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December 2005

Online at http://mpra.ub.uni-muenchen.de/11982/ MPRA Paper No. 11982, posted 07. December 2008 / 07:40

# Bandwidth Selection in Nonparametric Kernel Testing ${ }^{1}$ 

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#### Abstract

We propose a sound approach to bandwidth selection in nonparametric kernel testing. The main idea is to find an Edgeworth expansion of the asymptotic distribution of the test concerned. Due to the involvement of a kernel bandwidth in the leading term of the Edgeworth expansion, we are able to establish closed-form expressions to explicitly represent the leading terms of both the size and power functions and then determine how the bandwidth should be chosen according to certain requirements for both the size and power functions. For example, when a significance level is given, we can choose the bandwidth such that the power function is maximized while the size function is controlled by the significance level. Both asymptotic theory and methodology are established. In addition, we develop an easy implementation procedure for the practical realization of the established methodology and illustrate this on two simulated examples and a real data example.


Keywords: Choice of bandwidth parameter, Edgeworth expansion, nonparametric kernel testing, power function, size function.

[^0]
## 1. Introduction

Consider a nonparametric regression model of the form

$$
\begin{equation*}
Y_{i}=m\left(X_{i}\right)+e_{i}, i=1,2, \ldots, n \tag{1.1}
\end{equation*}
$$

where $\left\{X_{i}\right\}$ is a sequence of strictly stationary time series variables, $\left\{e_{i}\right\}$ is a sequence of independent and identically distributed (i.i.d.) errors with $E\left[e_{1}\right]=0$ and $0<$ $E\left[e_{1}^{2}\right]=\sigma^{2}<\infty, m(\cdot)$ is an unknown function defined over $\mathbb{R}^{d}$ for $d \geq 1$, and $n$ is the number of observations. We assume that $\left\{X_{i}\right\}$ and $\left\{e_{j}\right\}$ are independent for all $1 \leq i \leq j \leq n$.

To avoid the so-called "curse of dimensionality" problem, we mainly consider the case of $1 \leq d \leq 3$ in this paper. For the case of $d \geq 4$, various dimension reduction estimation and specification methods have been discussed extensively in several monographs, such as Fan and Gijbels (1996), Hart (1997), Fan and Yao (2003), and Li and Racine (2007).

There is a vast literature on testing a parametric regression model (null hypothesis) versus a nonparametric model, especially for the case of i.i.d. $X_{i}$ 's (random or fixed design case). Many goodness-of-fit testing procedures are based on evaluating a distance between a parametric estimate of the regression function $m$ (assuming the null hypothesis is true) and a nonparametric estimate of that function. Among the popular choices for a nonparametric kernel estimator for $m$ are the Nadaraya-Watson estimator, the Gasser-Müller estimator and a local linear (polynomial) estimator. Earlier papers following this approach of evaluating such a distance include Härdle and Mammen (1993), Weihrather (1993) and González-Manteiga and Cao (1993), among others. Härdle and Mammen (1993) consider a weighted $L_{2}$-distance between a parametric estimator and a nonparametric Nadaraya-Watson estimator of the regression function. The asymptotic distribution of their test statistic under the null hypothesis depends on the unknown error variance (the conditional error variance function). Weihrather (1993) instead uses a Gasser-Müller nonparametric estimator
in the fixed design regression case, divides by an estimator of the error variance and considers a discretized version of the $L_{2}$-distance. González-Manteiga and Cao (1993) also consider the fixed design regression case but rely on minimum distance estimation of the parametric model, seeking for minimizing a weighted $L_{2}$-type distance between the parametric model and a pilot nonparametric estimator.

Another approach to the same testing problem is introduced in Dette (1999) who focusses on the integrated conditional variance function, and uses as a test statistic the difference of a parametric estimator and a nonparametric (Nadaraya-Watson based) estimator of this integrated variance. It is shown that this estimator (asymptotically) corresponds to test statistics based on a weighted $L_{2}$-distance between a parametric and nonparametric estimator of the regression function, as in the above mentioned papers, using an appropriate weight function in defining the $L_{2}$-distance. Dette (1999) studies the asymptotic distribution of the test statistic under fixed alternatives. Such kind of alternatives are to be distinguished from the so-called sequences of local alternatives, where the difference between the regression function under the alternative and the one under the null hypothesis depends on the sample size $n$ and decreases with $n$. The latter setup is the one considered in our study.

The above papers and several more recent goodness-of-fit tests (see for example Zhang and Dette (2004) and references therein) have in common that they rely on nonparametric kernel type regression estimators and that the resulting test statistics are of a similar form (at least in first-order asymptotics), and all depend on a bandwidth parameter. The choice of the bandwidth parameter in such goodness-of-fit testing procedures is the main concern in the present paper. Roughly speaking one can distinguish in the literature two approaches to deal with this bandwidth parameter choice in nonparametric and semiparametric kernel methods used for constructing model specification tests for the mean function of model (1.1). A first approach is to use an estimation-based optimal bandwidth value, such as a cross-validation bandwidth. A second approach is to consider a set of suitable values for the bandwidth
and proceed further from there.
Existing studies based on the first approach include Härdle and Mammen (1993) for testing nonparametric regression with i.i.d. designs and errors, Hjellvik and Tjøstheim (1995), and Hjellvik, Yao and Tjøstheim (1998) for testing linearity in dependent time series cases, Li (1999) for specification testing in econometric time series cases, Chen, Härdle and Li (2003) for using empirical likelihood-based tests, Juhl and Xiao (2005) for testing structural change in nonparametric time series regression, and others. As pointed out in the literature, such choices cannot be justified in both theory and practice since estimation-based optimal values may not be optimal for testing purposes.

Nonparametric tests involving the second approach of choising either a set of suitable bandwidth values for the kernel case or a sequence of positive integers for the smoothing spline case include Fan (1996), Fan, Zhang and Zhang (2001), and Horowitz and Spokoiny (2001). The practical implementation of choosing such sets or sequences is however problematic. This is probably why Horowitz and Spokoiny (2001) develop their theoretical results based on a set of suitable bandwidths on the one hand, but choose their practical bandwidth values based on the assessment of the power function of their test on the other hand. Apart from using such test statistics based on nonparametric kernel, nonparametric series, spline smoothing and wavelet methods, there are test statistics constructed and studied based on empirical distributions. Such studies have recently been summarized in Zhu (2005).

To the best of our knowledge, the idea of choosing the appropriate smoothing parameter such that the size of the test under consideration is preserved while maximizing the power against a given alternative was only first explored analytically by Kulasekera and Wang (1997), in which the authors propose using a nonparametric kernel test to check whether the mean functions of two data sets can be identical in a nonparametric fixed design setting. In some other closely related studies, various discussions have been given on the comparison of power values of the same test at
different bandwidths or different tests at the same bandwidth. Such studies include Hart (1997), Hjellvik, Yao and Tjøstheim (1998), Hunsberger and Follmann (2001), and Zhang and Dette (2004). The last paper compares three main types of nonparametric kernel tests proposed in Härdle and Mammen (1993), Zheng (1996), and Fan, Zhang and Zhang (2001).

On the issue of size correction, there have recently been some studies. For example, Fan and Linton (2003) develop an Edgeworth expansion for the size function of their test and then propose using corrected asymptotic critical values to improve the smallmedium sample size properties of the size of their test. Some other related studies include Nishiyama and Robinson (2000), Horowitz (2003), Nishiyama and Robinson (2005), who develop some useful Edgeworth expansions for bootstrap distributions of partial-sum type of tests for improving the size performance.

The current paper is motivated by such existing studies, especially by Kulasekera and Wang (1997), Fan and Linton (2003), Dette and Spreckelsen (2004), and Zhang and Dette (2004), to develop a solid theory to support a power function-based bandwidth selection procedure such that the power of the proposed test is maximized while the size is under control when using nonparametric kernel testing in parametric specification of a nonparametric regression model of the form (1.1) associated with the hypothesis form of (1.2) below.

