} \right) = \frac{1}{8} \geq C,
\]
on choosing $C \leq 1/8$, where $\frac{2}{n} \left(n\rho_n^2\right)^{-1}$ appears because $|\gamma(P) - \gamma(Q)| = \left(\xi n\rho_n^2\right)^{-1} \geq \epsilon \left(n\rho_n^2\right)^{-1}$ for a small $\epsilon$.

We next suppose $n^{2+\epsilon} \rho_n = O(1)$. Let $P$ and $Q$ be the same as above except that in $Q$,

$$g_Q(x, q) = -\xi \eta^s \varphi_q \left(\frac{q - \gamma_q}{\eta}\right) \cdot \Delta m_Q(x, q) = \xi \eta^s, \quad \gamma_Q = \xi,$$

where $\eta = n^{1/(2s+1)}$, $\varphi_q$ is an infinitely differentiable function in $q$ satisfying (i) $\varphi_q(v) = 0$ for $v \geq 0$, (ii) $\varphi_q(v) = 1$, for $v \leq -1$, and (iii) $\varphi_q(v) \in (0, 1)$ for $v \in (-1, 0)$. It is not hard to check that $g_Q(x, q) \in C_{\alpha}(B, X \times [0, 1])$ for some $B > 0$. By similar steps above, we can show

$$\mathbb{E}_P[\Delta^2] \leq \frac{\xi^2}{2n}$$

when $n$ is large enough. By choosing $\xi$ appropriately, we have

$$\inf_{\tilde{\gamma}} \sup_{P \in \mathcal{P}(s, B)} \mathbb{P}(\left|\tilde{\gamma} - \gamma(P)\right| > \frac{\xi}{2}) \geq C,$$

where we choose $\epsilon \leq \xi$. ■

**Proof of Theorem 6**. We apply Theorem 2.7 of Kim and Polland (1990) to derive the asymptotic distribution of $n\rho_n^2(\hat{\gamma} - \gamma_0)$. Note that $n\rho_n^2(\hat{\gamma} - \gamma_0) = \arg \max_v \Sigma_n \left(\hat{Q}_n(\gamma_0) - \bar{Q}_n(\gamma_0)\right) =: \arg \max_{\gamma_0}(C(v))$, where $\gamma_0 = \gamma_0 + \frac{v}{n\rho_n^2}$.

(i) $C_n(v) \rightarrow C(v) \in C_{\text{max}}(\mathbb{R})$, where

$$C(v) = \Sigma^{1/2}W(v) - 2k_+(0)f_q(\gamma_0)D |v|,$$

where $W(v) := W_1(-v)1(v \leq 0) + W_2(v)1(v > 0)$ is a two-sided Brownian motion, $D = \lim_{n \rightarrow \infty} D_n$, and $\Sigma(v) = \lim_{n \rightarrow \infty} \Sigma_n$ with $\Sigma_n$ defined in Proposition 8. $C_{\text{max}}(\mathbb{R})$ is defined as the subset of continuous functions $x(\cdot) \in \mathcal{B}_{\text{loc}}(\mathbb{R})$ for which (i) $x(t) \rightarrow -\infty$ as $|t| \rightarrow \infty$ and (ii) $x(t)$ achieves its maximum at a unique point in $\mathbb{R}$, and $\mathcal{B}_{\text{loc}}(\mathbb{R})$ is the space of all locally bounded real functions on $\mathbb{R}$, endowed with the uniform metric on compacta. The weak convergence can be proved by combining Proposition 8 and Lemma 12. We now check $C(v) \in C_{\text{max}}(\mathbb{R})$. It is not hard to check $C(v)$ is continuous, has a unique minimum (see Lemma 2.6 of Kim and Pollard (1990)), and $\lim_{|v| \rightarrow \infty} C(v) = -\infty$ almost surely (which is true since $\lim_{|v| \rightarrow \infty} W_1(v)/|v| = 0$ almost surely).

(ii) $n\rho_n^2(\hat{\gamma} - \gamma_0) = O_p(1)$. This is proved in Proposition 8.

So

$$n\rho_n^2(\hat{\gamma} - \gamma_0) \xrightarrow{d} \arg \max_v \{C(v)\}.$$

Making the change-of-variables $v = \frac{V_1}{f_q(\gamma_0)D^2}r$, and noting the distributional equality $W_1(a^2r) \xrightarrow{d} aW_1(r)$,
we can rewrite the asymptotic distribution as

\[
\begin{align*}
\arg\max_v \{ C(v) \} &= \frac{V_1}{f_q(\gamma_0)D^2} \arg\max_r \left\{ C \left( \frac{V_1}{f_q(\gamma_0)D^2} r \right) \right\} \\
&= \frac{V_1}{f_q(\gamma_0)D^2} \arg\max_r \left\{ \begin{array}{ll}
4k_+(0) \sqrt{f_q(\gamma_0)V_1 W_1 \left( -\frac{V_1}{f_q(\gamma_0)D^2} r \right)} & - 2k_+(0)f_q(\gamma_0)D^2 \frac{V_1}{f_q(\gamma_0)D^2} r, \quad \text{if } r \leq 0, \\
4k_+(0) \sqrt{f_q(\gamma_0)V_2 W_2 \left( \frac{V_1}{f_q(\gamma_0)D^2} r \right)} & - 2k_+(0)f_q(\gamma_0)D^2 \frac{V_1}{f_q(\gamma_0)D^2} r, \quad \text{if } r > 0,
\end{array} \right.
\end{align*}
\]

\[
\begin{align*}
&= \frac{V_1}{f_q(\gamma_0)D^2} \arg\max_r \left\{ \begin{array}{ll}
\frac{V_1}{V_1 D^2} W_1 (-r) - \frac{1}{2} \frac{V_1}{V_1} \left| r \right|, & \text{if } r \leq 0, \\
\frac{V_1}{V_1 D^2} W_2 (r) - \frac{1}{2} \frac{V_1}{V_1} \left| r \right|, & \text{if } r > 0,
\end{array} \right.
\end{align*}
\]

\[
= \frac{V_1}{f_q(\gamma_0)D^2} \arg\max_r \left\{ \begin{array}{ll}
\frac{V_1}{V_1 D^2} W_1 (-r) - \frac{1}{2} \left| r \right|, & \text{if } r \leq 0, \\
\frac{V_1}{V_1 D^2} W_2 (r) - \frac{1}{2} \left| r \right|, & \text{if } r > 0,
\end{array} \right.
\]

\[
=: \frac{V_1}{f_q(\gamma_0)D^2} \Lambda (\lambda).
\]

\[\text{Proof of Corollary 3.}\quad \text{We mimic the proof of Theorem 6. Note that } nh^{d-1} \Delta_0^2 (\tilde{\gamma} - \gamma_0) = \arg \max_v nh^d \left( \tilde{Q}_n(\gamma_0^\nu) - Q_n(\gamma_0) \right) =: \arg \max \{ C_o(v) \}, \text{ where } \gamma_0^\nu = \gamma_0 + \frac{v}{nh^{d-1} \Delta_0^2}.
\]

(i) \( C_o(v) \sim C_o(v) \in C_{\max}(\mathbb{R}), \) where

\[
C_o(v) = \Sigma_o(v)^{1/2} W(v) - 2k_+(0)f(x_o, \gamma_0^\nu)^2 |v|.
\]

The weak convergence can be proved by combining Proposition 11 and Lemma 17

(ii) \( nh^{d-1} \Delta_0^2 (\tilde{\gamma} - \gamma_0) = O_p(1). \) This is proved in Proposition 10

So

\[
nh^{d-1} \Delta_0^2 (\tilde{\gamma} - \gamma_0) \overset{d}{\rightarrow} \arg \max_v \{ C_o(v) \}.
\]

Making the change-of-variables \( v = \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} r, \) and noting the distributional equality \( W_\epsilon(a^2 r) = a W_\epsilon(r), \) we can rewrite the asymptotic distribution as

\[
\begin{align*}
\arg\max_v \{ C_o(v) \} &= \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} \arg\max_r \left\{ C_o \left( \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} r \right) \right\} \\
&= \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} \arg\max_r \left\{ \begin{array}{ll}
4k_+(0) \sqrt{f(x_o, \gamma_0)\sigma^2(x_o) W_1 \left( \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} r \right)} & - 2k_+(0)f(x_o, \gamma_0)^2 \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} r, \quad \text{if } r \leq 0, \\
4k_+(0) \sqrt{f(x_o, \gamma_0)\sigma^2(x_o) W_2 \left( \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} r \right)} & - 2k_+(0)f(x_o, \gamma_0)^2 \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} r, \quad \text{if } r > 0,
\end{array} \right.
\end{align*}
\]

\[
= \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} \arg\max_r \left\{ \begin{array}{ll}
f(x_o, \gamma_0)\sigma^2(x_o) W_1 (-r) - \frac{1}{2} f(x_o, \gamma_0) \sigma^2(x_o) \left| r \right|, & \text{if } r \leq 0, \\
f(x_o, \gamma_0)\sigma^2(x_o) W_2 (r) - \frac{1}{2} f(x_o, \gamma_0) \sigma^2(x_o) \left| r \right|, & \text{if } r > 0,
\end{array} \right.
\]

\[
=: \frac{\sigma^2(x_o)}{f(x_o, \gamma_0)} \Lambda (\lambda_o, \kappa).
\]

\[\text{Proof of Corollary 4.}\quad \text{From the proof of Theorem 6 and the CMT, we have}
\]

\[
nh \left( \tilde{Q}_n(\gamma_0^\nu) - Q_n(\gamma_0) \right) \overset{d}{\rightarrow} \sup_v \{ C(v) \},
\]

\[8\]
where
\[
\sup_v \{C(v)\} = \sup_r \left\{ \begin{array}{ll}
4k_+(0) \sqrt{f_q(\gamma_0) W_1} - 2k_+(0) f_q(\gamma_0) D W_1 | \frac{V_1}{f_q(\gamma_0) D} r |, & \text{if } r \leq 0, \\
4k_+(0) \sqrt{f_q(\gamma_0) W_2} - 2k_+(0) f_q(\gamma_0) D W_1 | \frac{V_1}{f_q(\gamma_0) D} r |, & \text{if } r > 0,
\end{array} \right.
\]

\[
= 4k_+(0) \sup_r \left\{ \begin{array}{ll}
\frac{V_1}{D} W_1 (-r) - \frac{1}{2} | r |, & \text{if } r \leq 0, \\
\sqrt{W_2 (r) - \frac{1}{2} | r |}, & \text{if } r > 0,
\end{array} \right.
\]

\[
= 4k_+(0) \frac{V_1}{D} \sup_r \left\{ \begin{array}{ll}
W_1 (-r) - \frac{1}{2} | r |, & \text{if } r \leq 0, \\
\sqrt{W_2 (r) - \frac{1}{2} | r |}, & \text{if } r > 0.
\end{array} \right.
\]

So \( \sup_v \{C(v)\} = 4k_+(0) \frac{V_1}{D} \max \{M_1, M_2\} =: 4k_+(0) \frac{V_1}{D} M \), where \( M_1 = \sup_{r \leq 0} \{W_1 (-r) - \frac{1}{2} | r |\} \), \( M_2 = \sup_{r \geq 0} \{\sqrt{W_2 (r) - \frac{1}{2} | r |}\} \), and \( M_1 \) and \( M_2 \) are independent. From Bhattacharya and Brockwell (1976), \( M_1 \) follows the standard exponential function, and \( M_2 \) follows an exponential distribution with mean \( \lambda \). It follows that
\[
P (M \leq x) = P (M_1 \leq x, M_2 \leq x) = P (M_1 \leq x) P (M_2 \leq x) = (1 - e^{-x})(1 - e^{-x/\lambda}).
\]

By Slutsky’s theorem, the required result follows.

**Proof of Theorem 7** Because
\[
\hat{e}_i = y_i - \lambda i - x_i' \hat{\beta} - x_i' \hat{\delta} 1 (q_i \leq \tilde{\gamma}) = u_i + \left[ m_i - x_i' \hat{\beta} - x_i' \hat{\delta} 1 (q_i \leq \tilde{\gamma}) \right] \equiv u_i + D_i,
\]
we decompose \( I_n^{(1)} \) as
\[
I_n^{(1)} = \frac{n h^{d/2}}{n(n-1)} \sum_i \sum_{j \neq i} [D_i D_j + u_i u_j + 2u_i D_j] K_{h,ij}
\equiv I_{1n}^{(1)} + I_{2n}^{(1)} + I_{3n}^{(1)}.
\]

We complete the proof by examining \( I_{1n}^{(1)}, I_{2n}^{(1)}, I_{3n}^{(1)} \), and showing that \( I_{1n}^{(1)} \sim \Sigma^{(1)} + o_p(1) \) under \( H_0^{(1)} \) and the local alternative and \( I_{2n}^{(1)} \) is \( O_p(1) \) under \( H_1^{(1)} \). Throughout this proof, \( z_i = (x_i, q_i, u_i)' \) and \( \mathbb{E}_i \left[ \cdot \right] = \mathbb{E} \left[ \cdot | x_i, q_i \right] \).

It is shown in Proposition 12 of this paper that \( I_{1n}^{(1)} = O_p(h^{d/2}) \) under \( H_0^{(1)} \) and converges to \( \Delta \) under the local alternative. It can also be shown that \( I_{3n}^{(1)} = O_p(h^{d/2}) \) under \( H_0^{(1)} \) and is dominated by \( I_n \) under the alternative, see, e.g., Zheng (1996). Proposition 14 shows that \( I_{2n}^{(1)} \sim N_0 \left( 0, \Sigma^{(1)} \right) \), and Proposition 15 shows the results related to \( I_{3n}^{(1)} \). The proof is then complete.

**Proof of Theorem 8** First, decompose \( I_n^{(2)} \) by using (10):
\[
I_n^{(2)} = \frac{n h^{d/2}}{n(n-1)} \sum_i \sum_{j \neq i} 1_i^j 1_j^i \left\{ (m_i - \tilde{m}_i)(m_j - \tilde{m}_j) + u_i u_j + \tilde{u}_i \tilde{u}_j + 2u_i (m_j - \tilde{m}_j) - 2\tilde{u}_i (m_j - \tilde{m}_j) - 2u_i \tilde{u}_j \right\} K_{h,ij}
\equiv I_{1n}^{(2)} + I_{2n}^{(2)} + I_{3n}^{(2)} + 2I_{4n}^{(2)} - 2I_{5n}^{(2)} - 2I_{6n}^{(2)}.
\]

We complete the proof by examining \( I_{1n}^{(2)}, \ldots, I_{6n}^{(2)} \), and showing that \( I_{1n}^{(2)} \sim \Sigma^{(2)} + o_p(1) \) under both \( H_0^{(2)} \) and \( H_1^{(2)} \). Throughout this proof, \( z_i = (x_i, q_i, u_i)' \) and \( \mathbb{E}_i \left[ \cdot \right] = \mathbb{E} \left[ \cdot | x_i, q_i \right] \). We show that \( I_{2n}^{(2)} \) contributes to the asymptotic distribution under the null, and \( I_{1n}^{(2)} \) contributes to the power under the local alternative. All other terms will not contribute to the asymptotic distribution under either the null or the alternative; that proof just extends Propositions 3, 4, 5 and 6 in Appendix B of Porter and Yu (2011), so it is omitted here.
The remaining part of the proof concentrates on \( I_{in}^{(2)} \) and \( J_{2n}^{(2)} \), and we only briefly mention the results for the other terms since these are obtained in a similar fashion.

First, \( I_{2n}^{(2)}, I_{3n}^{(2)} \) and \( J_{6n}^{(2)} \) are invariant under \( H_0^{(2)} \) and \( H_1^{(2)} \). It can be shown that \( I_{3n}^{(2)} \) and \( J_{6n}^{(2)} \) are both \( O_p(1) \). Proposition 14 shows that \( I_{2n}^{(2)} \rightarrow_d N \left( 0, \Sigma^{(2)} \right) \).

Under \( H_0^{(2)} \), Proposition 13 shows that \( I_{1n}^{(2)} = O_{P_m} (1) \), and it can also be shown that \( I_{4n} \) and \( I_{5n} \) are both \( O_{P_m} (1) \) uniformly in \( m(\cdot) \in H_0 \).

Under \( H_1^{(2)} \), it can be shown that \( I_{4n}^{(2)} \) and \( I_{5n}^{(2)} \) are dominated by \( I_{1n}^{(2)} \), and Proposition 13 shows that \( I_{1n}^{(2)} = O_p \left( nh^{d/2} b \right) \) under \( H_1^{(2)} \). The local power can be easily obtained from the proof of Proposition 13.

Finally, Proposition 16 shows that \( v_n^{(2)} = \Sigma (2) + o_p (1) \) under both \( H_0^{(2)} \) and \( H_1^{(2)} \). So the proof is complete.

**Proof of Theorem 9** This proof is similar but more tedious than the proofs of Theorem 7 and 8. Note that \( \Phi (z) \) is a continuous function. By Pólya’s theorem, it suffices to show that for any fixed value of \( z \in \mathbb{R} \),

\[
\left| P \left( T_n^{(i)} \leq z | F_n \right) - \Phi (z) \right| = o_p (1).
\]

For the first test, let

\[
D_i^* = \mathbf{x}_i^T \hat{\beta} + \mathbf{x}_i^T \hat{\delta} 1 (q_i \leq \hat{\gamma}) - \mathbf{x}_i^T \hat{\beta}^* - \mathbf{x}_i^T \hat{\delta}^* 1 (q_i \leq \hat{\gamma}^*),
\]

where \( \hat{\beta}^*, \hat{\delta}^*, \hat{\gamma}^* \) is the least squares estimator using the data \( \{ y_i^n, x_i, q_i \}_{i=1}^n \). Then

\[
I_{1n}^{(1)*} = \frac{nh^{d/2}}{m(n-1)} \sum_{i,j} \left[ D_i^* D_j^* + u_i^* u_j^* + 2u_i^* D_j^* \right] K_{h,ij}
\]

\[
\equiv I_{1n}^{(1)*} + I_{2n}^{(1)*} + I_{3n}^{(1)*}.
\]

The theorem is proved if we can show that \( I_{in}^{(1)*} | F_n = o_p (1) \) for \( i = 1 \) and \( 3 \) and \( I_{2n}^{(1)*} / v_n^{(1)*} | F_n \rightarrow N (0, 1) \) in probability. The first part can be proved as in the proof of Theorem 4 and, for the second part, see the discussion below.