To state the main results of this paper, we introduce some notational details. The main interest of this paper is to test a parametric null hypothesis of the form

$$
\begin{align*}
& \mathcal{H}_{0}: m(x)=m_{\theta_{0}}(x) \quad \text { versus a sequence of alternatives of the form } \\
& \mathcal{H}_{1}: m(x)=m_{\theta_{1}}(x)+\Delta_{n}(x) \text { for all } x \in \mathbb{R}^{d} \tag{1.2}
\end{align*}
$$

where both $\theta_{0}, \theta_{1} \in \Theta$ are unknown parameters and $\Theta$ is a parameter space of $\mathbb{R}^{p}$, and $\Delta_{n}(x)$ is a sequence of nonparametrically unknown functions over $\mathbb{R}^{d}$. With $\Delta_{n}(x)$ not being equal to zero, the function $m_{\theta_{1}}(x)$ in $\mathcal{H}_{1}$ is in fact the projection of the true function on the null model.

Note that $m(x)$ under $\mathcal{H}_{1}$ in (1.2) is semiparametric when $\left\{\Delta_{n}(x)\right\}$ is unknown nonparametrically. Note also that instead of requiring (1.2) for all $x \in \mathbb{R}^{d}$, it may be assumed that (1.2) holds with probability one for $x=X_{i}$. Some first-order asymptotic properties for both the size and power functions of a nonparametric kernel test for the case where $\Delta_{n}(\cdot) \equiv \Delta(\cdot)$, corresponding to a class of fixed alternatives (not depending on $n$ ), have already been discussed in the literature, such as Dette and Spreckelsen (2004). This paper focuses on studying higher-order asymptotic properties of such kernel tests for the case where $\left\{\Delta_{n}(\cdot)\right\}$ is a sequence of local alternatives in the sense that $\lim _{n \rightarrow \infty} \Delta_{n}(x)=0$ for all $x \in \mathbb{R}^{d}$.

Let $K(\cdot)$ be the probability kernel density function and $h$ be the bandwidth involved in the construction of a nonparametric kernel test statistic denoted by $\widehat{T}_{n}(h)$. To implement the kernel test in practice, we propose a new bootstrap simulation procedure to approximate the $1-\alpha$ quantile of the distribution of the kernel test by a bootstrap simulated critical value $l_{\alpha}$. Let $\alpha_{n}(h)=P\left(\widehat{T}_{n}(h)>l_{\alpha} \mid \mathcal{H}_{0}\right)$ and $\beta_{n}(h)=P\left(\widehat{T}_{n}(h)>l_{\alpha} \mid \mathcal{H}_{1}\right)$ be the respective size and power functions. In Theorem 2.2 we show that

$$
\begin{align*}
& \alpha_{n}(h)=1-\Phi\left(l_{\alpha}-s_{n}\right)-\kappa_{n}\left(1-\left(l_{\alpha}-s_{n}\right)^{2}\right) \phi\left(l_{\alpha}-s_{n}\right)+o\left(\sqrt{h^{d}}\right),  \tag{1.3}\\
& \beta_{n}(h)=1-\Phi\left(l_{\alpha}-r_{n}\right)-\kappa_{n}\left(1-\left(l_{\alpha}-r_{n}\right)^{2}\right) \phi\left(l_{\alpha}-r_{n}\right)+o\left(\sqrt{h^{d}}\right), \tag{1.4}
\end{align*}
$$

where $s_{n}=p_{1} \sqrt{h^{d}}, r_{n}=p_{2} n \delta_{n}^{2} \sqrt{h^{d}}, \kappa_{n}=p_{3} \sqrt{h^{d}}$, and $\Phi(\cdot)$ and $\phi(\cdot)$ denote respectively the cumulative distribution and density function of the standard Normal random variable, in which all $p_{i}$ 's are positive constants and $\delta_{n}^{2}=\int \Delta_{n}^{2}(x) \pi^{2}(x) d x$ with $\pi(\cdot)$ being the marginal density function of $\left\{X_{i}\right\}$.

Our aim is to choose a bandwidth $h_{\text {ew }}$ such that $\beta_{n}\left(h_{\text {ew }}\right)=\max _{h \in H_{n}(\alpha)} \beta_{n}(h)$ with $H_{n}(\alpha)=\left\{h: \alpha-c_{\min }<\alpha_{n}(h)<\alpha+c_{\min }\right\}$ for some small $0<c_{\min }<\alpha$. Our detailed study in Section 3 shows that $h_{\text {ew }}$ is proportional to $\left(n \delta_{n}^{2}\right)^{-\frac{3}{2 d}}$. Such established relationship between $\delta_{n}$ and $h_{\text {ew }}$ shows us that the choice of an optimal rate of $h_{\text {ew }}$ depends on that of an order of $\delta_{n}$.

If $\delta_{n}$ is chosen proportional to $n^{-\frac{d+12}{6(d+4)}}$ for a sequence of local alternatives under $\mathcal{H}_{1}$, then the optimal rate of $h_{\text {ew }}$ is proportional to $n^{-\frac{1}{d+4}}$, which is the order of a nonparametric cross-validation estimation-based bandwidth frequently used for testing purposes. When considering a sequence of local alternatives with $\delta_{n}=O\left(n^{-\frac{1}{2}} \sqrt{\log \log n}\right)$ being chosen as the optimal rate for testing in this kind of kernel testing (Horowitz and Spokoiny 2001), the optimal rate of $h_{\text {ew }}$ is proportional to $(\log \log n)^{-\frac{3}{2 d}}$.

The rest of the paper is organised as follows. Section 2 points out that existing nonparametric kernel tests can be decomposed with quadratic forms of $\left\{e_{i}\right\}$ as leading terms in the decomposition. This motivates the discussion about establishing Edgeworth expansions for such quadratic forms. In Section 3, we apply the Edgeworth expansions to study both the size and power functions of a representative kernel test. Section 4 presents several examples of implementation. Some concluding remarks are made in Section 5. Mathematical assumptions and proofs are provided in the appendix.

## 2. Nonparametric kernel testing

As mentioned in the introductory section, various authors have discussed and studied nonparametric kernel test statistics based on a (weighted) $L_{2}$-distance function between a nonparametric kernel estimator and a parametric counterpart of the mean function. It can be shown that the leading term of each of these nonparametric kernel test statistics is of a quadratic form (see, for example, Chen, Härdle and Li 2003)

$$
\begin{equation*}
P_{n}(h)=\sum_{i=1}^{n} \sum_{j=1}^{n} e_{i} w\left(X_{i}\right) L_{h}\left(X_{i}-X_{j}\right) w\left(X_{j}\right) e_{j}, \tag{2.1}
\end{equation*}
$$

where $L_{h}(\cdot)=\frac{1}{n \sqrt{h^{d}}} L(\dot{\bar{h}}), L(x)=\int K(y) K(x+y) d y$, and $w(\cdot)$ is a suitable weight function probably depending on either $\pi(\cdot), \sigma^{2}(\cdot)$ or both, in which $K(\cdot)$ is a probability kernel function, $h$ is a bandwidth parameter and both are involved in a nonparametric kernel estimation of $m(\cdot)$.