For the second test, denote \( m_i^* = \hat{y}_i \) and define \( \hat{m}_i^* \) and \( \hat{u}_i^* \) by

\[
\hat{m}_i^* = \frac{1}{n-1} \sum_{j \neq i} m_j^* L_{b,ij} / \hat{f}_i,
\]

and

\[
\hat{u}_i^* = \frac{1}{n-1} \sum_{j \neq i} u_j^* L_{b,ij} / \hat{f}_i.
\]

Then using \( \hat{e}_i^* = y_i^* - \hat{y}_i = m_i^* + u_i^* - (\hat{m}_i^* + \hat{u}_i^*) \), we get

\[
I_{1n}^{(2)*} = \frac{nh^{d/2}}{m(n-1)} \sum_{i,j} \left\{ \frac{1}{\hat{f}_i^*} \right\} \left\{ (m_i^* - \hat{m}_i^*) (m_j^* - \hat{m}_j^*) + u_i^* u_j^* + \hat{u}_i^* \hat{u}_j^* + 2u_i^* (m_j^* - \hat{m}_j^*) - 2u_i^* (m_j^* - \hat{m}_j^*) - 2u_i^* \hat{u}_j^* \right\} K_{h,ij}
\]

\[
\equiv I_{1n}^{(2)*} + I_{2n}^{(2)*} + I_{3n}^{(2)*} + 2I_{4n}^{(2)*} - 2I_{5n}^{(2)*} - 2I_{6n}^{(2)*}.
\]

The theorem is proved if we can show that \( I_{in}^{(2)*} | F_n = o_p (1) \) for \( i = 1, 3, 4, 5, 6 \) and \( I_{2n}^{(2)*} / v_n^{(2)*} | F_n \rightarrow N (0, 1) \) in probability. The first part is similar to that of Theorem 8 under \( H_0^{(2)} \). However, note that \( m^* (\cdot) | F_n \) as defined above satisfies \( H_0^{(2)} \) even if \( m(\cdot) \) is from \( H_1^{(2)} \); see Gu et al. (2007) for a similar analysis in testing omitted variables. But there is some differences in showing the second part.
First, because $u_i^s | F_n$ are mean zero and mutually independent and have variance $c_i^2$,

$$\frac{nh^{3/2}}{n(n-1)} \sum_i \sum_{j \neq i} 1_i^1 1_j^1 u_i^s u_j^s K_{h,ij} = \frac{2nh^{1/2}}{n(n-1)} \sum_i \sum_{j > i} 1_i^1 1_j^1 u_i^s u_j^s K_{h,ij} \equiv \sum_i \sum_{j > i} U_{n,ij}^*$$

is a second order degenerate $U$-statistic with conditional variance

$$\frac{2h^2}{n(n-1)} \sum_i \sum_{j \neq i} 1_i^1 1_j^1 c_i^2 c_j^2 K_{h,ij} = v_n^2.$$ 

Because $U_{n,ij}^*$ depends on $i$ and $j$, we use the central limit theorem of de Jong (1987) for generalized quadratic forms rather than Hall (1984) to find the asymptotic distribution of $I_{2n}^{(2)*}$. From his Proposition 3.2, we know $I_{2n}^{(2)*} / v_{n}^{(2)} | F_n \rightarrow N(0,1)$ in probability as long as

$$G_1^* = \sum_i \sum_{j > i} \mathbb{E}^* \left[ U_{n,ij}^* \right] = o_p \left( v_{n}^{(2)} \right),$$

$$G_{II}^* = \sum \sum_{j > i} \sum \sum_{j > i} \mathbb{E}^* \left[ U_{n,ij}^* U_{n,jl}^* + U_{n,ij}^* U_{n,jl}^* + U_{n,ij}^* U_{n,jl}^* \right] = o_p \left( v_{n}^{(2)} \right),$$

$$G_{IV}^* = \sum \sum_{j > i} \sum \sum_{j > i} \mathbb{E}^* \left[ U_{n,ij}^* U_{n,jl}^* U_{n,jl}^* U_{n,jl}^* + U_{n,ij}^* U_{n,jl}^* U_{n,jl}^* U_{n,jl}^* \right] = o_p \left( v_{n}^{(2)} \right).$$

It is straightforward to show that

$$G_1^* = O_p \left( (nhd)^{-1} \right), G_{II}^* = O_p(n^{-1}), G_{IV}^* = O_p(h^d),$$

see, e.g., the proof of Theorem 2 of Hsiao et al. (2007), so the result follows by $v_{n}^{(2)} = O_p(1)$. Next, it is easy to check that $\mathbb{E}^* \left[ v_{n}^{(2)} \right] = v_{n}^{(2)} + o_p(1)$, and $\text{Var}^* \left( v_{n}^{(2)} \right) = o_p(1)$. Thus $I_{2n}^{(2)*} / v_{n}^{(2)} | F_n \rightarrow N(0,1)$ in probability. The analysis for $I_{1n}^{(1)*}$ is similar. ■

**Supplement B: Propositions**

**Proposition 1** $\hat{\gamma} - \gamma_0 = O_p(h)$.

**Proof.** We apply Lemma 4 of Porter and Yu (2015) to prove the result. By Lemma B.1 of Newey (1994), we have

$$\sup_{\gamma \in \Gamma} \left| \hat{Q}_n (\gamma) - Q_n (\gamma) \right| = O_p \left( \sqrt{\ln n / nh^d} \right) \rightarrow 0,$$

where

$$Q_n (\gamma) = \int \left[ f_{x} \left\{ \sum_{x \neq h} \mathbb{E} \left[ x \right] \right\} \right] f_0 \right) \ dx \ dx .$$

Let $\mathcal{N}_n = [\gamma_0 - h, \gamma_0 + h]$ and $\mathcal{N}_n = \text{Arg Max}_{\gamma \in \Gamma} Q_n (\gamma)$; then it is easy to show that $\sup_{\gamma \in \Gamma \setminus \mathcal{N}_n} Q_n (\gamma) = O (h^2)$. But for $\gamma \in \mathcal{N}_n$, the result is different. Specifically, let $\gamma = \gamma_0 + ah$, $a \in (0, 1)$ . Then

$$Q_n (\gamma) = \int \left[ f_{x} \left\{ \sum_{x \neq h} \mathbb{E} \left[ x \right] \right\} \right] f_0 \right) \ dx \ dx .$$

The differences of the first and the third terms in brackets are still $O (h^2)$, so the second term will dominate.
With Assumption I, we have \(\int_{\mathbb{R}} \left[ f K^z(u, x) (1, x, \gamma_0) \delta_0 f(x, \gamma_0) du \right]^2 f(x) dx > C\). Under Assumption K, if \(a \in (0, 1)\), then \(\int_{-1}^a k_-(u) du \leq 1\) and \(\int_{-1}^a k_-(u) du\) is a nonincreasing function of \(a\) for \(a \in (0, 1)\). As a result, \(Q_n(\gamma)\) is a nonincreasing function of \(a\) for \(a \in (0, 1)\) up to \(O(h^2)\). Similarly, \(Q_n(\gamma)\) is a nondecreasing function of \(a\) for \(a \in (-1, 0)\) up to \(O(h^2)\). So \(Q_n(\gamma)\) is maximized at some \(\gamma_n \in \mathcal{N}_n\) such that \(Q_n(\gamma_n) > \sup_{\gamma \in \Gamma \cap \mathcal{N}_n} |Q_n(\gamma)| + C/2\) for \(n\) large enough. The result of interest then follows.

**Proposition 2** \( \hat{\gamma} - \gamma_0 = O_p \left( (n/h)^{-1/2} \right) \).

**Proof.** We apply the standard shelling method to obtain the result. Specifically, for each \(n\), the parameter space is partitioned into the ‘shells’ \(S_{i,n} = \{ \gamma : 2^{i-1} \leq (n/h)^{1/2} |\gamma - \gamma_0| \leq 2^i \}\) with \(i\) ranging over the integers. If \((n/h)^{1/2} |\hat{\gamma} - \gamma_0|\) is larger than \(2^l\) for a given integer \(L\), then \(\hat{\gamma}\) is in one of the shells \(S_{i,n}\) with \(i \geq L\). In that case the supremum of the map \(\gamma \rightarrow \hat{Q}_n(\gamma) - \hat{Q}_n(\gamma_0)\) over this shell is nonnegative by the property of \(\hat{\gamma}\). Note that

\[
P \left( (n/h)^{1/2} |\hat{\gamma} - \gamma_0| > 2^L \right) \
\leq P \left( \max_{2^{L-1} \leq (n/h)^{1/2} |\hat{\gamma} - \gamma_0| < (n/h)^{1/2}} \sum_{i=1}^{n} \hat{\Delta}_i^2(\gamma) - \frac{1}{n} \sum_{i=1}^{n} \hat{\Delta}_i^2(\gamma_0) \geq 0 \right) + P (|\hat{\gamma} - \gamma_0| \geq h) \
\leq \sum_{i=L}^{\frac{1}{2} \log_2(nh)} P \left( \sup_{S_{i,n}} \max_{2^{L-1} \leq (n/h)^{1/2} |\hat{\gamma} - \gamma_0| < (n/h)^{1/2}} \sum_{i=1}^{n} \hat{\Delta}_i^2(\gamma) - \frac{1}{n} \sum_{i=1}^{n} \hat{\Delta}_i^2(\gamma_0) \geq 0 \right) + P (|\hat{\gamma} - \gamma_0| \geq h) \
\leq \sum_{i=L}^{\frac{1}{2} \log_2(nh)} P \left( \sup_{S_{i,n}} \max_{2^{L-1} \leq (n/h)^{1/2} |\hat{\gamma} - \gamma_0| < (n/h)^{1/2}} \sum_{i=1}^{n} \hat{\Delta}_i^2(\gamma) - \frac{1}{n} \sum_{i=1}^{n} \hat{\Delta}_i^2(\gamma_0) 1(\Delta_i > 0) \right) \
+ \sum_{i=L}^{\frac{1}{2} \log_2(nh)} P \left( \sup_{S_{i,n}} \max_{2^{L-1} \leq (n/h)^{1/2} |\hat{\gamma} - \gamma_0| < (n/h)^{1/2}} \sum_{i=1}^{n} \hat{\Delta}_i^2(\gamma) - \frac{1}{n} \sum_{i=1}^{n} \hat{\Delta}_i^2(\gamma_0) 1(\Delta_i < 0) \right) \
+ P (|\hat{\gamma} - \gamma_0| \geq h) 
= T1 + T2 + T3.

As \(T3\) converges to zero by Proposition I and \(T2\) is similar to \(T1\), we only use \(T1\) to illustrate the derivations in the following discussion. Since

\[
T1 \leq \sum_{i=L}^{\frac{1}{2} \log_2(nh)} P \left( \sup_{S_{i,n}} \max_{2^{L-1} \leq (n/h)^{1/2} |\hat{\gamma} - \gamma_0| < (n/h)^{1/2}} \sum_{i=1}^{n} \hat{\Delta}_i(\gamma) - \hat{\Delta}_i(\gamma_0) 1(\Delta_i > 0) \right) \
+ \sum_{i=L}^{\frac{1}{2} \log_2(nh)} P \left( \sup_{S_{i,n}} \max_{2^{L-1} \leq (n/h)^{1/2} |\hat{\gamma} - \gamma_0| < (n/h)^{1/2}} \sum_{i=1}^{n} \hat{\Delta}_i(\gamma) + \hat{\Delta}_i(\gamma_0) 1(\Delta_i > 0) < 0 \right),
\]

we focus on the first term because the second term is easier to analyze given that \(\Delta_i > 0\). To simplify notations, \(1(\Delta_i > 0)\) is neglected in the remaining proof. Notice that

\[
\frac{1}{n} \sum_{i=1}^{n} \left( \hat{\Delta}_i(\gamma) - \hat{\Delta}_i(\gamma_0) \right) \
= \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \left( y_{ij} K_{h,j}^{\gamma_0} - y_{ij} K_{h,j}^{\gamma_0} \right) - \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \left( y_{ij} K_{h,j}^{\gamma_0} - y_{ij} K_{h,j}^{\gamma_0} \right) \
= \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \left( m_{ij} K_{h,j}^{\gamma_0} - m_{ij} K_{h,j}^{\gamma_0} \right) - \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \left( m_{ij} K_{h,j}^{\gamma_0} - m_{ij} K_{h,j}^{\gamma_0} \right) \
+ \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \left( u_{ij} K_{h,j}^{\gamma_0} - u_{ij} K_{h,j}^{\gamma_0} \right) - \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \left( u_{ij} K_{h,j}^{\gamma_0} - u_{ij} K_{h,j}^{\gamma_0} \right) 
= D1 + D2,
\]

where \(m_{ij} = g_j + (1, x_j', q_j) \delta_0 1(q_j \leq \gamma_0)\) with \(g_j = g(x_j, q_j)\). Suppose \(\gamma_0 < \gamma < \gamma_0 + h\), then for some \(C > 0\),
we have

\[ D1 = \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} g_{ij} \left( K_{h,ij}^{-} - K_{0,ij}^{-} \right) - \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} g_{ij} \left( K_{h,ij}^{+} - K_{0,ij}^{+} \right) \]

\[
+ \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \left( 1, x_{j}, q_{j} \right) \delta_{0} \left( K_{h,ij}^{-} - K_{0,ij}^{-} \right) 1 \left( q_{j} \leq \gamma_{0} \right) \]

\[
\leq -C \left( \frac{\gamma_{0}}{h} \right)^{2} \]

with probability approaching 1 by calculating the mean and variance of \( D1 \) in its U-projection, where the difference of the first two terms contribute only \( O_{p} \left( \left| \gamma - \gamma_{0} \right|^{2} \right) \), and the third term contributes to

\[-C \left( \frac{\gamma_{0}}{h} \right)^{2} \]

because for each \( i \), \( K_{h,ij}^{-} \) covers \( O_{p} \left( n \left( \gamma - \gamma_{0} \right) \right) \) terms less than \( K_{h,ij}^{-} \) given that \( \gamma > \gamma_{0} \) and \( k_{\pm} \left( 0 \right) = 0 \).

In consequence, for \( \eta \leq h \),

\[
P \left( \sup_{\left| \gamma - \gamma_{0} \right| < \eta} \frac{1}{n} \sum_{i=1}^{n} \left( \hat{\Delta}_{i} \left( \gamma \right) - \bar{\Delta}_{i} \left( \gamma_{0} \right) \right) > 0 \right) \leq P \left( \sup_{\left| \gamma - \gamma_{0} \right| < \eta} D2 > C \left( \frac{\gamma_{0}}{h} \right)^{2} \right). \]

Next,

\[ D2 = \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} u_{ij} \left( K_{h,ij}^{-} - K_{0,ij}^{-} \right) 1 \left( q_{j} \leq \gamma_{0} \right) + \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} u_{ij} \left( K_{h,ij}^{+} - K_{0,ij}^{+} \right) 1 \left( q_{j} > \gamma \right) \]

\[
+ \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \left( 1, x_{j}, q_{j} \right) \delta_{0} \left( K_{h,ij}^{-} + K_{h,ij}^{+} \right) 1 \left( \gamma_{0} < q_{j} \leq \gamma \right) =: P1 + P2 + P3, \] (30)

and we can apply Lemma 8.4 of Newey and McFadden (1994) to bound \( D2 \). Since the first two terms are similar, we just check the first term and the last term.

For the first term, set

\[
m_{n} \left( z_{i}, z_{j} \right) = u_{ij} \left( K_{h,ij}^{-} - K_{0,ij}^{-} \right) 1 \left( q_{j} \leq \gamma_{0} \right), \]

where \( z_{i} = \left( u_{i}, x_{i}', q_{i} \right)' \), and \( m_{n} \left( z_{i}, z_{j} \right) = 0 \) for any \( i = j \). So we have

\[
n^{-2} \sum_{i=1}^{n} \sum_{j=1}^{n} m_{n} \left( z_{i}, z_{j} \right) = P1, \]

\[
E \left[ m_{n} \left( z_{i}, z_{j} \right) \right] = 0, \ E \left[ m_{n} \left( z_{i}, z_{j} \right) \right] = 0, \ E \left[ m_{n} \left( z_{1}, z_{1} \right) \right]/n = 0, \]

and

\[
E \left[ \left\| m_{n} \left( z_{1}, z_{2} \right) \right\|^{2} \right] = E \left[ E \left[ \left( \frac{u}{\sqrt{K}} K_{x} \left( \frac{x_{j}-x_{i}}{h} \right) x_{i} \right) k^{-} \left( \frac{u_{j}-\gamma_{0}}{h} \right) - k^{-} \left( \frac{u_{j}-\gamma_{0}}{h} \right) \right] 1 \left( q_{j} \leq \gamma_{0} \right) \left\| x_{i} \right\| \right] \]

\[
E \left[ E \left[ \frac{u}{\sqrt{K}} K_{x} \left( \frac{x_{j}-x_{i}}{h} \right) x_{i} \right] k^{-} \left( \frac{u_{j}-\gamma_{0}}{h} \right) - k^{-} \left( \frac{u_{j}-\gamma_{0}}{h} \right) \right]^{2} 1 \left( q_{j} \leq \gamma_{0} \right) \left\| x_{i} \right\| \right] \]

\[
E \left[ \frac{1}{h} \int_{-\infty}^{\gamma_{0}} \int_{x_{i}} f_{x_{j}} \frac{u_{j}}{\sqrt{K}} K_{x} \left( \frac{x_{j}-x_{i}}{h} \right) x_{i} \right] k^{-} \left( \frac{u_{j}-\gamma_{0}}{h} \right) - k^{-} \left( \frac{u_{j}-\gamma_{0}}{h} \right) \left[ f \left( u_{j} \left| x_{j}, q_{j} \right. \right) f \left( x_{j}, q_{j} \right) du_{j} dx_{j} dq_{j} \right] \]

\[
E \left[ \frac{1}{h} \int_{-\infty}^{\gamma_{0}} \int_{x_{i}} \sigma^{2} \left( x_{i} + u_{x}h, \gamma_{0} + u_{q}h \right) K_{x} \left( u_{x}, x_{i} \right) \left[ k^{-} \left( \frac{u_{q}+\gamma_{0}}{h} \right) - k^{-} \left( u_{q} \right) \right] \right] f \left( x_{i} + u_{x}h, \gamma_{0} + u_{q}h \right) du_{x} du_{q} \]

\[
\leq C \left( \frac{\gamma_{0}}{h} \right)^{2} \]

under Assumption I and the fact that \( \left| k^{-} \left( \cdot \right) \right| < \infty \) and \( K_{x} \left( \cdot \right) < \infty \) over their supports, where \( E \left[ \cdot \right] = E \left[ \cdot \left| z_{1} \right. \right], \)
\( \sigma^2(x_j, q_j) = \mathbb{E} [u_j^2 | x_j, q_j] \). Hence, by Proposition 1,

\[
\left( \mathbb{E} \left[ \| m_n(z_1, z_2) \|^2 \right] \right)^{1/2} / n \leq C \left( \frac{2a - \gamma}{n} \right)^{1/2} = O_p \left( \frac{1}{n^{1/2}} \right) = o_p(1)
\]

under Assumption H. As a result,

\[
P_1 = \frac{1}{n(n-1)} \sum_{i=1}^n \sum_{j=1, j\neq i}^n u_j \left( K_{h, ij}^{\gamma^-} - K_{h, ij}^{\gamma^+} \right) 1(q_j \leq \gamma_0) = n^{-1} \sum_{j=1}^n \mathbb{E} \left[ m_n(z_i, z_j) | z_j \right] + o_p(1)
\]

where

\[
\begin{align*}
\mathbb{E}_j \left[ m_n(z_i, z_j) \right] &= \frac{u_j}{n} 1(q_j \leq \gamma_0) \left[ k^- \left( \frac{q_j - \gamma}{h} \right) - k^- \left( \frac{q_i - \gamma}{h} \right) \right] \int_x K^x \left( \frac{x_j - x}{h} \right) f(x_i) \, dx_i \\
&= \frac{u_j}{n} 1(q_j \leq \gamma_0) \left[ k^- \left( \frac{q_j - \gamma}{h} \right) - k^- \left( \frac{q_i - \gamma}{h} \right) \right] \int_x K^x \left( u_x, x_i \right) f(x_i - u_xh) \, du_x \\
&\leq C \frac{u_j}{n} 1(q_j \leq \gamma_0) \left[ k^- \left( \frac{q_j - \gamma}{h} \right) - k^- \left( \frac{q_i - \gamma}{h} \right) \right] f(x_i).
\end{align*}
\]

Hence,

\[
\text{Var} (P1) = Var \left( n^{-1} \sum_{j=1}^n \mathbb{E}_j \left[ m_n(z_i, z_j) \right] + o_p(1) \right)
\leq n^{-1} \int_{-\infty}^\infty \int_x \frac{u_j^2}{n^2} \left[ k^- \left( \frac{q_i - \gamma}{h} \right) - k^- \left( \frac{q_j - \gamma}{h} \right) \right]^2 f(x_j)^2 f(u_j | x_j, q_j)^2 f(x_j, q_j) f(u_j, x_j) \, dx_j du_j \, dq_j
\leq C \frac{1}{nh} \int_{-\infty}^\infty \sigma^2(x_j, \gamma_0 + u_q h) \left[ k^- \left( u_j + \frac{\gamma - \gamma_0}{h} \right) - k^- \left( u_j \right) \right]^2 f(x_j, \gamma_0 + u_q h) f(x_j)^2 \, dx_j \, du_q
\leq C \frac{1}{nh} \left( \frac{\gamma - \gamma_0}{h} \right)^2.
\]

Similarly for \( P_3 \), we can set

\[
P_3 = \frac{1}{n(n-1)} \sum_{i=1}^n \sum_{j=1, j\neq i}^n u_j \left( K_{h, ij}^{\gamma^-} + K_{h, ij}^{\gamma^+} \right) 1(q_j < \gamma \leq \gamma_0) = n^{-2} \sum_{i=1}^n \sum_{j=1}^n m_n(z_i, z_j)
\]

with

\[
\mathbb{E} \left[ m_n(z_i, z_j) \right] = 0, \quad \mathbb{E}_i \left[ m_n(z_i, z_j) \right] = 0, \quad \mathbb{E} \left[ \| m_n(z_1, z_2) \| \right] / n = 0,
\]

and

\[
\begin{align*}
\mathbb{E} \left[ \| m_n(z_1, z_2) \|^2 \right]
&= \mathbb{E} \left[ \left( \frac{u_j}{h} K^x \left( \frac{x_j - x_i}{h} \right) \left[ k^- \left( \frac{q_j - \gamma}{h} \right) - k^+ \left( \frac{q_i - \gamma}{h} \right) \right] \int_x K^x \left( \frac{x_j - x}{h} \right) f(x_i) \, dx_i \right)^2 \right] \\
&= \mathbb{E} \left[ \frac{u_j}{h} \int_{-\infty}^{\gamma_0} \int_{x_j}^{\gamma_0} \sigma^2(x_j, u_x h, \gamma_0 + u_q h) K^x(u_x, x_i)^2 \left[ k^- \left( u_j + \frac{\gamma - \gamma_0}{h} \right) - k^+ \left( u_q \right) \right]^2 f(x_j + u_x h, \gamma_0 + u_q h) \, du_x \, du_q \right]
\leq C \frac{1}{h^2} \left( \frac{\gamma - \gamma_0}{h} \right)^3.
\end{align*}
\]

By Proposition 1, we have

\[
\left( \mathbb{E} \left[ \| m_n(z_1, z_2) \|^2 \right] \right)^{1/2} / n \leq C \left( \frac{\gamma - \gamma_0}{n} \right)^{3/2} = O_p(1) = o_p(1)
\]

under Assumption H. As a result,

\[
P_3 = \frac{1}{n(n-1)} \sum_{i=1}^n \sum_{j=1, j\neq i}^n u_j \left( K_{h, ij}^{\gamma^-} - K_{h, ij}^{\gamma^+} \right) 1(q_j \leq \gamma_0) = n^{-1} \sum_{j=1}^n \mathbb{E}_j \left[ m_n(z_i, z_j) \right] + o_p(1)
\]
with
\[ E_j \left[ m_n (z_i, z_j) \right] = \frac{u_i}{h^n} \left[ k^- \left( \frac{q_i - \gamma}{h} \right) - k^+ \left( \frac{q_i - \gamma}{h} \right) \right] \left( x_i - \gamma \right) f (x_i) \, dx_i \]
\[ = \frac{u_i}{h^n} \left[ k^- \left( \frac{q_i - \gamma}{h} \right) - k^+ \left( \frac{q_i - \gamma}{h} \right) \right] \left( x_i - \gamma \right) f (x_i) \, du_x \]
\[ \leq C \frac{u_i}{h^n} \left[ k^- \left( \frac{q_i - \gamma}{h} \right) - k^+ \left( \frac{q_i - \gamma}{h} \right) \right] \left( x_i - \gamma \right) f (x_i) \, dx_i \]

Hence,
\[ Var (P3) \approx Var \left( n^{-1} \sum_{j=1}^{n} E_j \left[ m_n (z_i, z_j) \right] \right) \]
\[ \leq C \frac{1}{nh} \left[ \frac{\gamma_0 - \gamma}{h} \right]^2 \left[ \frac{\gamma_0 - \gamma}{h} \right]^2 \left( x_i - \gamma \right) f (x_i) \, dx_i \, du_q \]
\[ \leq C \frac{1}{nh} \left[ \frac{\gamma_0 - \gamma}{h} \right]^3 \left( x_i - \gamma \right) f (x_i) \, dx_i \, du_q \]

Since, conditional on \( x_i \), the three summations on the right hand side of (39) are independent, we obtain
\[ Var (P2) = Var (P1) + Var (P2) + Var (P3) \]
\[ \leq C \frac{1}{nh} \left( \frac{\gamma_0 - \gamma}{h} \right)^2 \left( \frac{\gamma_0 - \gamma}{h} \right)^2 \leq C \frac{(\gamma - \gamma_0)^2}{n^h h^3} \]
uniformly for \( \gamma - \gamma_0 \leq \eta \). In consequence,
\[ \mathbb{P} \left( \sup_{|\gamma - \gamma_0| < \eta} \frac{1}{n} \sum_{i=1}^{n} \left( \hat{\Delta}_i (\gamma) - \hat{\Delta}_i (\gamma_0) \right) > 0 \right) \]
\[ \leq C \mathbb{V} \left( \sup_{|\gamma - \gamma_0| < \eta} \left( \hat{\Delta}_i (\gamma) - \hat{\Delta}_i (\gamma_0) \right) \right) \left( \frac{\gamma - \gamma_0}{h} \right)^2 \leq C \frac{(\gamma - \gamma_0)^2}{n^h h^3} \left( \frac{\gamma - \gamma_0}{h} \right)^4 = \frac{Ch}{n(\gamma_0 - \gamma)^2} \]
by Markov's inequality. So
\[ \sum_{l=L}^{L + \frac{\log \eta (nh)}{2}} \mathbb{P} \left( \sup_{|\gamma - \gamma_0| < \eta} \frac{1}{n} \sum_{i=1}^{n} \left( \hat{\Delta}_i (\gamma) - \hat{\Delta}_i (\gamma_0) \right) \right) \left( \Delta_l > 0 \right) > 0 \]
\[ \leq C \sum_{l=L}^{L + \frac{\log \eta (nh)}{2}} \frac{Ch}{n(\gamma_0 - \gamma)^2} = C \sum_{l \geq L} \frac{1}{n^{1/2}} \to 0 \]
as \( L \to \infty \). The proof is completed. \[ \square \]

**Proposition 3** For \( v \) on any compact set,
\[ \mathbb{E} \left[ n h \left( \hat{Q}_n (\gamma_0) - \check{Q}_n (\gamma_0) \right) \right] = -v^2 \mathbb{E} [\Delta f (x_i) \Delta_i f (x_i) | q_i = \gamma_0, f_q (\gamma_0) k' (0) + o(1), \]
\[ Cov \left( n h \left( \hat{Q}_n (\gamma_0^1) - \check{Q}_n (\gamma_0) \right) , n h \left( \hat{Q}_n (\gamma_0^2) - \check{Q}_n (\gamma_0) \right) \right) \]
\[ = 4 \mathbb{E} \left[ \Delta f^2 (x_i) (\sigma_1^2 (x_i) + \sigma_2^2 (x_i)) | q_i = \gamma_0 \right] f_q (\gamma_0) v_1 v_2 \zeta (0) + o(v_1 v_2), \]
and
\[ \frac{1}{\sqrt{Var (n h \left( \hat{Q}_n (\gamma_0) - \check{Q}_n (\gamma_0) \right))}} \left( n h \left( \hat{Q}_n (\gamma_0) - \check{Q}_n (\gamma_0) \right) - \mathbb{E} \left[ n h \left( \hat{Q}_n (\gamma_0) - \check{Q}_n (\gamma_0) \right) \right] \right) \xrightarrow{d} N (0, 1), \]
where \( \gamma_0 = \gamma_0 + \frac{\mu}{\sqrt{n h}} \).
Proof. Note that
\[ \hat{Q}_n(\gamma_0^v) - \tilde{Q}_n(\gamma_0) = \frac{1}{n} \sum_{i=1}^{n} \left( \hat{\Delta}_i(\gamma_0^v)^2 - \hat{\Delta}_i(\gamma_0)^2 \right) = \frac{1}{n} \sum_{i=1}^{n} \left( \hat{\Delta}_i(\gamma_0^v) + \hat{\Delta}_i(\gamma_0) \right) \left( \hat{\Delta}_i(\gamma_0^v) - \hat{\Delta}_i(\gamma_0) \right). \]

By Lemma 1, we have
\[ \hat{\Delta}_i(\gamma_0^v) - \hat{\Delta}_i(\gamma_0) = \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} y_j K_{h,i,j}^{-} - \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} y_j K_{h,i,j}^{+} \]
and by Lemma 3, we have
\[ \hat{\Delta}_i(\gamma_0^v) + \hat{\Delta}_i(\gamma_0) = \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} y_j K_{h,i,j}^{-} + \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} y_j K_{h,i,j}^{+}. \]

Next we focus on the other term. For simplicity, let \( v > 0 \). Now,
\[ \hat{\Delta}_i(\gamma_0^v) - \hat{\Delta}_i(\gamma_0) = \left( \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} y_j K_{h,i,j}^{-} \right) - \left( \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} y_j K_{h,i,j}^{+} \right) - \left( \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} y_j K_{h,i,j}^{+} \right) + \left( \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} y_j K_{h,i,j}^{-} \right) \]
\[ =: \frac{1}{n-1} \sum_{j=1,j\neq i}^{n} (T_{1ij} + T_{2ij} + T_{3ij} + T_{4ij} + T_{5ij} + T_{6ij}), \]
where
\[ T_{1ij} = g(x_j,q_j) \left( K_{h,i,j}^{+} - K_{h,i,j}^{-} \right), \]
\[ T_{2ij} = -g(x_j,q_j) \left( K_{h,i,j}^{+} - K_{h,i,j}^{-} \right), \]
\[ T_{3ij} = u_j \left( K_{h,i,j}^{+} - K_{h,i,j}^{-} \right), \]
\[ T_{4ij} = u_j \left( K_{h,i,j}^{+} - K_{h,i,j}^{-} \right), \]
\[ T_{5ij} = - (1, x_j', q_j) \delta_0 K_{h,i,j}^{+} 1(\gamma_0 - h \leq q_j \leq \gamma_0), \]
\[ T_{6ij} = - (1, x_j', q_j) \delta_0 K_{h,i,j}^{-} 1(\gamma_0 - h \leq q_j \leq \gamma_0). \]

By Lemma 4 we have
\[ \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1,j\neq i}^{n} 2\Delta f(x_i) (T_{1ij} + T_{2ij}) \approx 0, \]
by Lemma 2 we have
\[ \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1,j\neq i}^{n} 2\Delta f(x_i) (T_{3ij} + T_{4ij}) \approx \frac{2}{nh} \sum_{i=1}^{n} \Delta f(x_i) f(x_i) u_i \left( k^- \left( \frac{q_j - \gamma_0}{h} \right) - k^- \left( \frac{q_j - \gamma_0}{h} \right) \right) \left( k^+ \left( \frac{q_j - \gamma_0}{h} \right) - k^+ \left( \frac{q_j - \gamma_0}{h} \right) \right), \]
and by Lemma 3 we have
\[ \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1,j\neq i}^{n} 2\Delta f(x_i) (T_{5ij} + T_{6ij}) \approx \frac{2}{nh} \sum_{i=1}^{n} \Delta f(x_i) f(x_i) (1, x_j', q_i) \delta_0 \left[ k^- \left( \frac{q_j - \gamma_0}{h} \right) 1(\gamma_0 - h \leq q_j \leq \gamma_0) - k^- \left( \frac{q_j - \gamma_0}{h} \right) 1(\gamma_0 - h \leq q_j \leq \gamma_0) \right]. \]
Combining results we have
\[ nh \left( \hat{Q}_n(\gamma_0^v) - \tilde{Q}_n(\gamma_0) \right) \]
\[ = 2 \sum_{i=1}^{n} \Delta f(x_i) f(x_i) u_i \left( k^- \left( \frac{q_j - \gamma_0}{h} \right) - k^- \left( \frac{q_j - \gamma_0}{h} \right) \right) \left( k^+ \left( \frac{q_j - \gamma_0}{h} \right) - k^+ \left( \frac{q_j - \gamma_0}{h} \right) \right) \]
\[ + 2 \sum_{i=1}^{n} \Delta f(x_i) f(x_i) (1, x_j', q_i) \delta_0 \left[ k^- \left( \frac{q_j - \gamma_0}{h} \right) 1(\gamma_0 - h \leq q_j \leq \gamma_0) - k^- \left( \frac{q_j - \gamma_0}{h} \right) 1(\gamma_0 - h \leq q_j \leq \gamma_0) \right] + o_p(1) \]
\[ =: S_1 + S_2 + o_p(1). \]
Hence,
\[
\mathbb{E} \left[ nh \left( \hat{Q}_n(\gamma) - \hat{Q}_n(\gamma_0) \right) \right] \approx \mathbb{E} [ S_2(v) ] \\
= 2n \int_{x_i} \int_{\gamma_0}^{\gamma} \Delta f (x_i) f (x_i) (1, x'_j, q_j) \delta_0 k^\dagger \left( \frac{\eta - \gamma_0}{h} \right) \eta (x_i, q_i) dq_i dx_i \\
- 2n \int_{x_i} \int_{\gamma_0}^{\gamma} \Delta f (x_i) f (x_i) (1, x'_j, q_j) \delta_0 k^\dagger \left( \frac{\eta - \gamma_0}{h} \right) \eta (x_i, q_i) dq_i dx_i \\
= 2n \int_{x_i} \int_{\gamma_0}^{\gamma} \Delta f (x_i) f (x_i) (1, x'_j, \gamma_0 + u_j h) \delta_0 k^\dagger \left( u_j - \frac{\eta}{n h} \right) \eta (x_i, \gamma_0 + u_j h) du_q dx_i \\
- 2n \int_{x_i} \int_{\gamma_0}^{\gamma} \Delta f (x_i) f (x_i) (1, x'_j, \gamma_0 + u_j h) \delta_0 k^\dagger (u_j) \eta (x_i, \gamma_0 + u_j h) du_q dx_i \\
\approx -v^2 f_q (\gamma_0) \mathbb{E} [ \Delta f (x_i) \Delta_i f (x_i) | q_i = \gamma_0 ] k'_+ (0),
\]
where the first equality comes from the zero conditional mean property of the error term, and the last one applies \(|k'_+ (0)| = k'_+ (0)|. By Lemma 4 and Lemma 5, as well as the exogeneity property of error term \(u_i, \) we have
\[
Cov \left( nh \left( \hat{Q}_n(\gamma) - \hat{Q}_n(\gamma_0) \right) , nh \left( \hat{Q}_n(\gamma) - \hat{Q}_n(\gamma_0) \right) \right) \\
= Cov \left( S_1 (v_1) , S_1 (v_2) \right) + Cov \left( S_2 (v_1) , S_2 (v_2) \right) + Cov \left( S_1 (v_1) , S_2 (v_2) \right) + Cov \left( S_1 (v_2) , S_2 (v_2) \right) \\
\approx 4 \mathbb{E} \left[ \Delta f (x_i) f^2 (x_i) (\sigma^2_+ (x_i) + \sigma^2_ (x_i)) | q_i = \gamma_0 \right] f_q (\gamma_0) n_1 x_1 \xi_1.
\]
Roughly speaking, \(S_2(v)\) contributes to the mean process of \(nh \left( \hat{Q}_n(\gamma) - \hat{Q}_n(\gamma_0) \right) , \) and \(S_1(v)\) contributes to the variance process.

To show the weak convergence, we apply the Lyapunov CLT by checking the Lyapunov condition. Specifically, we show that
\[
\sum_{i=1}^n \mathbb{E} \left[ \Delta f (x_i) f (x_i) | q_i = \gamma_0 \right] = O_p \left( \left( \frac{nh}{(nh)^2} \right)^{1/2} \right) = o_p (1)
\]
as \(n \to \infty. \)

**Proposition 4** \(\hat{\Sigma} (\gamma_0) \d \to \Sigma.\)

**Proof.** By standard arguments, we have
\[
\mathbb{E} \left[ \frac{1}{n} \sum_{i=1}^n k_h (q_i - \gamma_0) \Delta f^2 (x_i) f^2 (x_i) u^2_i | q_i = \gamma_0 \right] \xrightarrow{p} 0
\]
and
\[
\mathbb{E} \left[ \frac{1}{n} \sum_{i=1}^n k_h (q_i - \gamma_0) \Delta f (x_i) \Delta f (x_i) | q_i = \gamma_0 \right] \xrightarrow{p} 0.
\]
So all we need to show is
\[
\mathbb{E} \left[ \frac{1}{n} \sum_{i=1}^n k_h (q_i - \gamma_0) \tilde{\Delta}^2 (\gamma_0) \tilde{f}^2 (x_i) \tilde{u}_i (\gamma_0) \right] \xrightarrow{p} 0
\]
and
\[
\mathbb{E} \left[ \frac{1}{n} \sum_{i=1}^n k_h (q_i - \gamma_0) \tilde{f}^{-1} (x_i, \gamma_0) \tilde{f} (x_i) - \frac{1}{n} \sum_{i=1}^n k_h (q_i - \gamma_0) \Delta f (x_i) \Delta f (x_i) \right] \xrightarrow{p} 0,
\]
which are implied by
\[
\tilde{f} (x_i, \gamma_0) - f (x_i, \gamma_0) \xrightarrow{p} 0, \tilde{\Delta} (\gamma_0) - \Delta f (x_i) \xrightarrow{p} 0, \tilde{f} (x_i) - f (x_i) \xrightarrow{p} 0 \text{ and } \tilde{u}_i (\gamma_0) - u_i \xrightarrow{p} 0
\]
uniformly in \( x_i \in \mathcal{X} \).

In the following, we take \( \hat{u}_i(\gamma_0) - u_i \xrightarrow{p} 0 \) for illustration since others are easier to show. By Lemma B.3 of Newey (1994),

\[
\sup_{x_i} \left| \frac{1}{n-1} \sum_{j=1, j \neq i}^n K_h(x_j - x_i) k_h^2(q_j - \gamma_0) y_j - \mathbb{E} \left[ y_j | x_i, \gamma_0^\pm \right] f_i(x_i, \gamma_0) \right| = O_p\left( \sqrt{\ln n/nh^2 + h} \right) = o_p\left( 1 \right),
\]

which implies

\[
\sup_{x_i} |\hat{m}_\pm(x_i, \gamma_0) - m_\pm(x_i, \gamma_0)| \xrightarrow{p} 0.
\]

As a result,

\[
\sup_{x_i, |q_i - \gamma_0| \leq h} |\hat{u}_i(\gamma_0) - u_i| = \sup_{x_i, |q_i - \gamma_0| \leq h} |y_i - \hat{m}_-(x_i, \gamma_0) 1(q_i \leq \gamma_0) - \hat{m}_+(x_i, \gamma_0) 1(q_i > \gamma_0) - u_i| = \sup_{x_i, |q_i - \gamma_0| \leq h} |m-(x_i, q_i) - \hat{m}_-(x_i, \gamma_0) 1(q_i \leq \gamma_0) + [m+(x_i, q_i) - \hat{m}_+(x_i, \gamma_0)] 1(q_i > \gamma_0)| \leq \sup_{x_i} |m_\pm(x_i, \gamma_0) - \hat{m}_\pm(x_i, \gamma_0)| + \sup_{x_i, |q_i - \gamma_0| \leq h} |m_\pm(x_i, q_i) - m_\pm(x_i, \gamma_0)| \xrightarrow{p} 0.
\]

\[\blacksquare\]

**Proposition 5**  \( \hat{\Sigma} - \hat{\Sigma}(\gamma_0) \xrightarrow{p} 0. \)

**Proof.** To derive the result, we need to show

\[
\frac{1}{n} \sum_{i=1}^n k_h(q_i - \tilde{\gamma}) \hat{\Delta}^2_i(\tilde{\gamma}) \hat{f}^2(x_i) \hat{u}_i^2 - \frac{1}{n} \sum_{i=1}^n k_h(q_i - \gamma_0) \hat{\Delta}^2_i(\gamma_0) \hat{f}^2(x_i) \hat{u}_i(\gamma_0)^2 \xrightarrow{p} 0
\]

and

\[
\frac{1}{n} \sum_{i=1}^n k_h(q_i - \tilde{\gamma}) \hat{\Delta}^2_i(\tilde{\gamma}) \hat{f}^{-1}(x_i, \tilde{\gamma}) \hat{f}(x_i) - \frac{1}{n} \sum_{i=1}^n k_h(q_i - \gamma_0) \hat{\Delta}^2_i(\gamma_0) \hat{f}^{-1}(x_i, \gamma_0) \hat{f}(x_i) \xrightarrow{p} 0.
\]

Take the first result for illustration since the second is simpler. First, we show that

\[
\frac{1}{n} \sum_{i=1}^n k_h(q_i - \gamma_0) \hat{\Delta}^2_i(\gamma_0) \hat{f}^2(x_i) \hat{u}_i^2 - \frac{1}{n} \sum_{i=1}^n k_h(q_i - \gamma_0) \hat{\Delta}^2_i(\gamma_0) \hat{f}^2(x_i) \hat{u}_i(\gamma_0)^2 \xrightarrow{p} 0. \tag{31}
\]

Since

\[
\sup_{x_i} \left| \hat{f}(x_i, \gamma) - \hat{f}(x_i, \gamma_0) \right| = \left| \frac{1}{n-1} \sum_{j=1, j \neq i}^n K_h(x_j - x_i) k_h(q_j - \gamma) - \frac{1}{n-1} \sum_{j=1, j \neq i}^n K_h(x_j - x_i) k_h(q_j - \gamma_0) \right| = \left| \frac{1}{n-1} \sum_{j=1, j \neq i}^n K_h(x_j - x_i) \left( k\left( \frac{q_j - \gamma}{h} \right) - k\left( \frac{q_j - \gamma_0}{h} \right) \right) \right| = O\left( \frac{1}{n^2 h} \right) = o_p\left( 1 \right), \tag{32}
\]

and, by a similar argument as that in (32),

\[
\sup_{x_i} \left| \hat{m}_\pm(x_i, \gamma) \hat{f}_\pm(x_i, \gamma) - \hat{m}_\pm(x_i, \gamma_0) \hat{f}_\pm(x_i, \gamma_0) \right| = o_p\left( 1 \right),
\]

we have

\[
\sup_{x_i} \left| \hat{\Delta}_i(\gamma) - \hat{\Delta}_i(\gamma_0) \right| \leq \sup_{x_i} \left| \hat{m}_-(x_i, \gamma) \hat{f}_-(x_i, \gamma) - \hat{m}_-(x_i, \gamma_0) \hat{f}_-(x_i, \gamma_0) \right| + \sup_{x_i} \left| \hat{m}_+(x_i, \gamma) \hat{f}_+(x_i, \gamma) - \hat{m}_+(x_i, \gamma_0) \hat{f}_+(x_i, \gamma_0) \right| \xrightarrow{p} 0. \tag{33}
\]
With the results in (32), (33) and 1 (\(\gamma_0 < q_i \leq \gamma\)) = \(O_p\left(\frac{1}{\sqrt{n/h}}\right) = o_p(1)\), we have

\[
\sup_{x_i} |\hat{m}_i - \hat{\mu}_i(\gamma_0)| \\
= \sup_{x_i} |\hat{m}_-(x_i, \hat{\gamma}) - \hat{m}_+(x_i, \gamma_0)| 1(q_i \leq \hat{\gamma}) + \hat{m}_+(x_i, \gamma_0) 1(q_i > \hat{\gamma}) - \hat{m}_+ (x_i, \gamma_0) 1(q_i > \gamma_0)| \\
\leq \sup_{x_i} |\hat{m}_-(x_i, \hat{\gamma}) - \hat{m}_-(x_i, \gamma_0)| 1(q_i \leq \hat{\gamma}) + \sup_{x_i} |\hat{m}_-(x_i, \gamma_0)| 1(\gamma_0 < q_i \leq \hat{\gamma}) \\
+ \sup_{x_i} |\hat{m}_+(x_i, \hat{\gamma}) - \hat{m}_+(x_i, \gamma_0)| 1(q_i > \hat{\gamma}) + \sup_{x_i} |\hat{m}_+(x_i, \gamma_0)| 1(\gamma < q_i \leq \gamma_0) \to 0.
\]

(34)

Combining (32), (33), (31) is obtained.