In this paper, we concentrate on a second group of nonparametric kernel test statistics using a different distance function. Rewrite model (1.1) into a notational version of the form under $\mathcal{H}_{0}$

$$
\begin{equation*}
Y=m_{\theta_{0}}(X)+e, \tag{2.2}
\end{equation*}
$$

where $X$ is assumed to be random and $\theta_{0}$ is the true value of $\theta$ under $\mathcal{H}_{0}$. Obviously, $E[e \mid X]=0$ under $\mathcal{H}_{0}$. Existing studies (Zheng 1996; Li and Wang 1998; Li 1999; Fan and Linton 2003; Dette and Spreckelsen 2004; Juhl and Xiao 2005) propose using a distance function of the form

$$
\begin{equation*}
E[e E(e \mid X) \pi(X)]=E\left[\left(E^{2}(e \mid X)\right) \pi(X)\right] \tag{2.3}
\end{equation*}
$$

where $\pi(\cdot)$ is the marginal density function of $X$.
This suggests using a normalized kernel-based sample analogue of (2.3) of the form

$$
\begin{equation*}
T_{n}(h)=\frac{1}{n \sqrt{h^{d}} \sigma_{n}} \sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} e_{i} K\left(\frac{X_{i}-X_{j}}{h}\right) e_{j}, \tag{2.4}
\end{equation*}
$$

where $\sigma_{n}^{2}=2 \mu_{2}^{2} \nu_{2} \int K^{2}(u) d u$ with $\mu_{k}=E\left[e_{1}^{k}\right]$ for $k \geq 1$ and $\nu_{l}=E\left[\pi^{l}\left(X_{1}\right)\right]$ for $l \geq 1$.
It can easily be seen that $T_{n}(h)$ is the leading term of the following quadratic form

$$
\begin{equation*}
Q_{n}(h)=\frac{1}{n \sqrt{h^{d}} \sigma_{n}} \sum_{i=1}^{n} \sum_{j=1}^{n} e_{i} K\left(\frac{X_{i}-X_{j}}{h}\right) e_{j} . \tag{2.5}
\end{equation*}
$$

In summary, both equations (2.1) and (2.5) can be generally written as

$$
\begin{equation*}
R_{n}(h)=\sum_{i=1}^{n} \sum_{j=1}^{n} e_{i} \phi_{n}\left(X_{i}, X_{j}\right) e_{j}, \tag{2.6}
\end{equation*}
$$

where $\phi_{n}(\cdot, \cdot)$ may depend on $n$, the bandwidth $h$ and the kernel function $K$.
Thus, it is of general interest to study asymptotic distributions and their Edgeworth expansions for quadratic forms of type (2.6). To present the main ideas of establishing Edgeworth expansions for such quadratic forms, we focus on $T_{n}(h)$ in the rest of this paper. This is because the main technology for establishing an Edgeworth
expansion for the asymptotic distribution of each of such tests is the same as that for $T_{n}(h)$.

Since $T_{n}(h)$ involves some unknown quantities, we estimate it by a stochastically normalized version of the form

$$
\begin{equation*}
\widehat{T}_{n}(h)=\frac{\sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} \widehat{e}_{i} K\left(\frac{X_{i}-X_{j}}{h}\right) \widehat{e}_{j}}{n \sqrt{h^{d}} \widehat{\sigma}_{n}}, \tag{2.7}
\end{equation*}
$$

where $\widehat{e}_{i}=Y_{i}-m_{\widehat{\theta}}\left(X_{i}\right)$ and $\widehat{\sigma}_{n}^{2}=2 \widehat{\mu}_{2}^{2} \widehat{\nu}_{2} \int K^{2}(u) d u$ with $\widehat{\mu}_{2}=\frac{1}{n} \sum_{i=1}^{n} \widehat{e}_{i}^{2}$ and $\widehat{\nu}_{2}=\frac{1}{n} \sum_{i=1}^{n} \hat{\pi}^{2}\left(X_{i}\right)$, in which $\hat{\theta}$ is a $\sqrt{n}$-consistent estimator of $\theta_{0}$ under $\mathcal{H}_{0}$ and $\widehat{\pi}(x)=\frac{1}{n \hat{b}_{\widehat{c v}}^{d}} \sum_{i=1}^{n} K\left(\frac{x-X_{i}}{\widehat{b}_{\mathrm{cv}}}\right)$ is the conventional nonparametric kernel density estimator with $\widehat{b}_{\mathrm{cv}}$ being a bandwidth parameter chosen by cross-validation (see for example Silverman 1986).

Similarly to existing results (Li 1999), it may be shown that for each given $h$

$$
\begin{equation*}
\widehat{T}_{n}(h)=T_{n}(h)+o_{P}\left(\sqrt{h^{d}}\right) . \tag{2.8}
\end{equation*}
$$

Thus, we may use the distribution of $\widehat{T}_{n}(h)$ to approximate that of $T_{n}(h)$. Let $l_{\alpha}^{e}$ $(0<\alpha<1)$ be the $1-\alpha$ quantile of the exact finite-sample distribution of $\widehat{T}_{n}(h)$. Because $l_{\alpha}^{e}$ may not be evaluated in practice, we therefore suggest choosing either a non-random approximate $\alpha$-level critical value, $l_{\alpha}$, or a stochastic approximate $\alpha$-level critical value, $l_{\alpha}^{*}$ by using the following simulation procedure:

- We generate $Y_{i}^{*}=m_{\widehat{\theta}}\left(X_{i}\right)+\sqrt{\widehat{\mu}_{2}} e_{i}^{*}$ for $1 \leq i \leq n$, where $\left\{e_{i}^{*}\right\}$ is a sequence of i.i.d. random samples drawn from a pre-specified distribution, such as $N(0,1)$. Use the data set $\left\{\left(X_{i}, Y_{i}^{*}\right): i=1,2, \ldots, n\right\}$ to estimate $\hat{\theta}$ by $\widehat{\theta}^{*}$ and compute $\widehat{T}_{n}(h)$. Let $l_{\alpha}$ be the $1-\alpha$ quantile of the distribution of

$$
\begin{equation*}
\widehat{T}_{n}^{*}(h)=\frac{\sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} \widehat{e}_{i}^{*} K\left(\frac{X_{i}-X_{j}}{h}\right) \widehat{e}_{j}^{*}}{n \sqrt{h^{d}} \widehat{\sigma}_{n}^{*}}, \tag{2.9}
\end{equation*}
$$

where $\widehat{e}_{i}^{*}=Y_{i}^{*}-m_{\widehat{\theta}^{*}}\left(X_{i}\right)$ and $\widehat{\sigma}_{n}^{* 2}=2 \widehat{\mu}_{2}^{* 2} \widehat{\nu}_{2} \int K^{2}(u) d u$ with $\widehat{\mu}_{2}^{*}=\frac{1}{n} \sum_{i=1}^{n} \widehat{e}_{i}^{* 2}$.
In the simulation process, the original sample $\mathcal{X}_{n}=\left(X_{1}, \cdots, X_{n}\right)$ acts in the resampling as a fixed design even when $\left\{X_{i}\right\}$ is a sequence of random regressors.

- Repeat the above step $M$ times and produce $M$ versions of $\widehat{T}_{n}^{*}(h)$ denoted by $\widehat{T}_{n, m}^{*}(h)$ for $m=1,2, \ldots, M$. Use the $M$ values of $\widehat{T}_{n, m}^{*}(h)$ to construct their empirical distribution function. The bootstrap distribution of $\widehat{T}_{n}^{*}(h)$ given $\mathcal{W}_{n}=\left\{\left(X_{i}, Y_{i}\right): 1 \leq i \leq n\right\}$ is defined by $P^{*}\left(\widehat{T}_{n}^{*}(h) \leq x\right)=P\left(\widehat{T}_{n}^{*}(h) \leq x \mid \mathcal{W}_{n}\right)$. Let $l_{\alpha}^{*}(0<\alpha<1)$ satisfy $P^{*}\left(\widehat{T}_{n}^{*}(h) \geq l_{\alpha}^{*}\right)=\alpha$ and then estimate $l_{\alpha}$ by $l_{\alpha}^{*}$.

Note that both $l_{\alpha}=l_{\alpha}(h)$ and $l_{\alpha}^{*}=l_{\alpha}^{*}(h)$ depend on $h$. It should be pointed out that the choice of a pre-specified distribution does not have much impact on both the theoretical and practical results. In addition, we may also use a wild bootstrap procedure to generate a sequence of resamples for $\left\{e_{i}^{*}\right\}$.

Note also that the above simulation is based on the so-called regression bootstrap simulation procedure discussed in the literature, such as Li and Wang (1998), Franke, Kreiss and Mammen (2002), and Li and Racine (2007). When $X_{i}=Y_{i-1}$, we may also use a recursive simulation procedure, commonly-used in the literature. See for example, Hjellvik and Tjøstheim (1995), and Franke, Kreiss and Mammen (2002).