Secondly, we show

\[
\sup_{\gamma \in \Gamma} |\hat{\gamma} - \gamma_0| = \hat{K}_h^{\gamma} \hat{\gamma} - K_h^{\gamma} + \hat{K}_h^{\gamma} = \hat{K}_h^{\gamma} - K_h^{\gamma} - \frac{1}{h} \sum_{i=1}^n h_i^2 |\hat{\gamma} - \gamma_0| \Delta_i^2(\hat{\gamma}) \Delta_i^2(\hat{\gamma}) + o_p(1)
\]

Hence the required result is derived. ■

**Proposition 6** \(\hat{\gamma} - \gamma_0 = O_p(h)\).

**Proof.** The proof mimics that of Proposition 1. By replacing Assumption K, G and H with K', G' and H', we now have

\[
\sup_{\gamma \in \Gamma} |\rho_n^{-2} \hat{Q}_n(\gamma) - Q_n(\gamma)| = O_p\left(\frac{1}{\sqrt{n/h}}\right) \to 0,
\]

where \(Q_n(\gamma)\) contains only the middle term in Proposition 1 and the first and third terms disappear because their difference is \(O(h^{2\gamma_0^2} / n^2) = o(1)\). Now, for \(\gamma \in \Gamma \setminus N_n\), \(\sup_{\gamma \in \Gamma \setminus N_n} Q_n(\gamma) = o(1)\). For \(\gamma \in N_n\), \(Q_n(\gamma)\) is a nondecreasing (nonincreasing) function of \(a\) for \(a \in (-1, 0)\) (\(a \in (0, 1)\)) up to \(o(1)\) and \(\sup_{\gamma \in N_n} |Q_n(\gamma)| = O(1)\). So \(Q_n(\gamma)\) is maximized at some \(\gamma_n \in N_n\) such that \(Q_n(\gamma_n) > \sup_{\gamma \in \Gamma \setminus N_n} |Q_n(\gamma)| + C/2\) for \(n\) large enough. The result of interest is then derived. ■

**Proposition 7** \(\hat{\gamma} - \gamma_0 = O_p\left((n\rho_n^{-2})^{-1}\right)\).

**Proof.** This proof mimics that of Proposition 2 with the term \(\sqrt{n/h}\) replaced by \(n\rho_n^{-2}\); Suppose \(\gamma_0 < \gamma < \gamma_0 + h\), now we have

\[
D_1 = \frac{1}{n(n-1)} \sum_{i=1}^n \sum_{j=1,j \neq i}^n g_j \left(K_{h,ij} - K_{h,ij}^0\right) - \frac{1}{n(n-1)} \sum_{i=1}^n \sum_{j=1,j \neq i} g_j \left(K_{h,ij}^+ - K_{h,ij}^0\right) + \frac{1}{n(n-1)} \sum_{i=1}^n \sum_{j=1,j \neq i} g_j \left(K_{h,ij}^0 - K_{h,ij}^0\right) 1(q_j \leq \gamma_0)
\]

\[
\leq -C \rho_n \left|\frac{\gamma - \gamma_0}{h}\right|
\]

for some \(C > 0\) with probability approaching 1 by calculating the mean and variance of \(D_1\) in its U-projection, where the difference of the first two terms contribute only \(O_p(|\gamma - \gamma_0| h^2)\), the third term contributes to \(-C \rho_n \left|\frac{\gamma - \gamma_0}{h}\right|\). Since \(\rho_n / h^2 \to \infty\), for \(\eta \leq h\),

\[
P \left(\sup_{|\gamma - \gamma_0| < \eta} \frac{1}{n} \sum_{i=1}^n \left(\Delta_i(\gamma) - \Delta_i(\gamma_0)\right) > 0\right) \leq P \left(\sup_{|\gamma - \gamma_0| < \eta} D_2 > C \rho_n \left|\frac{\gamma - \gamma_0}{h}\right|\right).
\]

With a different kernel function in Assumption K' and the same formula of D2, we now have

\[
Var(P1) \leq C \frac{1}{n^h} \left(\frac{2\rho_n}{h}\right)^2, Var(P2) \leq C \frac{1}{n^h} \left(\frac{2\rho_n}{h}\right)^2 \text{ and } Var(P3) \leq C \frac{1}{n^h} \left|\frac{\gamma - \gamma_0}{h}\right|.
\]
Proof. We mimic the proof of Proposition 3. Now, \( |\hat{\Delta}_i (\gamma_0) - \Delta_f (x_i) | \xrightarrow{P} 0 \) uniformly in \( i \) and \( v \), where
\[
\Delta_f (x_i) := (1, x_i', \gamma_0) \delta_n f(x_i, \gamma_0) = O_p(\rho_n).
\]
Decompose \( \hat{\Delta}_i (\gamma_0) - \Delta_f (x_i) \) into the same six terms as in the proof of Proposition 3 only with \( \delta_0 \) replaced by \( \delta_n \). By Lemma 7, we have
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2 \Delta_f (x_i) (T_{1ij} + T_{2ij}) \approx 0,
\]
by Lemma 8, we have
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2 \Delta_f (x_i) (T_{3ij} + T_{4ij})
\approx \frac{2}{nh} \sum_{i=1}^{n} \Delta_f (x_i) f(x_i) u_i \left[ \left( k^+ \left( \frac{q_i - \gamma_0}{h} \right) - k^- \left( \frac{q_i - \gamma_0}{h} \right) \right) \right]
\]
\[
= S_2(v) / nh
\]
and by Lemma 9 we have
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2\Delta_f(x_i) (T_{ij} - T_{ii}) \\
\approx \frac{2}{nh} \sum_{i=1}^{n} \Delta_f(x_i) f(x_i) (1, x_i^T, q_i) \delta_n \left[ k^+ \left( \frac{q_i - \gamma_0}{h} \right) - k^- \left( \frac{q_i - \gamma_0}{h} \right) \right] =: S_1(v)/nh.
\]

As a result,
\[
\mathbb{E} \left[ nh \left( \hat{Q}_n(\gamma_0^v) - \bar{Q}_n(\gamma_0) \right) \right] \approx \mathbb{E} [S_2(v)] \approx 2k_+(0) f_0(\gamma_0) \mathbb{E} \left[ \Delta_f^2 f(x_i, q_i) f(x_i) | q_i = \gamma_0 \right] v/p_n^2.
\]

Combining these results and the fact that \(k_-(0) = k_+(0)\), we obtain the first equation in the proposition.

By Lemma 10 and Lemma 11 we have
\[
\text{Cov} \left( \frac{nh}{nh^2} \left( \hat{Q}_n(\gamma_0^v) - \bar{Q}_n(\gamma_0) \right) \right) = \text{Cov} \left( S_1(v_1), S_1(v_2) \right) + \text{Cov} \left( S_2(v_1), S_2(v_2) \right) + \text{Cov} \left( S_1(v_1), S_2(v_2) \right) + \text{Cov} \left( S_1(v_2), S_2(v_1) \right) \\
\approx \Sigma_n v_2.
\]

Just as in Method II, \(S_2(v)\) contributes to the mean process of \(nh \left( \hat{Q}_n(\gamma_0^v) - \bar{Q}_n(\gamma_0) \right)\), and \(S_1(v)\) contributes to the variance process.

To show fidelity convergence, we apply the Cramér-Wold device, combined with the Lyapunov CLT. Specifically, we check the Lyapunov condition that
\[
\mathbb{E} \left[ \sum_{i=1}^{n} \left( \Delta_f f(x_i) u_i \right)^4 \right] = \frac{2}{n} \mathbb{E} \left[ \Delta_f^4 f(x_i) u_i \right] = O_p \left( \frac{nh^2}{nh^2} \right) = o_p(1)
\]
as \(n \to \infty\). Then the proposition is proved. ■

**Proposition 9** \(\hat{\gamma} = \gamma_0 = O_p(h)\).

**Proof.** The proof is similar to that of Proposition 6. But now
\[
Q_n(\gamma) = \left[ \int_{1}^{0} \int K^x(u_x, x_o) k_-(u_q) m(x_o + u_x h, \gamma_0 + u_q h) f(x_o + u_x h, \gamma_0 + u_q h) du_x du_q \right]^2,
\]
and we require \(\Delta_o/h^s \to \infty\) to make the proof go through. ■

**Proposition 10** \(\hat{\gamma} = \gamma_0 = O_p((nh^{-1} \Delta_0^2)^{-1})\).

**Proof.** The proof is the same as that of Proposition 7. We only pay attention to the role that \(\Delta_o \to 0\) plays to make the proof go through. ■

**Proposition 11** On any compact set of \(v\),
\[
\mathbb{E} \left[ nh^d \left( \hat{Q}_n(\gamma_0^v) - \bar{Q}_n(\gamma_0) \right) \right] = -2k_+(0) f(x_o, \gamma_0)^2 |v| + o(v).
\]
Hence, with the zero mean assumption of the error term, we have

\[ \text{Cov} \left( nh^d \left( \tilde{Q}_n (\gamma_0^v) - \tilde{Q}_n (\gamma_0) \right), nh^d \left( \tilde{Q}_n (\gamma_0^{v_2}) - \tilde{Q}_n (\gamma_0) \right) \right) \]

\[ = \left\{ \begin{array}{ll}
\Sigma_o (v_2) v_2 + o (v_2), & \text{if } v_1 \geq v_2 \geq 0 \text{ or } v_1 \leq v_2 \leq 0, \\
0, & \text{otherwise,}
\end{array} \right. \]

and the finite-dimensional limit distributions of \( nh^d \left( \tilde{Q}_n (\gamma_0^v) - \tilde{Q}_n (\gamma_0) \right) \) are the same as those of \( C_o (v) \), where \( \gamma_0^v = \gamma_0 + \frac{v}{n h} \Sigma^v \),

\[ \Sigma_o (v) = \left\{ \begin{array}{ll}
16k_2^2 (0) \sigma_x^2 (x_0) f (x_0, \gamma_0)^3 \kappa^2, & \text{if } v \leq 0, \\
16k_2^2 (0) \sigma_x^2 (x_0) f (x_0, \gamma_0)^3 \kappa^2, & \text{if } v > 0,
\end{array} \right. \]

and \( C_o (v) \) is defined in the proof of Corollary 3.

**Proof.** Mimic the proof of Proposition 8. Now,

\[ \tilde{Q}_n (\gamma_0^v) - \tilde{Q}_n (\gamma_0) = \tilde{\Delta}_o (\gamma_0^v)^2 - \tilde{\Delta}_o (\gamma_0)^2 = \left( \tilde{\Delta}_o (\gamma_0^v) + \tilde{\Delta}_o (\gamma_0) \right) \left( \tilde{\Delta}_o (\gamma_0^v) - \tilde{\Delta}_o (\gamma_0) \right). \]

By Lemma B.1 of Newey (1994), we can show that \( \left| \tilde{\Delta}_o (\gamma_0^v) - \Delta f (x_o) \right| \xrightarrow{p} 0 \) uniformly in \( i \) and \( v \), where \( \Delta f (x_o) := (1, x_o, \gamma_0) \delta_0 f (x_o, \gamma_0) = O_p(1) \), so \( \left| \tilde{\Delta}_o (\gamma_0^v) + \tilde{\Delta}_o (\gamma_0) - 2 \Delta f (x_o) \right| \xrightarrow{p} 0 \) uniformly in \( i \) and \( v \). We then only need to focus on the other term. For simplicity, let \( v > 0 \). Now,

\[ \tilde{\Delta}_o (\gamma_0^v) - \tilde{\Delta}_o (\gamma_0) := \frac{1}{n} \sum_{j=1}^{n} (T_{1j} + T_{2j} + T_{3j} + T_{4j} + T_{5j} + T_{6j}), \]

where \( \{T_{kj}\}_{k=1} \) are defined similarly as in Proposition 8 only with \( K_{h,i,j}^{-} \) and \( K_{h,i,j}^{+} \) replaced by \( K_{h,i,j}^{-} \) and \( K_{h,i,j}^{+} \), respectively. By Lemma 13 we have

\[ \frac{1}{n} \sum_{j=1}^{n} 2 \Delta f (x_o) (T_{1j} + T_{2j}) \approx 0, \]

and by Lemma 14 we have

\[ \frac{1}{n} \sum_{j=1}^{n} 2 \Delta f (x_o) (T_{5j} + T_{6j}) \approx -2k_+ (0) \frac{v}{nh^d} f (x_0, \gamma_0)^2. \]

So

\[ nh^d \left( \tilde{Q}_n (\gamma_0^v) - \tilde{Q}_n (\gamma_0) \right) \]

\[ = h^d \sum_{j=1}^{n} 2 \Delta f (x_o) (T_{1j} + T_{2j}) + h^d \sum_{j=1}^{n} 2 \Delta f (x_o) (T_{3j} + T_{4j}) + h^d \sum_{j=1}^{n} 2 \Delta f (x_o) (T_{5j} + T_{6j}) \]

\[ = o_p (1) + S_1 (v) + S_2 (v). \]

Hence, with the zero mean assumption of the error term, we have

\[ \mathbb{E} \left[ nh^d \left( \tilde{Q}_n (\gamma_0^v) - \tilde{Q}_n (\gamma_0) \right) \right] = \mathbb{E} [S_2 (v)] = -2k_+ (0) f (x_o, \gamma_0)^2 v + o_p (v). \]

By Lemmas 13, 15 and 16 and the exogeneity property of the error terms, we have

\[ \text{Cov} \left( nh^d \left( \tilde{Q}_n (\gamma_0^v_1) - \tilde{Q}_n (\gamma_0) \right), nh^d \left( \tilde{Q}_n (\gamma_0^{v_2}) - \tilde{Q}_n (\gamma_0) \right) \right) \]

\[ = \text{Cov} \left( S_1 (v_1), S_1 (v_2) \right) + \text{Cov} \left( S_2 (v_1), S_2 (v_2) \right) + \text{Cov} \left( S_1 (v_1), S_2 (v_2) \right) + \text{Cov} \left( S_1 (v_2), S_2 (v_1) \right) \]

\[ \approx \Sigma_{ov} v_2. \]

To show the weak convergence, we apply the Lyapunov CLT by checking the Lyapunov condition. Specifically,
we show that
\[
\sum_{j=1}^{n} \mathbb{E} \left[ h^{d} \Delta_{i}^{4} (x_{o}) u^{i}_{j} \left( \left( K_{h,j}^{\gamma_{0}^{-}} - K_{h,j}^{\gamma_{0}^{-}} \right) - \left( K_{h,j}^{\gamma_{0}^{+}} - K_{h,j}^{\gamma_{0}^{+}} \right) \right)^{4} \right] = O_{p} \left( \frac{nh^{d} \Delta_{i}^{4}}{nh^{d} \Delta_{i}^{2}} \right) = o_{p} (1).
\]

\[\Box\]

**Proposition 12** \( I_{1n}^{(1)} \) is \( o_{p} (1) \) under \( H_{0}^{(1)} \), and is \( O_{p} (nh^{d/2}) \) under \( H_{1}^{(1)} \).

**Proof.** Note that
\[
I_{1n} = \frac{nh^{d/2}}{n(n-1)} \sum_{i} \sum_{j \neq i} D_{i} D_{j} K_{h,ij}
= \frac{nh^{d/2}}{n(n-1)} \sum_{i} \sum_{j \neq i} \left[ m_{i} - x_{i}^{j} \beta - x_{i}^{j} \delta_{1} (q_{i} \leq \gamma) \right] \left[ m_{j} - x_{j}^{i} \beta - x_{j}^{i} \delta_{1} (q_{j} \leq \gamma) \right] K_{h,ij}.
\]

Under \( H_{0}^{(1)} \), \( m_{i} = x_{i}^{j} \beta_{0} + x_{i}^{j} \delta_{0} (q_{i} \leq \gamma_{0}) \), so that
\[
\begin{align*}
m_{i} - x_{i}^{j} \beta - x_{i}^{j} \delta_{1} (q_{i} \leq \gamma) \\
= x_{i}^{j} \left( \beta_{0} - \beta \right) + x_{i}^{j} \left( \delta_{0} - \delta \right) 1 (q_{i} \leq \gamma \wedge \gamma_{0}) \\
+ x_{i}^{j} \delta_{0} 1 (\gamma < q_{i} \leq \gamma_{0}) - x_{i}^{j} \delta_{1} 1 (\gamma_{0} < q_{i} \leq \gamma).
\end{align*}
\]

As a result, \( I_{1n} \) has ten terms with a typical term of the form
\[
T_{1} = (\beta - \beta_{0})' \left[ \frac{nh^{d/2}}{n(n-1)} \sum_{i} \sum_{j \neq i} K_{h,ij} x_{i} x_{j}^{i} \right] (\beta - \beta_{0})
\]
or
\[
T_{2} = \delta_{0} \left[ \frac{nh^{d/2}}{n(n-1)} \sum_{i} \sum_{j \neq i} K_{h,ij} x_{i} x_{j}^{i} 1 (\gamma < q_{i} \leq \gamma_{0}) 1 (\gamma < q_{j} \leq \gamma_{0}) \right] \delta_{0}.
\]

Given that \( \beta - \beta_{0} = O_{p} (n^{-1/2}) \), \( \delta - \delta_{0} = O_{p} (n^{-1/2}) \), and \( \gamma - \gamma_{0} = O_{p} (n^{-1}) \), it is easy to show that \( T_{1} = O_{p} (h^{d/2}) \) and \( T_{2} = O_{p} (h^{d/2}) \) since \( \frac{1}{n(n-1)} \sum_{j \neq i} K_{h,ij} x_{i} x_{j}^{i} = O_{p} (1) \) and \( \frac{1}{n(n-1)} \sum_{j \neq i} K_{h,ij} x_{i} x_{j}^{i} 1 (\gamma < q_{i} \leq \gamma_{0}) 1 (\gamma < q_{j} \leq \gamma_{0}) = O_{p} (n^{-1}) \).

We now analyze \( I_{1n} \) under \( H_{1}^{(1)} \). There are three cases. Let
\[
(\beta'_{o}, \delta'_{o}, \gamma_{o})' = \arg \inf_{\beta, \delta, \gamma} \mathbb{E} \left[ (y - x' \beta - x' \delta 1 (q \leq \gamma))^{2} \right].
\]

If \( \delta_{o} = 0 \), then \( m(x, q) = x' \beta_{o} \) and the model degenerates to the case analyzed in Zheng (1996). If \( \delta_{xo} = 0 \) and \( \delta_{ao} + \gamma_{0} \delta_{qo} = 0 \), then \( m(x, q) \) takes the CTR form of Chan and Tsay (1998). It follows that \( \beta - \beta_{o} = O_{p} (n^{-1/2}) \), \( \delta - \delta_{o} = O_{p} (n^{-1/2}) \), and \( \gamma - \gamma_{o} = O_{p} (n^{-1/2}) \). If \( \delta_{xo} \neq 0 \) or \( \delta_{ao} + \gamma_{0} \delta_{qo} \neq 0 \), then \( \beta - \beta_{o} = O_{p} (n^{-1/2}) \), \( \delta - \delta_{o} = O_{p} (n^{-1/2}) \), and \( \gamma - \gamma_{o} = O_{p} (n^{-1}) \). See Yu (2017) for these results. We concentrate on the last case. Now,
\[
I_{1n} = \frac{nh^{d/2}}{n(n-1)} \sum_{i} \sum_{j \neq i} m_{i} m_{j} K_{h,ij} (1 + o_{p}(1)),
\]
where \( m_{i} = x_{i}^{j} \beta_{o} + x_{i}^{j} \delta_{o} 1 (q_{i} \leq \gamma_{o}) \), so we need only calculate \( \mathbb{E} [(m_{i} - \overline{m}_{i}) (m_{j} - \overline{m}_{j}) K_{h,ij}] \), which is equal to
\[
\begin{align*}
\int (m_{i} - \overline{m}_{i}) (m_{j} - \overline{m}_{j}) K_{h,ij} f_{i} f_{j} dx_{i} dq_{i} dx_{j} dq_{j} \\
\approx \int (m_{i} - \overline{m}_{i})^{2} K^{*}(u_{x}, x_{i}) k(u_{q}) f_{i}^{2} dx_{i} dq_{i} dx_{q} \\
= \int (m_{i} - \overline{m}_{i})^{2} f_{i}^{2} dx_{i} dq_{i},
\end{align*}
\]
The result follows. \( \Box \)
Proposition 13 \( I_{1n}^{(2)} \) is \( o_{P_m}(1) \) uniformly in \( m \) under \( H_0^{(2)} \), and is \( O_p(nh^{d/2}b) \) under \( H_1^{(2)} \).

Proof. Given that \( \hat{f}_i^{-1} = f_i^{-1} + o_p(1) \) and \( f_i \) is bounded uniformly over \( (x_i, q_i) \in \mathcal{X} \times \Gamma \),

\[
I_{1n}^{(2)} = \frac{nh^{d/2}}{m(n-1)^2} \sum_i \sum_{j 
eq i} 1_i 1_j (m_i - \hat{m}_i) \hat{f}_i (m_j - \hat{m}_j) \hat{f}_j K_{h,ij} \left( \hat{f}_i^{-1} \hat{f}_j^{-1} \right) 
= O_p \left( \frac{nh^{d/2}}{m(n-1)^2} \sum_i \sum_{j 
eq i} 1_i 1_j (m_i - \hat{m}_i) L_{b,i} f_i^{-1} (m_j - \hat{m}_j) L_{b,j} f_j^{-1} \right)
\]

(35)

Mimicking the proof of Proposition A.1 of Fan and Li (1996), we can show that under \( H_0^{(2)} \), \( I_{1n} = O_p \left( nh^{d/2}b^2 \right) = o_p(1) \). The only new result we need to employ is that \( \| \mathbb{E}_1 [(m_2 - m_1)L_{b,21}] \| = O_p(b^3) \), which is accomplished in Lemma 18.

We now analyze \( I_{1n} \) under \( H_1^{(2)} \). It can be shown that the case where \( i, j, l, k \) are all different from each other dominates in the formula of the second equality of (35), so

\[
I_{1n} \approx O_p \left( nh^{d/2}E \left[ 1_1^1 (m_1 - m_2) L_{b,12} f_1^{-1} (m_3 - m_4) L_{b,34} K_{h,13} f_3^{-1} f_4^{-1} \right] \right).
\]

Because \( h/b \to 0 \), we can treat \( (x_1, q_1) = (x_3, q_3) \). Specifically,

\[
\mathbb{E} \left[ 1_1^1 (m_1 - m_2) L_{b,12} f_1^{-1} (m_3 - m_4) L_{b,34} K_{h,13} f_3^{-1} \right] 
= \mathbb{E} \left[ \frac{1}{b^2} L^x \left( \frac{x_1 - x_2}{b}, x_1 + u_q h \right) l ( \frac{q_1 - q_2}{b} ) K^x (u_x, x_1) k (u_q) \right] 
\approx \mathbb{E} \left[ 1_1^1 (m_1 - m_2) L_{b,12} f_1^{-1} \right] 
= \mathbb{E} \left[ \frac{1}{b^2} \right] \left\{ \mathbb{E}_1 \left[ (m_1 - m_2) L_{b,12} \right] \right\} 
\approx O \left( b^2 \right)
\]

(34)

\[f_1^{\gamma_0} \int_{\gamma_0^-}^{\gamma_0^+} \left[ f \left( m(x_1, q_1) - m(x_1, q_1 + u_q b) l (u_q) \right) \right] f(x_1, q_0) dx_1 dq_1
\approx O \left( b^2 \right) + \int_{\gamma_0^-}^{\gamma_0^+} \int_{\gamma_0^-}^{\gamma_0^+} \left[ f \left( m(x_1, q_1) - m(x_1, q_1 + u_q b) l (u_q) \right) \right] f(x_1, q_0) dx_1 dq_1
\]

(35)

where \( u = (u_x, u_q) \). Under \( H_1^{(2)} \),

\[f_1^{\gamma_0} \int_{\gamma_0^-}^{\gamma_0^+} \left[ f \left( m(x_1, q_1) - m(x_1, q_1 + u_q b) l (u_q) \right) \right] f(x_1, q_0) dx_1 dq_1
\approx O \left( b^2 \right) + \int_{\gamma_0^-}^{\gamma_0^+} \int_{\gamma_0^-}^{\gamma_0^+} \left[ f \left( m(x_1, q_1) - m(x_1, q_1 + u_q b) l (u_q) \right) \right] f(x_1, q_0) dx_1 dq_1
\]

(36)

The result follows.
Proposition 14 \( I_{2n}^{(1)} \xrightarrow{d} N(0, \Sigma^{(1)}) \) and \( I_{2n}^{(2)} \xrightarrow{d} N(0, \Sigma^{(2)}) \).

Proof. We only prove the second result since the first is proved in a similar way.

\[
I_{2n}^{(2)} = \frac{n h^{d/2}}{n(n-1)} \sum_i \sum_{j \neq i} 1_{t_1}^t 1_{t_2}^t u_i u_j K_{h,i,j}
\]

\[
\equiv \frac{n h^{d/2}}{n(n-1)} \sum_i \sum_{j \neq i} H_n(z_i, z_j) \equiv n h^{d/2} U_n,
\]

where \( U_n \) is a second order degenerate U-statistic with kernel function \( H_n \). We can apply theorem 1 of Hall (1984) to find its asymptotic distribution. Two conditions need to be checked: (i) \( \mathbb{E}[H_n^2(z_1, z_2)] < \infty \); (ii)

\[
\frac{1}{n} \mathbb{E}[G_n^2(z_1, z_2) + n^{-1} \mathbb{E}[H_n^2(z_1, z_2)]] \to 0 \text{ as } n \to \infty,
\]

where \( G_n(z_1, z_2) = \mathbb{E}[H_n(z_3, z_1)H_n(z_3, z_2)[z_1, z_2] \]. Because these checks follow in a similar way to lemma 3.3a of Zheng (1996) they are omitted here to save space. In conclusion

\[
n U_n / \sqrt{2 \mathbb{E}[H_n^2(z_1, z_2)]} \xrightarrow{d} N(0, 1).
\]

It is easy to check that

\[
\mathbb{E} \left[ H_n^2(z_1, z_2) \right] = \mathbb{E} \left[ \frac{1}{n(n-1)} \sum_i \sum_{j \neq i} K_{h,i,j}^2 \right]
\]

\[
= \frac{1}{n(n-1)} \sum_i \sum_{j \neq i} \mathbb{E} \left[ K_{h,i,j}^2 \right] f(x_1, q_1) f(x_2, q_2) dx_1 dq_1
\]

\[
\approx \frac{1}{n(n-1)} \sum_i \sum_{j \neq i} \mathbb{E} \left[ K_{h,i,j}^2 \right] f(x_1, q_1) f(x_2, q_2) dx_1 dq_1
\]

\[
= \mathbb{E} \left[ H_n^2(z_1, z_2) \right] \xrightarrow{d} N(0, 1).
\]

Proposition 15 \( v_n^{(1)} = \Sigma^{(1)} + o_p(1) \) under \( H_0^{(1)} \) and the local alternative and \( v_n^{(1)2} = O_p(1) \) under \( H_1^{(1)} \).

Proof. It can be shown that

\[
\frac{k^d}{n(n-1)} \sum_i \sum_{j \neq i} K_{h,i,j}^2 \mathbb{E} \left[ \sigma_i^2 - \sigma_j^2 \right] = O \left( n^{-2} h^{-d} \right) - o(1).
\]

where \( \sigma_i^2 = \sigma^2(x_i, q_i) \). Under \( H_0^{(1)} \), \( m - m = 0 \). Under the local alternative, \( \mathbb{E} \left[ f(x, q) (m - m)^4 \right] = O \left( n^{-2} h^{-d} \right) = o(1) \), and under \( H_1^{(1)} \), \( \mathbb{E} \left[ f(x, q) (m - m)^4 \right] = O(1) \).

Proposition 16 \( v_n^{(2)} = \Sigma^{(2)} + o_p(1) \) under both \( H_0^{(2)} \) and \( H_1^{(2)} \).
Proof. By steps similar to the last proposition, we have

$$\frac{k^{d'}}{n^{(n-1)}} \sum_i \sum_{j \neq i} 1_i^T 1_j^T K_{i,j}^2 \epsilon_i^2 \epsilon_j^2 = \int k^{2d}(u) du \mathbb{E} \left[ 1_i^T f(x_i, q_i) \left( \sigma^2(x_i, q_i) + (m_i - \overline{m}_i)^2 \right)^2 \right] + o_p(1),$$

where \( \overline{m}_i \) is redefined as \( \mathbb{E}_i \{ \overline{m}_i \} \). Note that \( \mathbb{E} \left[ 1_i^T f(x_i, q_i) (m_i - \overline{m}_i)^4 \right] \) is at most \( O(b) \) since \( m_i - \overline{m}_i \) contributes only for \( q \in [\gamma - b, \gamma + b] \).

Supplement C: Lemmas

Lemma 1 \( \frac{1}{n^{(n-1)}} \sum_i \sum_{j=1, j \neq i}^{n} 2 \Delta f(x_i) (T_{1ij} + T_{2ij}) \approx 0 \).

Proof. For \( 2 \Delta f(x_j) T_{1ij} \), we have

$$\mathbb{E} \left[ 2 \Delta f(x_j) T_{1ij} \mid z_j \right] = \int_{x_j} 2 \Delta f(x_j) g(x_j, q_j) \frac{1}{n} K^x \left( \frac{x_j - x_i}{h} \right) \left( k - \left( \frac{q_j - \gamma_0}{h} \right) \right) f(x_j) dx_j d q_j$$

$$= 2 \Delta f(x_j) \int_{x_j} g(x_j, u_x, \gamma_0 + u_x h) K^x (u_x, x_j) k^{-1} (q_j - \gamma_0 + u_x h) du_x dx_j$$

$$= 2 \Delta f(x_j) \int_{x_j} g(x_j, u_x, \gamma_0 + u_x h) K^x (u_x, x_j) (u_x + \frac{1}{2} h) du_x dx_j$$

$$= 2 \mathbb{E} \left[ 2 \Delta f(x_j) T_{1ij} \mid z_j \right] \approx 0.$$

And

$$\mathbb{E} \left[ 4 \Delta^2 f(x_j) T_{1ij} \right] = \int_{x_j} \int_{x_j} \int_{x_j} \int_{x_j} K^x \left( \frac{x_j - x_i}{h} \right)^2 \left( k - \left( \frac{q_j - \gamma_0}{h} \right) \right)^2 f(x_j, q_j) dx_j d q_j f(x_i) dx_i$$

$$= \int_{x_j} \int_{x_j} \int_{x_j} \int_{x_j} \left( k - \left( \frac{q_j - \gamma_0}{h} \right) \right)^2 f(x_j, q_j) dx_j d q_j$$

$$= O \left( \frac{1}{h^{2n}} \frac{\overline{m}}{\sqrt{n}/h} \right),$$

we have

$$\frac{1}{n} \sum_{i=1}^{n} \mathbb{E} \left[ 2 \Delta f(x_i) T_{1ij} \mid z_i \right] = 2 \Delta f(x_i) \mathbb{G}_2 (x_i, \gamma_0) \frac{\overline{m}}{\sqrt{n}/h} + O \left( \frac{1}{h^{2n}} \frac{\overline{m}}{\sqrt{n}/h} \right),$$

$$\frac{1}{n} \sum_{i=1}^{n} \mathbb{E} \left[ 2 \Delta f(x_i) T_{1ij} \mid z_j \right] = 2 \mathbb{G} (x_j, q_j) \left( k - \left( \frac{q_j - \gamma_0}{h} \right) \right) \Delta f(x_j) f(x_j),$$

$$\frac{1}{n} \mathbb{E} \left[ 4 \Delta^2 f(x_i) T_{1ij} \right] \approx \frac{1}{n} O \left( \frac{1}{h^{2n}} \frac{1}{(nh)^{1/2}} \right) = o (1).$$
By Lemma 8.4 of Newey and McFadden (1994), we have

$$\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \mathbb{E}[2\Delta f(x_i)T_{1ij}] \approx 2\Delta f(x_i) g_{2}(x_i, \gamma_0) \frac{v}{\sqrt{n/h}},$$

where the extra terms \(\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta f(x_i)T_{1ij}] \approx \mathbb{E}[2\Delta f(x_i)T_{1ij}]\) and \(\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta f(x_i)T_{1ij}] \approx \mathbb{E}[2\Delta f(x_i)T_{1ij}]\) are \(o_p(1)\). Similarly for \(2\Delta f(x_i)T_{2ij}\), we have

$$\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \mathbb{E}[2\Delta f(x_i)T_{2ij}] \approx -2\Delta f(x_i) g_{2}(x_i, \gamma_0) \frac{v}{\sqrt{n/h}}.$$

Hence,

$$\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \mathbb{E}[2\Delta f(x_i)(T_{1ij} + T_{2ij})] \approx 0$$

and

$$\text{Var} \left( \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \mathbb{E}[2\Delta f(x_i)(T_{1ij} + T_{2ij})] \right) \leq \text{Var} \left( \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \mathbb{E}[2\Delta f(x_i)T_{1ij}] \right) + \text{Var} \left( \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} \mathbb{E}[2\Delta f(x_i)T_{2ij}] \right) = o(1).$$

So the result of interest is derived.

\[ \boxed{\text{Lemma 2}} \]

$$\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2\Delta f(x_i) (T_{3ij} + T_{4ij}) \approx \frac{1}{nh} \sum_{i=1}^{n} \Delta f(x_i) f(x_i) u_i \left( \left( k^+ \left( \frac{q_i - \gamma_0}{h} \right) \right) - \left( k^- \left( \frac{q_i - \gamma_0}{h} \right) \right) \right).$$

**Proof.** Since

$$\mathbb{E}[2\Delta f(x_i)T_{3ij}|z_i] = 0,$$

$$\mathbb{E}[2\Delta f(x_i)T_{3ij}] = \frac{2}{nh} \int_{x_i}^{x_i+h} \Delta f(x) \left( k^+ \left( \frac{q_i - \gamma_0}{h} \right) \right) - \left( k^- \left( \frac{q_i - \gamma_0}{h} \right) \right) f(x) \, dx,$$

and

$$\mathbb{E} \left[ 4\Delta^2 f(x_i)T_{3ij} \right] = \mathbb{E} \left[ 4\Delta^2 f(x_i)T_{3ij} \right] = \frac{2}{nh} \int_{x_i}^{x_i+h} \Delta f(x) \left( k^+ \left( \frac{q_i - \gamma_0}{h} \right) \right) - \left( k^- \left( \frac{q_i - \gamma_0}{h} \right) \right) f(x) \, dx,$$

we have

$$\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta f(x_i)T_{3ij}] \approx \frac{2}{nh} \sum_{i=1}^{n} u_i \left( k^- \left( \frac{q_i - \gamma_0}{h} \right) \right).$$
By Lemma 8.4 of Newey and McFadden (1994), we have
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2\Delta_f (x_i) T_{3ij} \approx \frac{2}{nh} \sum_{i=1}^{n} \Delta_f (x_i) u_i \left( k^\gamma \left( \frac{q_i - \gamma a}{h} \right) - k^\gamma \left( \frac{q_i - \gamma a}{h} \right) \right).
\]
A similar result can be derived for \( \frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2\Delta_f (x_i) T_{4ij} \) and the result of interest is then obtained.

**Lemma 3**
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2\Delta_f (x_i) (T_{5ij} + T_{6ij}) \approx \frac{2}{nh} \sum_{i=1}^{n} \Delta_f (x_i) f (x_i) \left( 1, x_i, q_i \right) \delta_0 \left( k^\gamma \left( \frac{q_i - \gamma a}{h} \right) 1(\gamma_0 - h \leq q_i \leq \gamma_0) - k^\gamma \left( \frac{q_i - \gamma a}{h} \right) 1(\gamma_0 - h \leq q_i \leq \gamma_0) \right).
\]

**Proof.** Taking \( T_5 \) to illustrate, we have
\[
\mathbb{E}[2\Delta_f (x_i) T_{5ij} | z_i] = -\frac{2}{nh} \int_{x_i} \Delta_f (x_i) \left( 1, x_i, q_i \right) \delta_0 K^\gamma \left( \frac{x-x_i}{h}, x_i \right) k^\gamma \left( \frac{q_i - \gamma a}{h} \right) 1(\gamma_0 - h \leq q_i \leq \gamma_0) f (x_i) dx_i
\]
\[
= -\frac{2}{nh} \left( 1, x_i, q_i \right) \delta_0 k^\gamma \left( \frac{q_i - \gamma a}{h} \right) 1(\gamma_0 - h \leq q_i \leq \gamma_0) \int_{x_i} \Delta_f (x_i, x_i, x_i) K^\gamma (u, x_i, x_i) f (x_i) dx_i
\]
\[
\approx -\frac{2}{nh} \left( 1, x_i, q_i \right) \delta_0 k^\gamma \left( \frac{q_i - \gamma a}{h} \right) 1(\gamma_0 - h \leq q_i \leq \gamma_0) \Delta_f (x_i) f (x_i),
\]
and
\[
\mathbb{E}
\left[
\frac{4}{n} \Delta_f^2 (x_i) T_{5ij}
\right]
\]
\[
= \int_{x_i} \int_{x_i} \int_{x_i} \frac{1}{n} \Delta_f^2 (x_i) \left[ \left( 1, x_i, q_i \right) \delta_0 \right]^2 K^\gamma \left( \frac{x-x_i}{h}, x_i \right) \left( k^\gamma \left( \frac{q_i - \gamma a}{h} \right) \right)^2 f (x_i) dx_i dx_j dx_k
\]
\[
= \int_{x_i} \int_{0}^{n} \int_{x_i} \left\{ \frac{4}{n} \Delta_f^2 (x_i) \left[ \left( 1, x_i, q_i \right) \delta_0 \right]^2 \right\} dx_i dx_j dx_k
\]
\[
= O \left( \frac{1}{n^2} \right).
\]

Hence,
\[
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta_f (x_i) T_{5ij} | z_i] \approx -2\Delta_f (x_i)^2,
\]
\[
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta_f (x_i) T_{5ij}] \approx -\frac{2}{n^2} \left( 1, x_i, q_i \right) \delta_0 k^\gamma \left( \frac{q_i - \gamma a}{h} \right) 1(\gamma_0 - h \leq q_i \leq \gamma_0) \Delta_f (x_i) f (x_i),
\]
\[
\frac{1}{n} \mathbb{E}
\left[
\frac{4}{n} \Delta_f^2 (x_i) T_{5ij}^2
\right]^{1/2}
\]
\[
= \frac{1}{n} O \left( \frac{1}{n^2} \right) = o(1).
\]
A similar result can be derived for \( \frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta_f (x_i) T_{6ij} | z_i] \). Then by Lemma 8.4 of Newey and McFadden (1994), we have
\[
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta_f (x_i) (T_{5ij} + T_{6ij})]
\]
\[
\approx \frac{2}{nh} \sum_{i=1}^{n} \left( 1, x_i, q_i \right) \delta_0 \left[ k^\gamma \left( \frac{q_i - \gamma a}{h} \right) 1(\gamma_0 - h \leq q_i \leq \gamma_0) - k^\gamma \left( \frac{q_i - \gamma a}{h} \right) 1(\gamma_0 - h \leq q_i \leq \gamma_0) \right] \Delta_f (x_i) f (x_i),
\]
where the extra terms \( \frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta_f (x_i) T_{5ij} | z_i] - \mathbb{E}[2\Delta_f (x_i) T_{5ij}] \) and \( \frac{1}{n} \sum_{i=1}^{n} \mathbb{E}[2\Delta_f (x_i) T_{6ij}] - \mathbb{E}[2\Delta_f (x_i) T_{6ij}] \) are \( o_p(1) \).
Lemma 4 Cov $(S_1(v_1), S_1(v_2)) \approx 4 \mathbb{E} \left[ \Delta^2_f(x_i)f(x_i)(\sigma^2_+(x_i) + \sigma^2_-(x_i)) | q_i = \gamma_0 \right] f_\gamma(\gamma_0)v_1v_2 \xi(1)$.