Since the choice of a simulation procedure does not affect the establishment of our theory, our main results are established based on the proposed simulation procedure. We now have the following results in Theorems 2.1 and 2.2; their proofs are provided in the appendix.

Theorem 2.1. Suppose that Assumptions A. 1 and A.2 listed in the appendix hold. Then under $\mathcal{H}_{0}$

$$
\begin{equation*}
\sup _{x \in R^{1}}\left|P^{*}\left(\widehat{T}_{n}^{*}(h) \leq x\right)-P\left(\widehat{T}_{n}(h) \leq x\right)\right|=O\left(\sqrt{h^{d}}\right) \tag{2.10}
\end{equation*}
$$

holds in probability with respect to the joint distribution of $\mathcal{W}_{n}$, and

$$
\begin{equation*}
P\left(\widehat{T}_{n}(h)>l_{\alpha}^{*}\right)=\alpha+O\left(\sqrt{h^{d}}\right) . \tag{2.11}
\end{equation*}
$$

For an equivalent test, Li and Wang (1998) establish some results weaker than (2.10). Fan and Linton (2003) consider some higher-order approximations to the size function of the test discussed in Li and Wang (1998).

For each $h$ we define the following size and power functions

$$
\begin{equation*}
\alpha_{n}(h)=P\left(\widehat{T}_{n}(h)>l_{\alpha} \mid \mathcal{H}_{0}\right) \quad \text { and } \quad \beta_{n}(h)=P\left(\widehat{T}_{n}(h)>l_{\alpha} \mid \mathcal{H}_{1}\right) . \tag{2.12}
\end{equation*}
$$

Correspondingly, we define $\left(\alpha_{n}^{*}(h), \beta_{n}^{*}(h)\right)$ with $l_{\alpha}$ replaced by $l_{\alpha}^{*}$.
Before we discuss how to choose an optimal bandwidth in Section 3, we give Edgeworth expansions of both the size and power functions in Theorem 2.2 below. In order to express the Edgeworth expansions, we need to introduce the following notation. Let

$$
\begin{equation*}
\kappa_{n}=\frac{\sqrt{h^{d}}\left(\frac{\mu_{3}^{2} K^{2}(0)}{n h^{d}}+\frac{4 \mu_{2}^{3} \nu_{3}}{3} K^{(3)}(0)\right)}{\sigma_{n}^{3}} \tag{2.13}
\end{equation*}
$$

where $\nu_{l}=E\left[\pi^{l}\left(X_{1}\right)\right]=\int \pi^{l+1}(x) d x$, and $K^{(3)}(\cdot)$ is the three-time convolution of $K(\cdot)$ with itself.

Theorem 2.2. (i) Suppose that Assumptions A. 1 and A. 2 listed in the appendix hold. Then

$$
\begin{align*}
& \alpha_{n}(h)=1-\Phi\left(l_{\alpha}-s_{n}\right)-\kappa_{n}\left(1-\left(l_{\alpha}-s_{n}\right)^{2}\right) \phi\left(l_{\alpha}-s_{n}\right)+o\left(\sqrt{h^{d}}\right),  \tag{2.14}\\
& \alpha_{n}^{*}(h)=1-\Phi\left(l_{\alpha}^{*}-s_{n}\right)-\kappa_{n}\left(1-\left(l_{\alpha}^{*}-s_{n}\right)^{2}\right) \phi\left(l_{\alpha}^{*}-s_{n}\right)+o\left(\sqrt{h^{d}}\right) \tag{2.15}
\end{align*}
$$

hold in probability with respect to the joint distribution of $\mathcal{W}_{n}$, where $\Phi(\cdot)$ and $\phi(\cdot)$ are the probability distribution and density functions of $N(0,1)$, respectively, and $s_{n}=C_{0}(m) \sqrt{h^{d}}$ with

$$
C_{0}(m)=\frac{\int\left(\frac{\partial m_{\theta_{0}}(x)}{\partial \theta}\right)^{\tau}\left(E\left[\left(\frac{m_{\theta_{0}}\left(X_{1}\right)}{\partial \theta}\right)\left(\frac{m_{\theta_{0}}\left(X_{1}\right)}{\partial \theta}\right)^{\tau}\right]\right)^{-1}\left(\frac{m_{\theta_{0}}(x)}{\partial \theta}\right) \pi^{2}(x) d x}{\sqrt{2 \nu_{2} \int K^{2}(v) d v}} .
$$

(ii) Suppose that Assumptions A.1-A.3 listed in the appendix hold. Then the following equations hold in probability with respect to the joint distribution of $\mathcal{W}_{n}$ :

$$
\begin{align*}
& \beta_{n}(h)=1-\Phi\left(l_{\alpha}-r_{n}\right)-\kappa_{n}\left(1-\left(l_{\alpha}-r_{n}\right)^{2}\right) \phi\left(l_{\alpha}-r_{n}\right)+o\left(\sqrt{h^{d}}\right),  \tag{2.16}\\
& \beta_{n}^{*}(h)=1-\Phi\left(l_{\alpha}^{*}-r_{n}\right)-\kappa_{n}\left(1-\left(l_{\alpha}^{*}-r_{n}\right)^{2}\right) \phi\left(l_{\alpha}^{*}-r_{n}\right)+o\left(\sqrt{h^{d}}\right), \tag{2.17}
\end{align*}
$$

where $r_{n}=n C_{n}^{2} \sqrt{h^{d}}$, in which

$$
\begin{equation*}
C_{n}^{2}=\frac{\int \Delta_{n}^{2}(x) \pi^{2}(x) d x}{\sigma^{2} \sqrt{2 \nu_{2} \int K^{2}(v) d v}} . \tag{2.18}
\end{equation*}
$$

Assumption A. 2 implies that the random quantity $C_{0}(m)$ is bounded in probability. As expected, the rate of $r_{n}$ depends on the form of $\Delta_{n}(\cdot)$.

To simplify the following expressions, let $z_{\alpha}$ be the $1-\alpha$ quantile of the standard normal distribution and $d_{j}=\left(z_{\alpha}^{2}-1\right) c_{j}$ for $j=1,2$, where

$$
\begin{equation*}
c_{1}=\frac{4 K^{(3)}(0) \mu_{2}^{3} \nu_{3}}{3 \sigma_{n}^{3}} \quad \text { and } \quad c_{2}=\frac{\mu_{3}^{2} K^{2}(0)}{\sigma_{n}^{3}} . \tag{2.19}
\end{equation*}
$$

Let $d_{0}=d_{1}-C_{0}(m)$. A corollary of Theorem 2.2 is given in Theorem 2.3 below.
Theorem 2.3. Suppose that the conditions of Theorem 2.2(i) hold. Then under $\mathcal{H}_{0}$

$$
\begin{align*}
& l_{\alpha} \approx z_{\alpha}+d_{0} \sqrt{h^{d}}+d_{2} \frac{1}{n \sqrt{h^{d}}} \text { in probability, }  \tag{2.20}\\
& l_{\alpha}^{*} \approx z_{\alpha}+d_{0} \sqrt{h^{d}}+d_{2} \frac{1}{n \sqrt{h^{d}}} \text { in probability. } \tag{2.21}
\end{align*}
$$

Theorem 2.3 shows that the size distortion of the proposed test is $d_{0} \sqrt{h^{d}}+d_{2} \frac{1}{n \sqrt{h^{d}}}$ when using the standard asymptotic normality in practice. A similar result has been obtained by Fan and Linton (2003). We show in addition that the bootstrap simulated critical value is approximated explicitly by $z_{\alpha}+d_{0} \sqrt{h^{d}}+d_{2} \frac{1}{n \sqrt{h^{d}}}$.

As the main objective of this paper, Section 3 below proposes a suitable selection criterion for the choice of $h$ such that while the size function is appropriately controlled, the power function is maximized at such $h$. A closed-form expression of the power function-based optimal bandwidth is given.