Proof. Without loss of generality, assume $v_1 \geq v_2 \geq 0$. Then

$$4n \text{Cov} \left( \Delta_f(x_i) f(x_i) u_i \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right), \Delta_f(x_i) f(x_i) u_i \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right) \right)$$

$$= 4n \mathbb{E} \left[ \Delta^2_f(x_i) f(x_i)^2 u_i^2 \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right)^2 \right].$$

$$= 4n \int_{x_1}^{x_2} \int_{x_1}^{x_2} \int_{x_1}^{x_2} \Delta^2_f(x_i) f(x_i)^2 u_i^2 \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right)^2 \right] f(u_i|x_i, q_i) f(x_i, q_i) du_i dq_i dx_i$$

$$= 4n \int_{x_1}^{x_2} \int_{x_1}^{x_2} \int_{x_1}^{x_2} \Delta^2_f(x_i) f(x_i)^2 u_i^2 \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right)^2 \right] f(u_i|x_i, q_i) f(x_i, q_i) du_i dq_i dx_i$$

$$= 4n \int_{x_1}^{x_2} \int_{x_1}^{x_2} \int_{x_1}^{x_2} \Delta^2_f(x_i) f(x_i)^2 u_i^2 \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right)^2 \right] f(u_i|x_i, q_i) f(x_i, q_i) du_i dq_i dx_i$$

$$= 4n \mathbb{E} \left[ \Delta^2_f(x_i) f(x_i)^2 \sigma^2_+(x_i) | q_i = \gamma_0 \right] f_\gamma(\gamma_0)v_1v_2 \xi(1),$$

where $\sigma^2_+(x_i) = \mathbb{E}[u_i^2|x_i, q_i = \gamma_0+]$ and $\xi(1) = \int_0^1 k'_+(t)^2 dt$. For a more detailed proof, we refer to that of Lemma C.4 of DH. Similarly,

$$4n \text{Cov} \left( \Delta_f(x_i) f(x_i) u_i \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right), \Delta_f(x_i) f(x_i) u_i \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right) \right)$$

$$\approx 4 \mathbb{E} \left[ \Delta^2_f(x_i) f(x_i)^2 \sigma^2_-(x_i) | q_i = \gamma_0 \right] f_\gamma(\gamma_0)v_1v_2 \xi(1)$$

and

$$4n \text{Cov} \left( \Delta_f(x_i) f(x_i) u_i \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right), \Delta_f(x_i) f(x_i) u_i \left( k^+ \left( \frac{q_i - \gamma_0^+}{h} \right) - k^- \left( \frac{q_i - \gamma_0^-}{h} \right) \right) \right)$$

$$\approx 4 \mathbb{E} \left[ \Delta^2_f(x_i) f(x_i)^2 \sigma^2_-(x_i) | q_i = \gamma_0 \right] f_\gamma(\gamma_0)v_1v_2 \xi(1) \frac{u_i}{\sqrt{nh}} = o_P(1),$$

where $\sigma^2_-(x_i) = \mathbb{E}[u_i^2|x_i, q_i = \gamma_0-]$. So the result of interest is obtained by summing up terms. ■

Lemma 5 Cov $(S_2(v_1), S_2(v_2)) = o(v_1v_2), Cov (S_1(v_1), S_2(v_2)) = o(v_1v_2)$ and Cov $(S_2(v_1), S_1(v_2)) = o(v_1v_2)$.

Proof. From Lemma 3 by tedious calculation, we can obtain

$$\text{Var} (S_2 (v_2)) = o(v_2).$$

Hence,

$$\text{Cov} (S_2(v_1), S_2(v_2)) \leq \sqrt{\text{Var} (S_2(v_1)) \text{Var} (S_2(v_2))} = o(v_1v_2),$$

$$\text{Cov} (S_1(v_1), S_2(v_2)) \leq \sqrt{\text{Var} (S_1(v_1)) \text{Var} (S_2(v_2))} = o(v_1v_2),$$

$$\text{Cov} (S_2(v_1), S_1(v_2)) \leq \sqrt{\text{Var} (S_2(v_1)) \text{Var} (S_1(v_2))} = o(v_1v_2).$$

■

Lemma 6 For any $\phi_1, \phi_2 > 0$, there exists $\eta > 0$ such that

$$P \left\{ \sup_{v_1, v_2 < \eta} \left| \frac{1}{nh} \left( \hat{Q}_n (\gamma_0^+ - \hat{Q}_n (\gamma_0)) - \frac{1}{nh} \left( \hat{Q}_n (\gamma_0^- - \hat{Q}_n (\gamma_0)) \right) \right) > \phi_1 \right\} < \phi_2,$$

where

$$\frac{1}{nh} \left( \hat{Q}_n (\gamma_0^+ - \hat{Q}_n (\gamma_0)) - \frac{1}{nh} \left( \hat{Q}_n (\gamma_0^- - \hat{Q}_n (\gamma_0)) \right) \right) = \mathbb{E} \left[ \frac{1}{nh} \left( \hat{Q}_n (\gamma_0^+ - \hat{Q}_n (\gamma_0)) \right) \right].$$
Proof. From Proposition 3 we have

$$\lim_{n \to \infty} \mathbb{E} \left[ \hat{n} \left( \hat{Q}_n \left( \gamma_0^u \right) - \hat{Q}_n \left( \gamma_0 \right) \right) \right] = \lim_{n \to \infty} \mathbb{E} \left[ \hat{n} \left( \hat{Q}_n \left( \gamma_0^u \right) - \hat{Q}_n \left( \gamma_0 \right) \right) \right]^2 + \lim_{n \to \infty} \mathbb{E} \left[ \hat{n} \left( \hat{Q}_n \left( \gamma_0^{u^2} \right) - \hat{Q}_n \left( \gamma_0 \right) \right) \right]^2 \leq C |v_1 - v_2|^2.$$

By Markov's inequality, the result follows. ■

**Lemma 7** $\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2 \Delta f(x_i)(T_{1ij} + T_{2ij}) \approx 0$.

**Proof.** The proof is similar to that of Lemma 4. For $2 \Delta f(x_i)T_{1ij}$, we now have

$$\mathbb{E} [2 \Delta f(x_i)T_{1ij} | X_i] = \int_{-\infty}^{\infty} \int_{x_i} f(x, x_i) \frac{1}{n} K(x - x_i) \left( k^{- \left( \frac{q_i - \gamma_0^u}{h} \right)} - k^{- \left( \frac{q_i - \gamma_0}{h} \right)} \right) f(x, x_i) dx_i dx_j$$

where $g(x, q_j) = g(x, q_j) f(x, q_j), g_1^{(n)}(x, q_j) = \frac{\partial g(x, q_j)}{\partial q_j}$ and $g_2^{(n)}(x, q_j) = \frac{\partial^2 g(x, q_j)}{\partial q_j^2}$. Since

$$\mathbb{E} [2 \Delta f(x_i)T_{1ij} | X_i] = \int_{x_i} 2 \Delta f(x_i) \left( k^{- \left( \frac{q_i - \gamma_0^u}{h} \right)} - k^{- \left( \frac{q_i - \gamma_0}{h} \right)} \right) f(x_i) dx_i$$

and

$$\mathbb{E} \left[ \Delta f(x_i)T_{2ij} \right] = \int_{x_i} \int_{x_i} \int_{x_j} \frac{1}{n} \Delta f(x_i) \left( k^{- \left( \frac{q_i - \gamma_0^u}{h} \right)} - k^{- \left( \frac{q_i - \gamma_0}{h} \right)} \right)^2 f(x_i) dx_i dx_j dx_i.$$
we have
\[
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left[2 \Delta f(x_i) T_{1ij} | X_i \right] \approx 2 \Delta f(x_i) \frac{n}{nh^2} \left[ g_2(x_i, \gamma_0) + \cdots + h^s g_2^{(s)}(x_i, \gamma_0) \right],
\]
\[
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left[2 \Delta f(x_i) T_{1ij} | X_j \right] = \frac{2}{n} g(x_j, q_j) \left(k- \left(\frac{q-\gamma_0}{h}\right) - k- \left(\frac{q-\gamma_0}{h}\right)\right) \Delta f(x_j) f(x_j),
\]
\[
\frac{1}{n} \mathbb{E}\left[4 \Delta f(x_i) T_{1ij}^2 \right]^{1/2} = \frac{1}{n} O \left(\frac{1}{\sqrt{nh^2}}\right) = o(1).
\]
By Lemma 8.4 of Newey and McFadden(1994), we have
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2 \Delta f(x_i) T_{1ij} \approx 2 \Delta f(x_i) \frac{n}{nh^2} \left[ g_2(x_i, \gamma_0) + \cdots + h^s g_2^{(s)}(x_i, \gamma_0) \right],
\]
where the extra terms \( \frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left[2 \Delta f(x_i) T_{1ij} | X_j \right] - \mathbb{E}\left[2 \Delta f(x_i) T_{1ij} \right] \) and \( \frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left[2 \Delta f(x_i) T_{1ij} | X_j \right] - \mathbb{E}\left[2 \Delta f(x_i) T_{1ij} \right] \) are \( o_p(1) \). Similarly for \( 2 \Delta f(x_i) T_{2ij} \), we have
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2 \Delta f(x_i) T_{2ij} \approx -2 \Delta f(x_i) \frac{n}{nh^2} \left[ g_2(x_i, \gamma_0) + \cdots + h^s g_2^{(s)}(x_i, \gamma_0) \right].
\]
Hence,
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2 \Delta f(x_i) (T_{1ij} + T_{2ij}) \approx 0
\]
and the result of interest is derived.

Lemma 8
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2 \Delta f(x_i) \left( T_{3ij} + T_{4ij} \right)
\]
\[
\approx \frac{2}{nh} \sum_{i=1}^{n} \Delta f(x_i) f(x_i) u_i \left[ \left(k- \left(\frac{q-\gamma_0}{h}\right) - k- \left(\frac{q-\gamma_0}{h}\right)\right) - \left(k+ \left(\frac{q-\gamma_0}{h}\right) - k+ \left(\frac{q-\gamma_0}{h}\right)\right)\right].
\]
Proof. The proof is the same as that of Lemma 2.

Lemma 9
\[
\frac{1}{n(n-1)} \sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n} 2 \Delta f(x_i) \left( T_{5ij} + T_{6ij} \right)
\]
\[
\approx \frac{2}{nh} \sum_{i=1}^{n} \Delta f(x_i) f(x_i) \left( 1, x_i, q_i, \delta_n \right) \left[ \left(k- \left(\frac{q-\gamma_0}{h}\right) \right) \left(1 \leq q_i \leq \gamma_0 \right) - \left(k- \left(\frac{q-\gamma_0}{h}\right) \right) \left(0 \leq q_i \leq \gamma_0 \right) \right].
\]
Proof. The proof is the same as that of Lemma 3.

Lemma 10
\[
\text{Cov} \left( S_1(v_1), S_1(v_2) \right) \approx \Sigma_n v_2.
\]
Proof. Without loss of generality, assume \( v_1 \geq v_2 \geq 0 \). Then
\[
4n \text{Cov} \left( \Delta f(x_i) f(x_i) u_i \left(k+ \left(\frac{q-\gamma_0}{h}\right) \right) - k+ \left(\frac{q-\gamma_0}{h}\right) \right), \Delta f(x_i) f(x_i) u_i \left(k+ \left(\frac{q-\gamma_0}{h}\right) \right)
\]
\[
= 4n \mathbb{E}\left[ \left(2 \Delta f(x_i) f(x_i)^2 u_i^2 \left(k+ \left(\frac{q-\gamma_0}{h}\right) \right) \right) \right] + 4n \mathbb{E}\left[ \left(2 \Delta f(x_i) f(x_i)^2 u_i^2 \left(k+ \left(\frac{q-\gamma_0}{h}\right) \right) \right) \right] f(u_i | x_i, q_i) f(x_i, q_i) du_i dq_i dx_i
\]
\[
= 4n \int_{x_i}^{\gamma_0} \int_{u_i}^{\gamma_0} \int_{q_i}^{\gamma_0} \Delta f(x_i) f(x_i)^2 u_i^2 \left(k+ \left(\frac{q-\gamma_0}{h}\right) \right) f(u_i | x_i, q_i) f(x_i, q_i) du_i dq_i dx_i
\]
\[
+ 4n \int_{x_i}^{\gamma_0} \int_{u_i}^{\gamma_0} \int_{q_i}^{\gamma_0} \Delta f(x_i) f(x_i)^2 u_i^2 \left(k+ \left(\frac{q-\gamma_0}{h}\right) \right) f(u_i | x_i, q_i) f(x_i, q_i) du_i dq_i dx_i
\]
\[
+ 4n \int_{x_i}^{\gamma_0} \int_{u_i}^{\gamma_0} \int_{q_i}^{\gamma_0} \Delta f(x_i) f(x_i)^2 u_i^2 \left(k+ \left(\frac{q-\gamma_0}{h}\right) \right) f(u_i | x_i, q_i) f(x_i, q_i) du_i dq_i dx_i
\]
\[
\approx 4k \left(0 \right)^2 \mathbb{E}\left[ \Delta f(x_i) f(x_i)^2 \sigma_+^2(x_i) | q_i = \gamma_0 \right] f_q(\gamma_0) / \rho_n^2 v_2,
\]
31
where \( \sigma^2_i(x_i) = \mathbb{E}[u_i^2|x_i, q_i = \gamma_0^+]. \) Similarly,

\[
4n \text{Cov} \left[ \Delta_f(x_i) f(x_i) u_i \left( k^- \left( \frac{2 \gamma_0^+}{h} \right) - k^- \left( \frac{2 \gamma_0^-}{h} \right) \right) , \Delta_f(x_i) f(x_i) u_i \left( k^- \left( \frac{2 \gamma^+_i}{h} \right) - k^- \left( \frac{2 \gamma^-_i}{h} \right) \right) \right] \\
\approx 4k_+ (0)^2 \mathbb{E} \left[ \Delta_f(x_i) f(x_i)^2 \sigma^2_i(x_i) | q_i = \gamma_0 \right] f_q(\gamma_0)/\rho_n^2 v_2.
\]

and

\[
4n \text{Cov} \left[ \Delta_f(x_i) f(x_i) u_i \left( k^\pm \left( \frac{2 \gamma^+_i}{h} \right) - k^\pm \left( \frac{2 \gamma^-_i}{h} \right) \right) , \Delta_f(x_i) f(x_i) u_i \left( k^\mp \left( \frac{2 \gamma^+_i}{h} \right) - k^\mp \left( \frac{2 \gamma^-_i}{h} \right) \right) \right] \\
\approx 4k_+ (0)^2 \mathbb{E} \left[ \Delta_f(x_i) f(x_i)^2 \sigma^2_i(x_i) | q_i = \gamma_0 \right] f_q(\gamma_0)/\rho_n^2 v_2.
\]

If \( v_1 \leq v_2 \leq 0, \) the result is similar except that the term \( \sigma^2_i(x_i) \) is replaced by \( \sigma^2(x_i). \)

If \( v_1v_2 < 0, \) then the four terms are all \( o(v_2). \)

Lemma 11 Cov \( (S_2(v_1), S_2(v_2)) = o(v_2), \) Cov \( (S_1(v_1), S_2(v_2)) = o(v_2) \) and Cov \( (S_2(v_1), S_1(v_2)) = o(v_2). \)

Proof. The proof idea is the same as that in Lemma 5.

Lemma 12 For any \( \phi_1, \phi_2 > 0, \) there exists \( \eta > 0 \) such that

\[
P \left\{ \sup_{|v_1 - v_2| < \eta} \left| \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_1)}) - \hat{Q}_n(\gamma_0) \right) - \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_2)}) - \hat{Q}_n(\gamma_0) \right) \right| > \phi \right\} < \phi_2,
\]

where

\[
\sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_i)}) - \hat{Q}_n(\gamma_0) \right) = n \left( \hat{Q}_n(\gamma_0^{(v_i)}) - \hat{Q}_n(\gamma_0) \right) - \mathbb{E} \left[ \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_i)}) - \hat{Q}_n(\gamma_0) \right) \right].
\]

Proof. Without loss of generality, assume \( v_1 \geq v_2 \geq 0. \) From Proposition 8 we have

\[
\lim_{n \to \infty} \mathbb{E} \left[ \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_i)}) - \hat{Q}_n(\gamma_0) \right) - \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_2)}) - \hat{Q}_n(\gamma_0) \right) \right]^2 \\
= \lim_{n \to \infty} \mathbb{E} \left[ \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_i)}) - \hat{Q}_n(\gamma_0) \right)^2 \right] + \lim_{n \to \infty} \mathbb{E} \left[ \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_2)}) - \hat{Q}_n(\gamma_0) \right)^2 \right] \\
- 2 \lim_{n \to \infty} \mathbb{E} \left[ \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_i)}) - \hat{Q}_n(\gamma_0) \right) \right] \mathbb{E} \left[ \sqrt{n} \hat{h} \left( \hat{Q}_n(\gamma_0^{(v_2)}) - \hat{Q}_n(\gamma_0) \right) \right] \\
\leq 16k_+ (0)^2 (v_1 + v_2 - 2v_2) V_1 f_q(\gamma_0) \\
\leq C |v_1 - v_2|.
\]

By Markov’s inequality, the result follows.

Lemma 13 \( \frac{1}{v} \sum_{j=1}^n 2 \Delta_f(x_o) (T_{ij} + T_{2j}) \approx 0. \)

Proof. First,

\[
\mathbb{E} \left[ 2 \Delta_f(x_o) T_{ij} \right] \\
= \int_{-\infty}^{+\infty} \int_{-\infty}^{+\infty} 2 \Delta_f(x_o) g(x_o, q_1) \frac{1}{h} K^x \left( \frac{x-x_o}{h} \right) \left( \frac{2 \gamma_0^-}{h} \right) - \frac{2 \gamma_0^+}{h} f(x_o, q_1) \right) dx_o dq_1 \\
= 2 \Delta_f(x_o) \int_{-\infty}^{+\infty} \int_{-\infty}^{+\infty} g(x_o + u_h, \gamma_0^+ + u_h h) K^x \left( u_h, x_o \right) \left( \frac{2 \gamma_0^-}{h} \right) - \frac{2 \gamma_0^+}{h} f(x_o + u_h, \gamma_0^+ + u_h h) \right) du_h dx_o \\
- 2 \Delta_f(x_o) \int_{-\infty}^{+\infty} \int_{-\infty}^{+\infty} g(x_o + u_h, \gamma_0^- + u_h h) K^x \left( u_h, x_o \right) \left( \frac{2 \gamma_0^-}{h} \right) - \frac{2 \gamma_0^+}{h} f(x_o + u_h, \gamma_0^- + u_h h) \right) du_h dx_o \\
= 2 \Delta_f(x_o) \int_{-\infty}^{+\infty} \int_{-\infty}^{+\infty} \left[ g(x_o, \gamma_0^+ + u_h h) + \cdots + g(x_o, \gamma_0^+ + u_h h) u_h h + o(u_h h) \right] K^x \left( u_h, x_o \right) \left( \frac{2 \gamma_0^-}{h} \right) - \frac{2 \gamma_0^+}{h} du_h dx_o \\
- 2 \Delta_f(x_o) \int_{-\infty}^{+\infty} \int_{-\infty}^{+\infty} \left[ g(x_o, \gamma_0^- + u_h h) + \cdots + g(x_o, \gamma_0^- + u_h h) u_h h + o(u_h h) \right] K^x \left( u_h, x_o \right) \left( \frac{2 \gamma_0^-}{h} \right) - \frac{2 \gamma_0^+}{h} du_h dx_o \\
\approx 2 \Delta_f(x_o) \int_{-\infty}^{+\infty} \left[ g(x_o, \gamma_0^+ + u_h h) + \cdots + g(x_o, \gamma_0^+ + u_h h) h^s \right] - \left( g(x_o, \gamma_0^+ + u_h h) + \cdots + g(x_o, \gamma_0^+ + u_h h) h^s \right) \left( \frac{2 \gamma_0^-}{h} \right) - \frac{2 \gamma_0^+}{h} du_h \right) \\
\approx 2 \Delta_f(x_o) \left( \frac{v}{nh^2} - \frac{v}{\Delta_f^2} \right) \left[ g(x_o, \gamma_0^+ + u_h h) + \cdots + g(x_o, \gamma_0^+ + u_h h) h^s \right] - \left( g(x_o, \gamma_0^+ + u_h h) + \cdots + g(x_o, \gamma_0^+ + u_h h) h^s \right) \left( \frac{2 \gamma_0^-}{h} \right) - \frac{2 \gamma_0^+}{h} du_h \right) \\
\approx 2 \Delta_f(x_o) \left( \frac{v}{nh^2} - \frac{v}{\Delta_f^2} \right)
\]
where \( g(x_j, q_j) = g(x_j, q_j)f(x_j, q_j), \) \( g^{(s)}(x_j, q_j) = \frac{\partial g(x_j, q_j)}{\partial x_j}, \) \( g_2(x_j, q_j) = \frac{\partial g(x_j, q_j)}{\partial q_j}, \) and \( g^{(s)}_1(x_j, q_j) = \frac{\partial g^{(s)}(x_j, q_j)}{\partial q_j}. \)

Similarly, we have

\[
\mathbb{E}[2\Delta f(x_0)T_{2j}] \approx -2\Delta f(x_0) \left( g_2(x_0, \gamma_0) + \cdots + g^{(s)}_{12}(x_0, \gamma_0) h^* \right) \frac{\nu}{n^2 - n^*}.
\]

Hence,

\[
\mathbb{E}\left[ \frac{1}{n} \sum_{j=1}^n 2\Delta f(x_0)\Delta f(x_0) (T_{1j} + T_{2j}) \right] \approx 0.
\]

Next,

\[
\frac{1}{n} \mathbb{E} \left[ 4\Delta f(x_0)T_{1j}^2 \right] = \frac{1}{n} \int_{x_j} \int_{x_j} \frac{4}{n} \Delta f(x_j) g(x_j, q_j)^2 \left( \frac{x_j - x}{h} \right)^2 \left( k^{-} \left( \frac{q_j - \gamma_j}{h} \right) - k^{-} \left( \frac{2q_j - \gamma_j}{h} \right) \right)^2 f(x_j, q_j) dx_j dq_j
\]

\[
= \frac{1}{n} \int_{x_j} \int_{x_j} \left\{ \frac{4}{n} \Delta f^2(x_j) g(x_j + u_x h, \gamma_0 + u_q h)^2 K^* \left( \frac{u_x}{n^*} \right)^2 \left( k^{-} \left( \frac{u_q + \frac{u_x - \gamma_j}{h}}{n^*} \right) - k^{-} (u_q) \right)^2 f(x_j + u_x h, \gamma_0 + u_q h) \right\} du_x du_q
\]

\[
= O \left( \frac{\Delta^2}{nh^* n^* h^*} \right) = o(1),
\]

so

\[
Var \left[ \frac{1}{n} \sum_{j=1}^n 2\Delta f(x_0)T_{1j} \right] = \frac{1}{n} Var \left[ 2\Delta f(x_0)T_{1j} \right] = \frac{1}{n} \mathbb{E} \left[ 4\Delta f^2(x_0)T_{1j}^2 \right] - \frac{1}{n} \mathbb{E} \left[ 2\Delta f(x_0)T_{1j} \right]^2 = o(1).
\]

With similar results derived for \( \frac{1}{n} \sum_{j=1}^n 2\Delta f(x_0)T_{2j} \) and

\[
\frac{1}{n} \mathbb{E} \left[ 4\Delta f(x_0)^2 T_{1j} T_{2j} \right] = \frac{1}{n} \int_{x_j} \int_{x_j} \left\{ \frac{4}{n^2} \Delta f^2(x_j) g(x_j + u_x h, \gamma_0 + u_q h)^2 \left( k^{-} \left( u_q - \frac{u_x - \gamma_j}{n^*} \right) - k^{-} (u_q) \right) \right\} du_x du_q
\]

\[
= O \left( \frac{\Delta^2}{nh^* n^* h^*} \right),
\]

we have

\[
\frac{1}{n} Cov \left[ 2\Delta f(x_0)T_{1j}, 2\Delta f(x_0)T_{2j} \right] = O \left( \frac{1}{n^* h^*} - O \left( \frac{1}{n h^*} \right) \right) = o(1).
\]

As a result,

\[
Var \left[ \frac{1}{n} \sum_{j=1}^n 2\Delta f(x_0)(T_{1j} + T_{2j}) \right] = o(1).
\]

By Markov’s inequality, the result of interest is obtained. \( \blacksquare \)

**Lemma 14** \( \frac{1}{n} \sum_{j=1}^n 2\Delta f(x_0)(T_{5j} + T_{6j}) \approx -2k_-(0) \frac{\nu}{nh^*} f(x_0, \gamma_0)^2. \)

**Proof.** Since

\[
\mathbb{E}[2\Delta f(x_0)T_{5j}] = -\int_{x_j} \int_{x_j} \frac{2}{n} \Delta f(x_j) K^* \left( \frac{x_j - x}{h} \right)^2 \left( k^{-} \left( \frac{q_j - \gamma_j}{h} \right) - k^{-} \left( \frac{2q_j - \gamma_j}{h} \right) \right) \frac{1}{h} 1(\gamma_0 - h \leq q_j \leq \gamma_0) f(x_j, q_j) dx_j dq_j
\]

\[
= -2\Delta f(x_0) \int x_j \int x_j u_x h, \gamma_0 + u_q h) \Delta f^2(u_x, \gamma_0) K^* (u_x, \gamma_0) \left( k^{-} (u_q) - k^{-} (u_q) \right) f(x_j + u_x h, \gamma_0 + u_q h) du_x du_q
\]

\[
= -2\Delta f(x_0) (1, x_0', \gamma_0) \delta_n f(x_0, \gamma_0) (1 + O(h)) \approx -2\Delta f(x_0)
\]
and, similarly,

\[
\mathbb{E}[2\Delta f(x_o)T_{\theta_j}] = \int_{I_j} (x_o, x'_o, q_j) \delta_n K^x \left( \frac{x-x_o}{h}, x_o \right) k^- \left( \frac{q_i-q_o}{h} \right) (\gamma_i - h \leq q_j \leq \gamma_0) f(x_j, q_j) dx_j dq_j
\]

\[
= 2\Delta f(x_o) \int_{1+}^{\infty} \frac{x}{h} \int_{-\infty}^{\infty} dx \left( 1 + x \right) \delta_n K^x (x, x) k^- \left( \frac{q_i-q_o}{h} \right) f(x_o + u_x h, \gamma_0 + u_q h) du_x du_q
\]

\[
= 2\Delta f(x_o)(1, x'_o, \gamma_0) \delta_n f(x_o, \gamma_o) \left( 1 - k_+ \right) \approx 2\Delta f(x_o)(1, -\infty) \frac{v}{n h^2 \Delta^2},
\]

the result of interest follows. 

Lemma 15 \( \text{Cov} (S_1(v_1), S_1(v_2)) \approx \Sigma_o v_2 \). 

Proof. Without loss of generality, assume \( v_1 \geq v_2 \geq 0 \). Then

\[
4nh^2 \text{Cov} \left( \Delta f(x_o) u_j \left( K_{h,j}^{\gamma_1^+} - K_{h,j}^{\gamma_0^+} \right), \Delta f(x_i) u_i \left( K_{h,i}^{\gamma_2^+} - K_{h,i}^{\gamma_0^+} \right) \right)
\]

\[
= 4nh^2 \mathbb{E} \left[ \Delta f(x_o) u_j^2 K^x \left( \frac{x-x_o}{h}, x_o \right) k^+ \left( \frac{q_i-q_o}{h} \right) k^+ \left( \frac{q_j-q_i}{h} \right) \right] \approx 4k_+^2 \approx 4k_+^2 \sigma^2_{\pm}(x_o, \gamma_0)^3 \kappa^2 v_2,
\]

and

\[
4nh^2 \text{Cov} \left( \Delta f(x_o) u_j \left( K_{h,j}^{\gamma_1^-} - K_{h,j}^{\gamma_0^-} \right), \Delta f(x_i) u_i \left( K_{h,i}^{\gamma_2^-} - K_{h,i}^{\gamma_0^-} \right) \right)
\]

\[
\approx 4k_+^2 \sigma^2_{\pm}(x_o, \gamma_0)^3 \kappa^2 v_2,
\]

and

\[
-4nh^2 \text{Cov} \left( \Delta f(x_o) u_j \left( K_{h,j}^{\gamma_1^\pm} - K_{h,j}^{\gamma_0^\pm} \right), \Delta f(x_i) u_i \left( K_{h,i}^{\gamma_2^\pm} - K_{h,i}^{\gamma_0^\pm} \right) \right)
\]

\[
\approx 4k_+^2 (\delta_{\pm} - 0) \sigma^2_{\pm}(x_o, \gamma_0)^3 \kappa^2 v_2.
\]

where \( \sigma^2_{\pm}(x_o) = \mathbb{E}[c_j^2] x_j = x_o, q_j = \gamma_{\pm}, \kappa^2 = \int K(u_x)^2 du_x. \)

If \( v_1 \leq v_2 \leq 0 \), the result is similar except that the term \( \sigma^2_{\pm}(x_o) \) is replaced by \( \sigma^2_{\pm}(x_o) \).

If \( v_1 v_2 < 0 \), then the four terms are all \( o(v_2) \).

Lemma 16 \( \text{Cov} (S_2(v_1), S_2(v_2)) = o(v_2) \), \( \text{Cov} (S_1(v_1), S_2(v_2)) = o(v_2) \) and \( \text{Cov} (S_2(v_1), S_1(v_2)) = o(v_2) \).

Proof. The proof idea is the same as that in Lemma 11.

Lemma 17 For any \( \phi_1, \phi_2 > 0 \), there exists \( \eta > 0 \) such that

\[
P \left\{ \sup_{|v_1 - v_2| < \eta} \left| \frac{1}{n h^d} \left( \tilde{Q}_n \left( \gamma_0^v \right) - \tilde{Q}_n \left( \gamma_0 \right) \right) - \frac{1}{n h^d} \left( \tilde{Q}_n \left( \gamma_0^v \right) - \tilde{Q}_n \left( \gamma_0 \right) \right) \right| > \phi_1 \right\} < \phi_2,
\]

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Proof. Without loss of generality, assume $v_1 \geq v_2 \geq 0$. From Proposition 11, we have

$$
\lim_{n \to \infty} \mathbb{E} \left[ n h^d \left( \tilde{Q}_n (\gamma^*_n) - \tilde{Q}_n (\gamma_0) \right) \right] = 16 k^2 (0) f (x_0, \gamma_0)^3 \kappa^2 \sigma^2_{+} (x_0) (v_1 + v_2 - 2v_2) \leq C |v_1 - v_2|.
$$

By Markov’s inequality, the result follows. ■

Lemma 18 $\| \mathbb{E} [(m_2 - m_1) \ell_{b, 21}] \| = O_p (b^\alpha)$.

Proof. We have

$$
\| \mathbb{E} [(m(x_2, q_2) - m(x_1, q_1)) \ell_{b, 21} | x_1, q_1] \| = \int (m(x_2, q_2) - m(x_1, q_1)) f(x_2, q_2) \frac{1}{b^s} L^2 \left( \frac{x_2 - x_1}{b}, x_1 \right) l \left( \frac{q_2 - q_1}{b} \right) dx_2 dq_2
$$

where $Q_m ((x_2, q_2), (x_1, q_1))$ is the $(s - 1)$-th order Taylor expansion of $m(x_2, q_2)$ at $m(x_1, q_1), R_m ((x_2, q_2), (x_1, q_1))$ is the remainder term, $Q_f ((x_2, q_2), (x_1, q_1))$ is $(\lambda - 1)$-th order Taylor expansion of $f(x_2, q_2)$ at $f(x_1, q_1)$, and $R_f ((x_2, q_2), (x_1, q_1))$ is the remainder term. From Assumption L,

$$
\int Q_m ((x_2, q_2), (x_1, q_1)) (f(x_1, q_1) + Q_f ((x_2, q_2), (x_1, q_1))) \frac{1}{b^s} L^2 \left( \frac{x_2 - x_1}{b}, x_1 \right) l \left( \frac{q_2 - q_1}{b} \right) dx_2 dq_2 = 0,
$$

so $\| \mathbb{E} [(m(x_2, q_2) - m(x_1, q_1)) \ell_{b, 21} | x_1] \|$ is bounded by

$$
\| \int R_m ((x_2, q_2), (x_1, q_1)) f(x_1, q_1) \frac{1}{b^s} L^2 \left( \frac{x_2 - x_1}{b}, x_1 \right) l \left( \frac{q_2 - q_1}{b} \right) dx_2 dq_2 \| + \| \int (m(x_2, q_2) - m(x_1, q_1)) R_f ((x_2, q_2), (x_1, q_1)) \frac{1}{b^s} L^2 \left( \frac{x_2 - x_1}{b}, x_1 \right) l \left( \frac{q_2 - q_1}{b} \right) dx_2 dq_2 \| \leq C b^\eta + C b^{\lambda + 1} \leq C b^\alpha,
$$

where $\eta = \min (\lambda + 1, s)$. ■

Supplement D: Parametric Tests for Threshold Effects when Instruments are present

This supplement discusses the asymptotic distribution of the Wald-type and score-type test statistics under the null and local alternatives when instruments are available. We also provide implementation details for the use of Hansen’s (1996) simulation method in the current context.
For $\gamma \in \Gamma$, define

\[
\Omega_1 (\gamma) = E \left[ z_i x_i' \xi_i^2 1(q_i \leq \gamma) \right], \quad \Omega_2 (\gamma) = E \left[ z_i x_i' \xi_i^2 1(q_i > \gamma) \right],
\]

\[
Q_1 (\gamma) = E \left[ z_i x_i' 1(q_i \leq \gamma) \right], \quad Q_2 (\gamma) = E \left[ z_i x_i' 1(q_i > \gamma) \right],
\]

\[
V_1 (\gamma) = \left[ Q_1 (\gamma)' \Omega_1 (\gamma)^{-1} Q_1 (\gamma) \right]^{-1}, \quad V_2 (\gamma) = \left[ Q_2 (\gamma)' \Omega_2 (\gamma)^{-1} Q_2 (\gamma) \right]^{-1},
\]

\[
\Omega = E \left[ z_i x_i' \xi_i^2 \right], \quad Q = E \left[ z_i x_i' \right], \quad V = \left[ Q' \Omega^{-1} Q \right]^{-1}.
\]

$S_1 (\gamma)$ is a mean-zero Gaussian process with covariance kernel $E \left[ S_1 (\gamma_1) S_1 (\gamma_2)' \right] = \Omega_1 (\gamma_1 \wedge \gamma_2), \quad S = \lim_{\gamma \rightarrow \infty} S_1 (\gamma)$, and $S_2 (\gamma) = S - S_1 (\gamma)$. $S(\gamma)$ is a mean zero Gaussian process with covariance kernel $E \left[ S_1 (\gamma_1) S_2 (\gamma_2)' \right] = \Omega_2 (\gamma_1 \wedge \gamma_2)$.

\[
H(\gamma_1, \gamma_2) = E \left[ (z_i 1(q_i \leq \gamma_1) - Q_1 (\gamma_1) V Q' \Omega^{-1} z_i) (z_i 1(q_i \leq \gamma_2) - Q_1 (\gamma_2) V Q' \Omega^{-1} z_i)' \xi_i^2 \right].
\]

Given the threshold point $\gamma$, the 2SLS estimators for $\beta_1$ and $\beta_2$ are

\[
\bar{\beta}_1 (\gamma) = \left( \hat{Q}_1 (\gamma)' \frac{1}{n} Z_{\leq \gamma}' Z_{\leq \gamma} \right)^{-1} \hat{Q}_1 (\gamma), \quad \bar{\beta}_2 (\gamma) = \left( \hat{Q}_2 (\gamma)' \frac{1}{n} Z_{> \gamma}' Z_{> \gamma} \right)^{-1} \hat{Q}_2 (\gamma),
\]

where $\hat{Q}_1 (\gamma) = n^{-1} \sum_{i=1}^{n} z_i x_i' 1(q_i \leq \gamma)$ and $\hat{Q}_2 (\gamma) = n^{-1} \sum_{i=1}^{n} z_i x_i' 1(q_i > \gamma)$. The residual from this equation is

\[
\bar{e}_i (\gamma) = y_i - x_i' \bar{\beta}_1 (\gamma) 1(q_i \leq \gamma) - x_i' \bar{\beta}_2 (\gamma) 1(q_i > \gamma).
\]

The GMM estimators for $\beta_1$ and $\beta_2$ are

\[
\tilde{\beta}_1 (\gamma) = \left( \hat{Q}_1 (\gamma)' \tilde{\Omega}_1^{-1} (\gamma) \hat{Q}_1 (\gamma) \right)^{-1} \left( \hat{Q}_1 (\gamma)' \tilde{\Omega}_1^{-1} (\gamma) \frac{1}{n} Z_{\leq \gamma}' Y \right),
\]

\[
\tilde{\beta}_2 (\gamma) = \left( \hat{Q}_2 (\gamma)' \tilde{\Omega}_2^{-1} (\gamma) \hat{Q}_2 (\gamma) \right)^{-1} \left( \hat{Q}_2 (\gamma)' \tilde{\Omega}_2^{-1} (\gamma) \frac{1}{n} Z_{> \gamma}' Y \right),
\]

where the weight matrices

\[
\tilde{\Omega}_1 (\gamma) = \frac{1}{n} \sum_{i=1}^{n} z_i x_i' \tilde{\xi}_i^2 (\gamma) 1(q_i \leq \gamma), \quad \tilde{\Omega}_2 (\gamma) = \frac{1}{n} \sum_{i=1}^{n} z_i x_i' \tilde{\xi}_i^2 (\gamma) 1(q_i > \gamma).
\]

The estimated covariance matrices for the GMM estimators are

\[
\tilde{V}_1 (\gamma) = \left( \hat{Q}_1 (\gamma)' \tilde{\Omega}_1^{-1} (\gamma) \hat{Q}_1 (\gamma) \right)^{-1}, \quad \tilde{V}_2 (\gamma) = \left( \hat{Q}_2 (\gamma)' \tilde{\Omega}_2^{-1} (\gamma) \hat{Q}_2 (\gamma) \right)^{-1}.
\]

When $H_0$ holds, $\delta = 0$, and then the 2SLS estimator for $\beta$ is

\[
\tilde{\beta} = \left( \hat{Q}' \left( \frac{1}{n} Z' Z \right)^{-1} \hat{Q} \right)^{-1} \left( \hat{Q}' \left( \frac{1}{n} Z' Z \right)^{-1} \frac{1}{n} Z' Y \right),
\]

where $\hat{Q} = n^{-1} \sum_{i=1}^{n} z_i x_i'$. Note here that the underlying assumption in this specification testing context is $E \left[ \xi | z \right] = 0$, so that the 2SLS estimator can be applied. Correspondingly, the GMM estimator for $\beta$ is

\[
\tilde{\beta} = \left( \hat{Q}' \tilde{\Omega}^{-1} \hat{Q} \right)^{-1} \left( \hat{Q}' \tilde{\Omega}^{-1} \frac{1}{n} Z' Y \right),
\]

and the residual is

\[
\tilde{e}_i = y_i - x_i' \tilde{\beta},
\]
where the weight matrix is
\[
\tilde{\Omega} = \frac{1}{n} \sum_{i=1}^{n} \mathbf{z}_i \mathbf{z}_i' \tilde{e}_i^2
\]
with \( \tilde{e}_i = y_i - \mathbf{x}_i' \hat{\beta} \). The estimated covariance matrix for the GMM estimator is
\[
\hat{V} = \left( \tilde{\Omega} \tilde{Q}^{-1} \tilde{Q} \right)^{-1}.
\]

**Wald-type Tests**

Define
\[
W_n(\gamma) = \left( \hat{V}_1(\gamma) + \hat{V}_2(\gamma) \right)^{-1/2} \sqrt{n} \left( \hat{\beta}_1(\gamma) - \hat{\beta}_2(\gamma) \right), \gamma \in \Gamma.
\]
The Wald-type test statistic is a functional of \( W_n(\gamma) \). Two test statistics are the most popular. The first is the Kolmogorov-Smirnov sup-type statistic
\[
K_n^w = \sup_{\gamma \in \Gamma} \| W_n(\gamma) \|
\]
and the second is the Cramér–von Mises average-type statistic
\[
C_n^w = \int_{\Gamma} \| W_n(\gamma) \| w(\gamma) d\gamma,
\]
where \( w(\gamma) \) in \( C_n^w \) is a known positive weight function with \( \int_{\Gamma} w(\gamma) d\gamma = 1 \). For example, \( w(\tau) = 1/|\Gamma| \) with \(|\Gamma|\) being the length of \( \Gamma \). But if we have some information on the locations where threshold effects are most likely to occur, we can impose larger weights on the neighborhoods of such locations. The choice of the norm \( \| \cdot \| \) is also an issue. The Euclidean norm \( \| \cdot \|_2 \) is obviously natural, e.g., CH use (the square of) this norm. Yu (2013b) suggests using the \( \ell_1 \) norm in testing quantile threshold effects, and Bai (1996) suggests using the \( \ell_\infty \) norm in structural change tests.