## 3. Power function-based bandwidth choice

We now employ the Edgeworth expansions established in Section 2 to choose a suitable bandwidth such that the power function $\beta_{n}(h)$ is maximized while the size function $\alpha_{n}(h)$ is controlled by a significance level. We thus define

$$
\begin{equation*}
h_{\mathrm{ew}}=\arg \max _{h \in H_{n}(\alpha)} \beta_{n}(h) \text { with } H_{n}(\alpha)=\left\{h: \alpha-c_{\min }<\alpha_{n}(h)<\alpha+c_{\mathrm{min}}\right\} \tag{3.1}
\end{equation*}
$$

for some arbitrarily small $c_{\text {min }}>0$.
We now start to discuss how to solve the optimization problem (3.1). It follows from (2.13) and (2.19) that

$$
\begin{equation*}
\kappa_{n}=\frac{\sqrt{h^{d}}\left(\frac{\mu_{3}^{2} K^{2}(0)}{n h^{d}}+\frac{4 \mu_{2}^{3} \nu_{3}}{3} K^{(3)}(0)\right)}{\sigma_{n}^{3}}=c_{1} \sqrt{h^{d}}+c_{2} \frac{1}{n \sqrt{h^{d}}} . \tag{3.2}
\end{equation*}
$$

Let $x=\sqrt{h^{d}}$. We rewrite $\kappa_{n}$ as $\kappa_{n}=c_{1} x+c_{2} n^{-1} x^{-1}$. Let $\gamma_{n}=\left(z_{\alpha}^{2}-1\right) \kappa_{n}$,

$$
\begin{align*}
& l_{\alpha}-r_{n} \approx z_{\alpha}+\gamma_{n}-r_{n}=z_{\alpha}+\left(d_{1}-n C_{n}^{2}\right) x+d_{4} x^{-1} \equiv z_{\alpha}+d_{3} x+d_{4} x^{-1}  \tag{3.3}\\
& l_{\alpha}-s_{n} \approx z_{\alpha}+\gamma_{n}-s_{n} \approx z_{\alpha}+\left(d_{1}-C_{0}(m)\right) x+d_{4} x^{-1}=z_{\alpha}+d_{0} x+d_{4} x^{-1} \tag{3.4}
\end{align*}
$$

where $d_{0}=d_{1}-C_{0}(m), d_{1}=\left(z_{\alpha}^{2}-1\right) c_{1}, d_{3}=d_{1}-n C_{n}^{2}$ and $d_{4}=c_{2}\left(z_{\alpha}^{2}-1\right) n^{-1}$. Note that $\lim _{n \rightarrow \infty} d_{4}=0$. Since Assumption A. 3 implies that $\lim _{n \rightarrow \infty} n C_{n}^{2}=+\infty$, we thus have

$$
\begin{equation*}
\lim _{n \rightarrow \infty} d_{3}=-\infty \quad \text { when } \quad \lim _{n \rightarrow \infty} n C_{n}^{2}=+\infty \tag{3.5}
\end{equation*}
$$

Due to this, we treat $d_{3}$ as a sufficiently large negative value when $n C_{n}^{2}$ is viewed as a sufficiently large positive value in the finite-sample analysis of this section.

Ignoring the higher-order terms (i.e. terms of order $o\left(x+n^{-1} x^{-1}\right)$ or smaller), we now re-write the power and size functions $\beta_{n}(h)$ and $\alpha_{n}(h)$ simply as functions of $x=\sqrt{h^{d}}$ as follows:

$$
\begin{align*}
\beta_{n}(h) & \approx 1-\Phi\left(l_{\alpha}-r_{n}\right)-\kappa_{n}\left(1-\left(l_{\alpha}-r_{n}\right)^{2}\right) \phi\left(l_{\alpha}-r_{n}\right) \\
& \approx 1-\Phi\left(z_{\alpha}+d_{3} x+d_{4} x^{-1}\right)-\left(c_{1} x+c_{2} n^{-1} x^{-1}\right) \\
& \times\left(1-\left(z_{\alpha}+d_{3} x+d_{4} x^{-1}\right)^{2}\right) \phi\left(z_{\alpha}+d_{3} x+d_{4} x^{-1}\right) \equiv \beta(x),  \tag{3.6}\\
\alpha_{n}(h) & \approx 1-\Phi\left(l_{\alpha}-s_{n}\right)-\kappa_{n}\left(1-\left(l_{\alpha}-s_{n}\right)^{2}\right) \phi\left(l_{\alpha}-s_{n}\right)
\end{align*}
$$

$$
\begin{align*}
& \approx 1-\Phi\left(z_{\alpha}+d_{0} x+d_{4} x^{-1}\right)-\left(c_{1} x+c_{2} n^{-1} x^{-1}\right) \\
& \times\left(1-\left(z_{\alpha}+d_{0} x+d_{4} x^{-1}\right)^{2}\right) \phi\left(z_{\alpha}+d_{0} x+d_{4} x^{-1}\right) \equiv \alpha(x) . \tag{3.7}
\end{align*}
$$

Our objective is then to find $x_{\text {ew }}=\sqrt{h_{\text {ew }}^{d}}$ such that

$$
\begin{equation*}
x_{\mathrm{ew}}=\arg \max _{x \in H_{n}(\alpha)} \beta(x) \quad \text { with } \quad H_{n}(\alpha)=\left\{x: \alpha-c_{\min }<\alpha(x)<\alpha+c_{\mathrm{min}}\right\}, \tag{3.8}
\end{equation*}
$$

where $c_{\min }$ is chosen as $c_{\min }=\frac{\alpha}{10}$ for example. Finding roots of $\beta^{\prime}(x)=0$ implies that the leading order of the unique real root of the equation is given approximately by

$$
\begin{equation*}
h_{\mathrm{ew}}=x_{\mathrm{ew}}^{\frac{2}{d}}=a_{1}^{-\frac{1}{2 d}} t_{n}^{-\frac{3}{2 d}}, \tag{3.9}
\end{equation*}
$$

where $t_{n}=n C_{n}^{2}, a_{1}=\frac{\sqrt{2} K^{(3)}(0)}{3\left(\sqrt{\int K^{2}(u) d u}\right)^{3}} \cdot c(\pi)$ with $c(\pi)=\frac{\int \pi^{3}(x) d x}{\left(\sqrt{\int \pi^{2}(x) d x}\right)^{3}}$, in which $C_{n}^{2}$ is as defined in Theorem 2.2(ii).

It can also be shown that $h_{\text {ew }}$ is the maximizer of the power function $\beta_{n}(h)$ at $h=h_{\text {ew }}$ such that

$$
\begin{equation*}
\left.\beta_{n}^{\prime \prime}(x)\right|_{x=\sqrt{h_{\mathrm{ew}}^{d}}}<0, \tag{3.10}
\end{equation*}
$$

at least for sufficiently large $n$. Detailed derivations of (3.9) and (3.10) are given in the supplemental material ${ }^{2}$.

Furthermore, the choice of $h_{\text {ew }}$ satisfies both Assumptions A.1(v) and A. 3 that

$$
\lim _{n \rightarrow \infty} n h_{\mathrm{ew}}^{d}=+\infty \quad \text { and } \quad \lim _{n \rightarrow \infty} n \sqrt{h_{\mathrm{ew}}^{d}} C_{n}^{2}=+\infty
$$

This implies that the choice of $h_{\mathrm{ew}}$ is valid to ensure $\lim _{n \rightarrow \infty} \beta_{n}\left(h_{\mathrm{ew}}\right)=1$.
When both $\sigma^{2}=\mu_{2}=E\left[e_{1}^{2}\right]$ and the marginal density function $\pi(\cdot)$ of $\left\{X_{i}\right\}$ are unknown in practice, we propose using an estimated version of $h_{\text {ew }}$ as follows:

$$
\begin{equation*}
\widehat{h}_{\mathrm{ew}}=\widehat{a}_{1}^{-\frac{1}{2 d}} \widehat{t}_{n}^{-\frac{3}{2 d}} \tag{3.11}
\end{equation*}
$$

[^1]where
\[

$$
\begin{aligned}
& \widehat{t}_{n}=n \widehat{C}_{n}^{2} \text { with } \widehat{C}_{n}^{2}=\frac{\frac{1}{n} \sum_{i=1}^{n} \widehat{\Delta}_{n}^{2}\left(X_{i}\right) \widehat{\pi}\left(X_{i}\right)}{\widehat{\mu}_{2} \sqrt{2 \widehat{\nu}_{2} \int K^{2}(v) d v} \text { and }} \\
& \widehat{a}_{1}=\frac{\sqrt{2} K^{(3)}(0)}{3\left(\sqrt{\int K^{2}(u) d u}\right)^{3}} \widehat{c}(\pi) \text { with } \widehat{c}(\pi)=\frac{\frac{1}{n} \sum_{i=1}^{n} \widehat{\pi}^{2}\left(X_{i}\right)}{\left(\sqrt{\frac{1}{n} \sum_{i=1}^{n} \widehat{\pi}\left(X_{i}\right)}\right)^{3}},
\end{aligned}
$$
\]

in which $\widehat{\mu}_{2}, \widehat{\nu}_{2}$ and $\widehat{\pi}(\cdot)$ are as defined in (2.7), and $\widehat{\Delta}_{n}(x)$ is given by

$$
\widehat{\Delta}_{n}(x)=\frac{\sum_{i=1}^{n} K\left(\frac{x-X_{i}}{\hat{b}_{\mathrm{cv}}}\right)\left(Y_{i}-m_{\widehat{\theta}}\left(X_{i}\right)\right)}{\sum_{i=1}^{n} K\left(\frac{x-X_{i}}{\hat{b}_{\mathrm{cv}}}\right)}
$$

with $\widehat{\theta}$ and $\widehat{b}_{\text {cv }}$ being the same as in (2.7).
Note also that $\widehat{h}_{\text {ew }}$ provides an optimal bandwidth irrespectively of whether one works under the null hypothesis $\mathcal{H}_{0}$ or under the alternative hypothesis $\mathcal{H}_{1}$. In other words, it can be used for computing not only the power under an alternative $\mathcal{H}_{1}$, but also the size under $\mathcal{H}_{0}$ in each case. Detailed discussion about this is given in the supplemental material ${ }^{3}$.

We conclude this section by summarizing the above discussion into the following proposition; its proof is given in the supplemental material ${ }^{4}$.

Proposition 3.1. Suppose that Assumptions A.1-A.3 listed in the appendix hold. Additionally, suppose that $\Delta_{n}(x)$ is continuously differentiable such that

$$
\lim _{n \rightarrow \infty} \sup _{x \in D_{\pi}} \frac{\left\|\Delta_{n}^{\prime}(x)\right\|}{\left|\Delta_{n}(x)\right|} \leq C<\infty \quad \text { and } \lim _{n \rightarrow \infty} \inf _{x \in \mathbb{R}^{d}}\left|\Delta_{n}(x)\right| \sqrt{n \widehat{b}_{\mathrm{cv}}^{d}}=\infty \text { in probability }
$$

for some $C>0$, where $D_{\pi}=\left\{x \in \mathbb{R}^{d}: \pi(x)>0\right\}$ and $\|\cdot\|^{2}$ denotes the Euclidean norm. Then

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \frac{\beta_{n}\left(\widehat{h}_{\mathrm{ew}}\right)}{\beta_{n}\left(h_{\mathrm{ew}}\right)}=1 \quad \text { in probability. } \tag{3.12}
\end{equation*}
$$

As pointed out in the introduction, implementation of each of existing nonparametric kernel tests involves either a single bandwidth chosen optimally for estimation

[^2]purposes or a set of bandwidth values. The proposed $\widehat{h}_{\text {ew }}$ is chosen optimally for testing purposes. Section 4 below shows how to implement the proposed test based on our bandwidth in practice and compares the finite-sample performance of the proposed choice with that of some closely relevant alternatives in the literature.

## 4. Examples of implementation

This section presents two simulated examples and one real data example to illustrate the proposed theory and methods in Sections 2 and 3 as well as to make comparisons with some closely relevant alternatives in the literature. Simulated example 4.1 below discusses the finite-sample performance of the proposed test $\widehat{T}_{n}\left(\widehat{h}_{\text {ew }}\right)$ with that of the alternative version where the test is coupled with a cross-validation (CV) bandwidth choice. Simulated example 4.2 below compares our test with some of the commonly used tests in the literature. Example 4.3 provides a real data example to show that the proposed test makes a clear difference. In the following finite-sample study in Examples 4.1-4.3 below, we consider the case where $\Delta_{n}(x)=c_{n} \Delta(x)$, in which $\left\{c_{n}\right\}$ is a sequence of positive real numbers satisfying $\lim _{n \rightarrow \infty} c_{n}=0$ and $\Delta(x)$ is an unknown function not depending on $n$.

Example 4.1. Consider a nonparametric time series regression model of the form

$$
\begin{equation*}
Y_{i}=\theta_{1} X_{i 1}+\theta_{2} X_{i 2}+c_{n}\left(X_{i 1}^{2}+X_{i 2}^{2}\right)+e_{i}, 1 \leq i \leq n, \tag{4.1}
\end{equation*}
$$

where $\left\{e_{i}\right\}$ is a sequence of Normal errors and both $X_{i 1}$ and $X_{i 2}$ are time series variables generated by

$$
\begin{equation*}
X_{i 1}=\alpha X_{i-1,1}+u_{i} \quad \text { and } \quad X_{i 2}=\beta X_{i-1,2}+v_{i}, 1 \leq i \leq n \tag{4.2}
\end{equation*}
$$

with $\left\{u_{i}\right\}$ and $\left\{v_{i}\right\}$ being i.i.d. random errors generated independently from Normal distributions as below.

Under $\mathcal{H}_{0}$, we generate a sequence of observations $\left\{Y_{i}\right\}$ with $\theta_{1}=\theta_{2}=1$ as the true parameters, i.e.

$$
\begin{equation*}
\mathcal{H}_{0}: \quad Y_{i}=X_{i 1}+X_{i 2}+e_{i}, \tag{4.3}
\end{equation*}
$$

where $\left\{e_{i}\right\}$ is a sequence of independent and identically distributed random errors generated from $N(0,1)$, and $\left\{X_{i 1}\right\}$ and $\left\{X_{i 2}\right\}$ are independently generated from

$$
\begin{equation*}
X_{i 1}=0.5 X_{i-1,1}+u_{i} \text { and } X_{i 2}=0.5 X_{i-1,2}+v_{i}, 1 \leq i \leq n \tag{4.4}
\end{equation*}
$$

with $X_{01}=X_{02}=0$ and $\left\{u_{i}\right\}$ and $\left\{v_{i}\right\}$ are sequences of independent and identically distributed random errors and generated independently from a $N(0,1)$.

Under $\mathcal{H}_{1}$, we are interested in two alternative models of the form

$$
\begin{equation*}
\mathcal{H}_{1}: Y_{i}=X_{i 1}+X_{i 2}+c_{n}\left(X_{i 1}^{2}+X_{i 2}^{2}\right)+e_{i}, \quad e_{i} \sim N(0,1) \tag{4.5}
\end{equation*}
$$

with $c_{n}$ being chosen as either $c_{1 n}=n^{-\frac{1}{2}} \sqrt{\log \log (n)}$ or $c_{2 n}=n^{-\frac{7}{18}}$.
In the testing procedure, the parameters $\theta_{1}$ and $\theta_{2}$ in the parametric model are estimated as discussed in Sections 1 and 2.

The reasoning for the above choice of $c_{j n}$ is as follows. The rate of $c_{1 n}=$ $n^{-\frac{1}{2}} \sqrt{\log \log (n)}$ should be an optimal rate of testing in this kind of nonparametric kernel testing problem as discussed in Horowitz and Spokoiny (2001). The rate of $c_{2 n}=n^{-\frac{7}{18}}$ implies that the optimal bandwidth $\widehat{h}_{\text {ew }}$ in (3.11) with $d=2$ is proportional to $n^{-\frac{1}{6}}$.