The following theorem states the asymptotic distribution of a general continuous functional \( g(\cdot) \) of \( W_n(\cdot) \) under the local alternative \( \delta_n = n^{-1/2} c \). The corresponding test statistic is denoted as \( g_n^w \).

**Theorem 10** If \( \delta_n = n^{-1/2} c \), \( E \left[ \| x \| \right] < \infty, E[y^4] < \infty, E[\varepsilon^4] \) and \( E \left[ \| z \|^4 \right] < \infty \), then
\[

g_n^w \underset{d}{\rightarrow} g_c^c = g(W^c),
\]

where
\[
W^c(\gamma) = (V_1(\gamma) + V_2(\gamma))^{-1/2} \left[ \left( V_1(\gamma) Q_1(\gamma)' \Omega_1(\gamma)^{-1} Q_1(\gamma) \right) - V_2(\gamma) Q_2(\gamma)' \Omega_2(\gamma)^{-1} Q_2(\gamma) \right] c
\]
\[
+ (V_1(\gamma) + V_2(\gamma))^{-1/2} \left[ \left( V_1(\gamma) Q_1(\gamma)' \Omega_1(\gamma)^{-1} S_1(\gamma) - V_2(\gamma) Q_2(\gamma)' \Omega_2(\gamma)^{-1} S_2(\gamma) \right) \right].
\]

**Proof.** Under the local alternative \( \delta_n = n^{-1/2} c \), \( Y = X_{\leq \gamma_0} (\beta + \delta_n) + X_{> \gamma_0} \beta + \varepsilon = X \beta + X_{\leq \gamma_0} \delta_n + \varepsilon \), so that
\[
\hat{\beta}_1(\gamma) = \left( X_{\leq \gamma}' Z_{\leq \gamma} \left( Z_{\leq \gamma}' Z_{\leq \gamma} \right)^{-1} Z_{\leq \gamma}' X_{\leq \gamma} \right)^{-1} \left( X_{\leq \gamma}' Z_{\leq \gamma} \left( Z_{\leq \gamma}' Z_{\leq \gamma} \right)^{-1} Z_{\leq \gamma}' Y \right)
\]
\[
= \beta + O_p(1) \frac{1}{n} \sum_{i=1}^{n} \mathbf{z}_i [x_i' \delta_n 1(q_i \leq \gamma_0 \wedge \gamma) + \varepsilon_i 1(q_i \leq \gamma)]
\]
\[
= \beta + o_p(1) \text{ uniformly in } \gamma \in \Gamma.
\]
where \( \beta \) is uniformly consistent for \( \beta \). As a result,

\[
\tilde{e}_i(\gamma) = y_i - x_i'\tilde{\beta}_1(\gamma) 1(q_i \leq \gamma) - x_i'\tilde{\beta}_2(\gamma) 1(q_i > \gamma)
\]

\[
= x_i'\beta + x_i'\delta_n 1(q_i \leq \gamma_0) + \tilde{e}_i - x_i' (\beta + o_p(1))
\]

\[
= \tilde{e}_i + o_p(||x_i||) \text{ uniformly in } \gamma \in \Gamma,
\]

so that

\[
\tilde{\Omega}_1(\gamma) = \frac{1}{n} \sum_{i=1}^{n} z_i z_i' \tilde{e}_i^2(\gamma) 1(q_i \leq \gamma)
\]

\[
= \frac{1}{n} \sum_{i=1}^{n} z_i z_i' (\tilde{e}_i + o_p(||x_i||))^2 1(q_i \leq \gamma) \xrightarrow{P} \Omega_1(\gamma)
\]

uniformly in \( \gamma \in \Gamma \) by a standard argument. Similarly, \( \tilde{\Omega}_2(\gamma) \xrightarrow{P} \Omega_2(\gamma) \) uniformly in \( \gamma \in \Gamma \). Now,

\[
\sqrt{n} \left( \tilde{\beta}_1(\gamma) - \beta \right) = \left[ \tilde{Q}_1(\gamma)' \tilde{\Omega}_1(\gamma)^{-1} \tilde{Q}_1(\gamma) \right]^{-1} \left[ \tilde{Q}_1(\gamma)' \tilde{\Omega}_1(\gamma)^{-1} \frac{1}{\sqrt{n}} \sum_{i \leq \gamma} (X_i \leq \gamma_0 \delta_n + \tilde{e}) \right],
\]

where \( \tilde{Q}_1(\gamma) \xrightarrow{P} Q_1(\gamma) \), \( \frac{1}{\sqrt{n}} \sum_{i \leq \gamma} X_i \leq \gamma_0 \delta_n \xrightarrow{P} Q_1(\gamma \wedge \gamma_0) c \) uniformly in \( \gamma \in \Gamma \), and \( \frac{1}{\sqrt{n}} \sum_{i \leq \gamma} \tilde{e} \xrightarrow{D} S_1(\gamma) \). Hence

\[
\sqrt{n} \left( \tilde{\beta}_1(\gamma) - \beta \right) \xrightarrow{D} V_1(\gamma) Q_1(\gamma)' \Omega_1(\gamma)^{-1} [Q_1(\gamma \wedge \gamma_0) c + S_1(\gamma)].
\]

Similarly,

\[
\sqrt{n} \left( \tilde{\beta}_2(\gamma) - \beta \right) \xrightarrow{D} V_2(\gamma) Q_2(\gamma)' \Omega_2(\gamma)^{-1} [Q_2(\gamma \vee \gamma_0) c + S_2(\gamma)].
\]

From the arguments above and by the CMT, \( \tilde{V}_1(\gamma) \xrightarrow{P} V_1(\gamma) \) and \( \tilde{V}_2(\gamma) \xrightarrow{P} V_2(\gamma) \) uniformly in \( \gamma \in \Gamma \). Finally, \( W_n(\gamma) \xrightarrow{D} W^c(\gamma) \) as specified in the theorem, where the second part of \( W^c(\gamma) \) is the process in Theorem 4 of CH. □

**Score-type Tests**

The score-type tests are based on

\[
T_n(\gamma) = \left[ \left( z_i 1(q_i \leq \gamma) - \tilde{Q}_1(\gamma) \tilde{V} Q^{\prime} \Omega^{-1} z_i \right) \left( z_i 1(q_i \leq \gamma) - \tilde{Q}_1(\gamma) \tilde{V} Q^{\prime} \Omega^{-1} z_i \right)^{\prime} \right]^{-1/2} \cdot n^{-1/2} \sum_{i=1}^{n} \left[ z_i 1(q_i \leq \gamma) - \tilde{Q}_1(\gamma) \tilde{V} Q^{\prime} \Omega^{-1} z_i \right] \tilde{e}_i, \gamma \in \Gamma.
\]

Note here that although \( \tilde{Q}_1(\gamma) \tilde{V} Q^{\prime} \Omega^{-1} n^{-1/2} \sum_{i=1}^{n} z_i \tilde{e}_i = o_p(1) \), \( z_i 1(q_i \leq \gamma) \) is centered by \( \tilde{Q}_1(\gamma) \tilde{V} Q^{\prime} \Omega^{-1} z_i \). This is because the effect of \( \tilde{\beta} \) will not disappear asymptotically so the asymptotic distribution of \( n^{-1/2} \sum_{i=1}^{n} z_i 1(q_i \leq \gamma) \tilde{e}_i \) differs from \( n^{-1/2} \sum_{i=1}^{n} z_i 1(q_i \leq \gamma) \tilde{e}_i \) under \( H_0 \). Recentering is to offset the effect of \( \tilde{\beta} \). Since only \( \tilde{\beta} \) is used in the construction of \( T_n(\gamma) \), this type of tests is constructed under \( H_0 \) and only one GMM estimator needs to be constructed. This significantly lightens the computation burden. Given \( T_n(\gamma) \), we can similarly construct the Kolmogorov-Smirnov sup-type statistic \( K_n^* \) and the Cramér-von Mises average-type statistic \( C_n^* \).

The following theorem states the asymptotic distribution of a general continuous functional \( g(\cdot) \) of \( T_n(\cdot) \) under the local alternative \( \delta_n = n^{-1/2} c \). The corresponding test statistic is denoted as \( g_n^* \).
Theorem 11 If \( \delta_n = n^{-1/2} \epsilon \), \( \mathbb{E} \left[ \| x \|^4 \right] < \infty \), \( \mathbb{E} \left[ q^4 \right] < \infty \), \( \mathbb{E} \left[ \epsilon^4 \right] \) and \( \mathbb{E} \left[ \| z \|^4 \right] < \infty \), then

\[
g_n^* \overset{d}{\to} g_c^* = g(T^c),
\]

where

\[
T^c(\gamma) = H(\gamma, \gamma)^{-1/2} \{ S(\gamma) + [Q_1(\gamma \wedge \gamma_0) - Q_1(\gamma) V Q' \Omega^{-1} Q_1(\gamma_0)] c \}.
\]

Proof. As in the last theorem, we can show \( \hat{\beta} \overset{p}{\to} \beta, \hat{\Omega} \overset{p}{\to} \Omega \), and \( \hat{V} \overset{p}{\to} V \) under the local alternative.

\[
n^{-1/2} \sum_{i=1}^n z_i 1(q_i \leq \gamma) \varepsilon_i
= n^{-1/2} \sum_{i=1}^n z_i 1(q_i \leq \gamma) \left( y_i - x_i' \beta \right)
= n^{-1/2} \sum_{i=1}^n z_i 1(q_i \leq \gamma) (y_i - x_i' \beta) - n^{-1} \sum_{i=1}^n z_i x_i' 1(q_i \leq \gamma) \sqrt{n} \left( \hat{\beta} - \beta \right)
= n^{-1/2} \sum_{i=1}^n z_i 1(q_i \leq \gamma) (x_i' \delta_n 1(q_i \leq \gamma_0) + \varepsilon_i) - \hat{Q}_1(\gamma) \sqrt{n} \left( \hat{\beta} - \beta \right),
\]

where \( n^{-1/2} \sum_{i=1}^n z_i 1(q_i \leq \gamma) x_i' \delta_n 1(q_i \leq \gamma_0) \overset{p}{\to} Q_1(\gamma \wedge \gamma_0) c, \hat{Q}_1(\gamma) \overset{p}{\to} Q_1(\gamma) \) uniformly in \( \gamma \in \Gamma \), and \( n^{-1/2} \sum_{i=1}^n z_i 1(q_i \leq \gamma) \varepsilon_i \overset{d}{\to} S_1(\gamma) \). Next,

\[
n^{-1/2} \sum_{i=1}^n \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{\Omega}^{-1} z_i \varepsilon_i
= \hat{Q}_1(\gamma) \left( \hat{Q} \hat{\Omega}^{-1} \hat{Q}^{-1} \right) + \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{\Omega}^{-1} \left( n^{-1} \sum_{i=1}^n z_i x_i' 1(q_i \leq \gamma_0) \right) c
+ \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{\Omega}^{-1} \left( n^{-1/2} \sum_{i=1}^n z_i \varepsilon_i \right),
\]

where the second term in the last equality converges in probability to \( Q_1(\gamma) V Q' \Omega^{-1} Q_1(\gamma_0) c \) uniformly in \( \gamma \in \Gamma \), and \( n^{-1/2} \sum_{i=1}^n z_i \varepsilon_i \overset{d}{\to} N(0, \Omega) \). In summary,

\[
n^{-1/2} \sum_{i=1}^n \left[ z_i 1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{\Omega}^{-1} z_i \right] \varepsilon_i
= n^{-1/2} \sum_{i=1}^n \left[ z_i 1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{\Omega}^{-1} z_i \right] \varepsilon_i
+ \left[ Q_1(\gamma \wedge \gamma_0) - Q_1(\gamma) V Q' \Omega^{-1} Q_1(\gamma_0) \right] c + o_p(1)
\sim \left[ Q_1(\gamma \wedge \gamma_0) - Q_1(\gamma) V Q' \Omega^{-1} Q_1(\gamma_0) \right] c,
\]

and it is not hard to show \( n^{-1} \sum_{i=1}^n \left( z_i 1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{\Omega}^{-1} z_i \right) \left( z_i 1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{\Omega}^{-1} z_i \right)^' \overset{p}{\to} H(\gamma, \gamma) \) uniformly in \( \gamma \in \Gamma \), so the results of the theorem follow.

\[
H(\gamma_1, \gamma_2) = \mathbb{E} \left[ \left( z_i 1(q_i \leq \gamma_1) - Q_1(\gamma_1) V Q' \Omega^{-1} z_i \right) \left( z_i 1(q_i \leq \gamma_2) - Q_1(\gamma_2) V Q' \Omega^{-1} z_i \right)^' \right].
\]

To understand \( S(\gamma) \) in \( T^c(\gamma) \), consider a simple case where \( x = (1, x')' \), \( q \) follows a uniform distribution on \([0, 1]\) and is independent of \((z', x', \varepsilon)'\). In this case,

\[
H(\gamma_1, \gamma_2) = (\gamma_1 \wedge \gamma_2) \Omega - \gamma_1 \gamma_2 Q V Q'.
\]
If \( d_z = d \), i.e., the model is just-identified, then

\[
H(\gamma_1, \gamma_2) = \mathbb{E} \left[ (z_i1(q_i \leq \gamma_1) - Q_1(\gamma_1) Q^{-1}_1 z_i)(z_i1(q_i \leq \gamma_2) - Q_1(\gamma_2) Q^{-1}_1 z_i) \varepsilon_i^2 \right] = \Omega_1(\gamma_1 \wedge \gamma_2) - Q_1(\gamma_1) \Omega^{-1}_1(\gamma_2) - \Omega_1(\gamma_1) Q^{-1}_1 Q_1(\gamma_2) + Q_1(\gamma_1) Q^{-1}_1 Q^{-1}_1 Q_1(\gamma_2),
\]

and we can let, for \( \gamma \in \Gamma \),

\[
T_n(\gamma) = \left[ n^{-1} \sum_{i=1}^{n} \left( z_i1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{Q}^{-1}_1 z_i \right) \left( z_i1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{Q}^{-1}_1 z_i \right) \varepsilon_i^2 \right]^{-1/2}
\]

\[
\cdot \left[ n^{-1/2} \sum_{i=1}^{n} \left( z_i1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{Q}^{-1}_1 z_i \right) \varepsilon_i, \gamma \in \Gamma. \tag{36}
\]

Combining these two cases, \( H(\gamma_1, \gamma_2) \) reduces to \( (\gamma_1 \wedge \gamma_2 - \gamma_1 \gamma_2) \Omega \), where \( d_z = d \). In other words, \( \Omega^{-1/2} S(\gamma) \) is a standard \( d \)-dimensional Brownian Bridge. Now, the local power is generated by \( [Q_1(\gamma_0) - Q_1(\gamma) V Q^\Omega^{-1} Q_1(\gamma_0)]c = (\gamma - \gamma_0 - \gamma_0) Q c \). Of course, the construction of \( T_n(\gamma) \) can be greatly simplified in this simple case, e.g., let

\[
T_n(\gamma) = \tilde{\Omega}^{-1/2} \cdot n^{-1/2} \sum_{i=1}^{n} \left( z_i1(q_i \leq \gamma) - \gamma z_i \right) \tilde{\varepsilon}_i,
\]

which converges to the standard \( d \)-dimensional Brownian Bridge. In linear regression, we need only replace \( z_i \) in all formula of (36) by \( x_i \).

### Simulating Critical Values

The asymptotic distributions in the above two theorems are nonpivotal, but the simulation method in Hansen (1996) can be extended to the present case. More specifically, let \( \{\xi^*_i\}_{i=1}^{n} \) be i.i.d. \( N(0, 1) \) random variables, and set

\[
W^*_n(\gamma) = \left( \bar{\tilde{V}}_1(\gamma) + \bar{\tilde{V}}_2(\gamma) \right)^{-1/2} \sqrt{n} \left( \tilde{\beta}^*_1(\gamma) - \tilde{\beta}^*_2(\gamma) \right), \gamma \in \Gamma,
\]

and, for \( \gamma \in \Gamma \),

\[
T^*_n(\gamma) = \left[ n^{-1} \sum_{i=1}^{n} \left( z_i1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{Q}^{-1}_1 z_i \right) \left( z_i1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{Q}^{-1}_1 z_i \right) \varepsilon_i^2 \right]^{-1/2}
\]

\[
\cdot \left[ n^{-1/2} \sum_{i=1}^{n} \left( z_i1(q_i \leq \gamma) - \hat{Q}_1(\gamma) \hat{V} \hat{Q} \hat{Q}^{-1}_1 z_i \right) \tilde{\varepsilon}_i, \xi^*_i, \gamma \in \Gamma. \tag{37}
\]

where \( \tilde{\beta}^*_1(\gamma) \) and \( \tilde{\beta}^*_2(\gamma) \) are similarly defined as \( \tilde{\beta}^*_1(\gamma) \) and \( \tilde{\beta}^*_2(\gamma) \) with the only difference being that \( y_i \) is replaced by \( \tilde{\varepsilon}_i(\gamma) \xi^*_i \); more specifically,

\[
\tilde{\beta}^*_1(\gamma) = \left( \hat{Q}_1(\gamma) \hat{Q}^{-1}_1(\gamma) \hat{Q}_1(\gamma) \right)^{-1} \left( \hat{Q}_1(\gamma) \hat{Q}^{-1}_1(\gamma) \frac{1}{n} \sum_{i=1}^{n} z_i1(q_i \leq \gamma) \tilde{\varepsilon}_i(\gamma) \xi^*_i \right),
\]

\[
\tilde{\beta}^*_2(\gamma) = \left( \hat{Q}_2(\gamma) \hat{Q}^{-1}_2(\gamma) \hat{Q}_2(\gamma) \right)^{-1} \left( \hat{Q}_2(\gamma) \hat{Q}^{-1}_2(\gamma) \frac{1}{n} \sum_{i=1}^{n} z_i1(q_i > \gamma) \tilde{\varepsilon}_i(\gamma) \xi^*_i \right).
\]

Our test rejects \( H_0 \) if \( g^*_n(g^*_n) \) is greater than the \( (1 - \alpha) \)th conditional quantile of \( g(W^*_n(\gamma)) (g(T^*_n(\gamma))) \). Equivalently, the \( p \)-value transformation can be employed. Take the score test as an example. Define \( p^*_n = 1 - F^*_n(g^*_n) \), and \( p_n = 1 - F_0(g^*_n) \), where \( F^*_n \) is the conditional distribution of \( g(T^*_n(\gamma)) \) given the

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original data, and $F_0$ is the asymptotic distribution of $g(T_n(\gamma))$ under the null. Our test rejects $H_0$ if $p_n^* \leq \alpha$. By stochastic equicontinuity of the $T_n(\gamma)$ process, we can replace $\Gamma$ by finite grids with the distance between adjacent grid points going to zero as $n \to \infty$. A natural choice of the grids for $\Gamma$ is the $q_i$’s in $\Gamma$. Also, the conditional distribution can be approximated by standard simulation techniques. More specifically, the following procedure is used.

**Step 1:** generate $\{\xi^*_i\}_{i=1}^n$ be i.i.d. $N(0, 1)$ random variables.

**Step 2:** set $T_n^*(\gamma_l)$ as in (37), where $\{\gamma_l\}_{l=1}^L$ is a grid approximation of $\Gamma$. Note here that the same $\{\xi^*_i\}_{i=1}^n$ are used for all $\gamma_l$, $l = 1, \ldots, L$.

**Step 3:** set $g_n^*(\gamma_l) = g(T_n^*(\gamma_l))$.

**Step 4:** repeat Step 1-3 $J$ times to generate $\{g_n^*(\gamma_l)\}_{j=1}^J$.

**Step 5:** if $p_n^* = J^{-1} \sum_{j=1}^J 1 \{g_n^* \geq g^*_n\} \leq \alpha$, we reject $H_0$; otherwise, accept $H_0$.

It can be shown that $p_n^* = p_n + o_p(1)$ under both the null and local alternative. Hence $p_n^* \to p_c = 1 - F_0(g^*_c)$ under the local alternative, and $p_n^* \to U$, the uniform distribution on $[0, 1]$, under the null. The proof is similar to that of Yu (2013b, 2016) and so it is omitted here.

### Additional References