Throughout this example, we choose $K(\cdot)$ as the standard normal density function. Let $\widehat{h}_{\mathrm{cv}}$ be chosen by a cross-validation criterion of the form

$$
\begin{equation*}
\widehat{h}_{\mathrm{cv}}=\arg \min _{h \in H_{\mathrm{cv}}} \frac{1}{n} \sum_{i=1}^{n}\left(Y_{i}-\widehat{m}_{-i}\left(X_{i 1}, X_{i 2} ; h\right)\right)^{2} \text { with } H_{\mathrm{cv}}=\left[n^{-1}, n^{\frac{1}{6}}\right] \tag{4.6}
\end{equation*}
$$

in which

$$
\widehat{m}_{-i}\left(X_{i 1}, X_{i 2} ; h\right)=\frac{\sum_{l=1, \neq i}^{n} K\left(\frac{X_{l 1}-X_{i 1}}{h}\right) K\left(\frac{X_{l 2}-X_{i 2}}{h}\right) Y_{l}}{\sum_{l=1, \neq i}^{n} K\left(\frac{X_{l 1}-X_{i 1}}{h}\right) K\left(\frac{X_{l 2}-X_{i 2}}{h}\right)} .
$$

Let $\widehat{h}_{0 \text { test }}$ be the corresponding version of $\widehat{h}_{\text {ew }}$ in (3.11) and $\widehat{h}_{0 \text { cv }}$ be the corresponding version of $\widehat{h}_{\mathrm{cv}}$ in (4.6) both computed under $\mathcal{H}_{0}$. Since $\left\{Y_{i}\right\}$ under $\mathcal{H}_{1}$ depends on the choice of $c_{n}$, thus the computing of both $\widehat{h}_{\text {ew }}$ of (3.11) and $\widehat{h}_{\text {cv }}$ of (4.6) under $\mathcal{H}_{1}$ depend on the choice of $c_{n}$. Let $\widehat{h}_{j \text { test }}$ be the corresponding versions of $\widehat{h}_{\text {ew }}$ in (3.11) and $\widehat{h}_{j \mathrm{cv}}$ be the corresponding versions of $\widehat{h}_{\mathrm{cv}}$ in (4.6) with $c_{n}=c_{j n}$ for $j=1,2$.

In order to compare the size and power properties of $\widehat{T}_{n}(h)$ with the most relevant alternatives, we introduce the following simplified notation: for $j=1,2$,

$$
\begin{aligned}
& \alpha_{01}=P\left(\widehat{T}_{n}\left(\widehat{h}_{0 \mathrm{cv}}\right)>l_{\alpha}^{*}\left(\widehat{h}_{0 \mathrm{cv}}\right) \mid \mathcal{H}_{0}\right), \quad \beta_{j 1}=P\left(\widehat{T}_{n}\left(\widehat{h}_{j \text { cv }}\right)>l_{\alpha}^{*}\left(\widehat{h}_{0 \mathrm{cv}}\right) \mid \mathcal{H}_{1}\right), \\
& \alpha_{02}=P\left(\widehat{T}_{n}\left(\widehat{h}_{0 \text { test }}\right)>l_{\alpha}^{*}\left(\widehat{h}_{0 \text { test }}\right) \mid \mathcal{H}_{0}\right), \quad \beta_{j 2}=P\left(\widehat{T}_{n}\left(\widehat{h}_{j \text { test }}\right)>l_{\alpha}^{*}\left(\widehat{h}_{0 \text { test }}\right) \mid \mathcal{H}_{1}\right) .
\end{aligned}
$$

We consider cases where the number of replications of each of the sample versions of $\alpha_{0 k}$ and $\beta_{j k}$ for $j, k=1,2$ was $M=1000$, each with $B=250$ number of bootstrapping resamples, and the simulations were done for the cases of $n=250,500$ and 750. The detailed results at the $1 \%, 5 \%$ and $10 \%$ significance level are given in Tables 4.1-4.3, respectively.

Table 4.1. Simulated size and power values at the $1 \%$ significance level

| Sample Size | Null Hypothesis Is True |  | Null Hypothesis Is False |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | $\alpha_{01}$ | $\alpha_{02}$ | $\beta_{11}$ | $\beta_{21}$ | $\beta_{12}$ | $\beta_{22}$ |
| 250 | 0.012 | 0.016 | 0.212 | 0.239 | 0.294 | 0.272 |
| 500 | 0.018 | 0.014 | 0.270 | 0.303 | 0.318 | 0.334 |
| 750 | 0.014 | 0.008 | 0.310 | 0.367 | 0.408 | 0.422 |

Table 4.2. Simulated size and power values at the $5 \%$ significance level

| Sample Size | Null Hypothesis Is True |  | Null Hypothesis Is False |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | $\alpha_{01}$ | $\alpha_{02}$ | $\beta_{11}$ | $\beta_{21}$ | $\beta_{12}$ | $\beta_{22}$ |
| 250 | 0.054 | 0.046 | 0.514 | 0.522 | 0.656 | 0.658 |
| 500 | 0.052 | 0.058 | 0.572 | 0.564 | 0.690 | 0.730 |
| 750 | 0.046 | 0.052 | 0.648 | 0.658 | 0.820 | 0.812 |

Table 4.3. Simulated size and power values at the $10 \%$ significance level

| Sample Size | Null Hypothesis Is True |  | Null Hypothesis Is False |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | $\alpha_{01}$ | $\alpha_{02}$ | $\beta_{11}$ | $\beta_{21}$ | $\beta_{12}$ | $\beta_{22}$ |
| 250 | 0.116 | 0.110 | 0.696 | 0.764 | 0.884 | 0.909 |
| 500 | 0.104 | 0.090 | 0.744 | 0.817 | 0.860 | 0.934 |
| 750 | 0.108 | 0.090 | 0.844 | 0.895 | 0.946 | 0.968 |

Tables 4.1-4.3 report comprehensive simulation results for both the sizes and power values of the proposed tests for models (4.3) and (4.4). Column 2 in each of Tables 4.1-4.3 shows that while the sizes for the test based on $\widehat{h}_{0 \text { cv }}$ are comparable with these given in column 3 based on $\widehat{h}_{0 \text { test }}$, the power values of the test based on $\widehat{h}_{j \text { test }}$ in columns 6 and 7 are always greater than these given in columns 4 and 5 based on $\widehat{h}_{j \text { cv }}$. This is not surprising, because the theory shows that each of $\widehat{h}_{j \text { test }}$ is chosen such that the resulting power function is maximized while the corresponding size function is under control by the significance level.

In addition, the test based on $\widehat{h}_{2 \text { test }}$ is almost uniformly more powerful than the best based on $\widehat{h}_{1 \text { test }}$, which is the second most powerful test. This is basically because $\widehat{h}_{2 \text { test }}$ is based on considering $\mathcal{H}_{1}$ with $c_{2 n}=n^{-\frac{7}{18}}$, which goes to zero slower than $c_{1 n}=n^{-\frac{1}{2}} \sqrt{\log \log (n)}$, and hence the distance between the alternative and the null is biggest in the former case (and therefore easier to detect). Meanwhile, the last columns of Tables 4.1-4.3 show that the test based on the bandwidth $\widehat{h}_{2 \text { test }}$ is still a powerful test even though the bandwidth is proportional to $n^{-\frac{1}{6}}$, which is the same as the optimal bandwidth based on a cross-validation estimation method. This shows that whether an estimation-based optimal bandwidth may be used for testing depends on whether the bandwidth is chosen optimally for testing purposes.

We finally want to stress that the proposed test based on either $\widehat{h}_{1 \text { test }}$ or $\widehat{h}_{2 \text { test }}$ has not only stable sizes even at a small sample size of $n=250$, but also reasonable power values even when the 'distance' between the null and the alternative has been made deliberately close at the rate of $\sqrt{n^{-1} \log \log (n)}=0.060$ for $n=500$ for example. We
can expect that the test would have bigger power values when the 'distance' is made wider. Overall, Tables 4.1-4.3 show that the established theory and methodology is workable in the small and medium-sample case.

Example 4.1 discusses the small and medium-sample comparison results for the proposed test with either testing-based optimal bandwidth or estimation-based (CV) bandwidth. Example 4.2 below considers comparing the small and medium-sample performance of the proposed test associated with the optimal bandwidth with some closely related nonparametric tests available in both the econometrics and statistics literature.

Example 4.2. Consider a linear model of the form

$$
\begin{equation*}
Y_{i}=\alpha_{0}+\beta_{0} X_{i}+e_{i}, 1 \leq i \leq n=250, \tag{4.7}
\end{equation*}
$$

where $\left\{X_{i}\right\}$ is a sequence of independent random variables sampled from $N(0,25)$ distribution truncated at its 5 th and 95 th percentiles, and $\left\{e_{i}\right\}$ is sampled from one of the three distributions: (i) $e_{i} \sim N(0,4)$; (ii) a mixture of Normals in which $\left\{e_{i}\right\}$ is sampled from $N(0,1.56)$ with probability 0.9 and from $N(0,25)$ with probability 0.1 ; and (iii) the Type I extreme value distribution scaled to have a variance of 4 . The mixture distribution is leptokurtic with a variance of 0.39 , and the Type I extreme value distribution is asymmetrical.

This is the same example as used in Horowitz and Spokoiny (2001) for the comparison with some of the commonly used tests in the literature, such as the Andrews' test proposed in Andrews (1997), the HM test proposed in Härdle and Mammen (1993), the HS test proposed in Horowitz and Spokoiny (2001) and the empirical likelihood (EL) test proposed in Chen, Härdle and Li (2003).

To compute the sizes of the test, choose $\alpha_{0}=\beta_{0}=1$ as the true parameters and then generate $\left\{Y_{i}\right\}$ from $Y_{i}=1+X_{i}+e_{i}$ under $\mathcal{H}_{0}$, and generate $\left\{Y_{i}\right\}$ from $Y_{i}=1+X_{i}+\frac{5}{\tau} \phi\left(\frac{X_{i}}{\tau}\right)+e_{i}$ under $\mathcal{H}_{1}$, where $\tau=1$ or 0.25 , and $\phi(\cdot)$ is the density function of the standard normal distribution.

The kernel function used here is $K(x)=\frac{15}{16}\left(1-x^{2}\right)^{2} I(|x| \leq 1)$. Choose $c_{n}=5 \tau^{-1}$ and $\Delta(x)=\phi\left(x \tau^{-1}\right)$ for the corresponding forms in (1.2). For $j=1,2$, let $c_{j n}=5 \tau_{j}^{-1}$ and $\Delta_{j}(x)=\phi\left(x \tau_{j}^{-1}\right)$ with $\tau_{1}=1$ and $\tau_{2}=0.25$. Let $\widehat{h}_{i \text { new }}$ be the corresponding version of $\widehat{h}_{\mathrm{ew}}$ of (3.11) based on $\left(c_{j n}, \Delta_{j}(x)\right)$ for $j=1,2$.

In order to make a fair comparison, we use the same number of the bootstrap resamples of $M=99$, the same number of replications of $M=1000$ under $\mathcal{H}_{0}$ and $M=250$ under $\mathcal{H}_{1}$ as in Table 1 of Horowitz and Spokoiny (2001). In Table 4.4 below, we add the size and power values to the last two columns for both the EL test and the proposed test- $\widehat{T}_{n}\left(\widehat{h}_{\text {inew }}\right)$ of this paper. The other parts of the table are obtained and tabulated similarly to Table 1 of Horowitz and Spokoiny (2001).

Table 4.4. Simulated size and power values at the $5 \%$ significance level

|  |  | Probability of Rejecting Null Hypothesis |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | Andrews | HM | HS | EL | $\widehat{T}_{n}\left(\widehat{h}_{\text {new }}\right)$ |  |
| Distribution | $\tau$ | Test | Test | Test | Test | Test |  |
|  |  | Null Hypothesis Is True |  |  |  |  |  |
| Normal |  | 0.057 | 0.060 | 0.066 | 0.053 | 0.049 |  |
| Mixture |  | 0.053 | 0.053 | 0.048 | 0.055 | 0.052 |  |
| Extreme |  | 0.063 | 0.057 | 0.055 | 0.057 | 0.052 |  |
|  |  | Null Hypothesis Is False |  |  |  |  |  |
| Normal | 1.0 | 0.680 | 0.752 | 0.792 | 0.900 | 0.907 |  |
| Mixture | 1.0 | 0.692 | 0.736 | 0.835 | 0.905 | 1.000 |  |
| Extreme | 1.0 | 0.600 | 0.760 | 0.820 | 0.924 | 0.935 |  |
| Normal | 0.25 | 0.536 | 0.770 | 0.924 | 0.929 | 0.993 |  |
| Mixture | 0.25 | 0.592 | 0.704 | 0.922 | 0.986 | 0.999 |  |
| Extreme | 0.25 | 0.604 | 0.696 | 0.968 | 0.989 | 0.989 |  |

Table 4.4 shows that the proposed test has better power properties than any of
the commonly used tests, while the size values are comparable with those of the competitors. The results further support the power-based bandwidth selection procedure proposed in Sections 2 and 3.

As discussed in the supplemental material, the proposed theory and methodology for model (1.1) can be applied to an extended model of the form

$$
\begin{equation*}
Y_{i}=m\left(X_{i}\right)+e_{i} \quad \text { with } \quad e_{i}=\sigma\left(X_{i}\right) \epsilon_{i}, 1 \leq i \leq n, \tag{4.8}
\end{equation*}
$$

where $\sigma(\cdot)$ satisfying $\inf _{x \in \mathbb{R}^{d}} \sigma(x)>0$ is unknown nonparametrically and $\left\{\epsilon_{i}\right\}$ is a sequence of i.i.d. random errors with zero mean and finite variance. In addition, $\left\{\epsilon_{i}\right\}$ and $\left\{X_{j}\right\}$ are assumed to be independent for all $1 \leq j \leq i \leq n$. A special case of model (4.8) is discussed in Example 4.3 below.

Example 4.3. This example examines the high frequency seven-day Eurodollar deposit rate sampled daily from 1 June 1973 to 25 February 1995. This provides us with $n=5505$ observations. Let $\left\{X_{i}: i=1,2, \cdots, n=5505\right\}$ be the set of Eurodollar deposit rate data. Figures 4.1 and 4.2 below plot the data values and the conventional nonparametric kernel density estimator

$$
\widehat{\pi}(x)=\frac{1}{n \widetilde{h}_{\mathrm{cv}}} \sum_{i=1}^{n} K\left(\frac{x-X_{i}}{\widetilde{h}_{\mathrm{cv}}}\right)
$$

respectively, where $K(x)=\frac{1}{\sqrt{2 \pi}} e^{-\frac{x^{2}}{2}}$ and $\widetilde{h}_{\mathrm{cv}}$ is the conventional normal-reference based bandwidth given by

$$
\begin{equation*}
\widetilde{h}_{\mathrm{cv}}=1.06 \cdot n^{-\frac{1}{5}} \sqrt{\frac{1}{n-1} \sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)^{2}} \text { with } \bar{X}=\frac{1}{n} \sum_{i=1}^{n} X_{i} \text {. } \tag{4.9}
\end{equation*}
$$

Note that $\widehat{b}_{\mathrm{cv}}$ of (2.7), $\widehat{h}_{\mathrm{cv}}$ of (4.6) and $\widetilde{h}_{\mathrm{cv}}$ of (4.9) are normally different from each other. In the case where $\left\{X_{i}\right\}$ follows an autoregressive model, they can be chosen the same. Thus, they are chosen the same in this example.

It has been assumed in the literature (see, for example, Ait-Sahalia 1996; Fan and Zheng 2003; Arapis and Gao 2006) that the Eurodollar data set $\left\{X_{i}\right\}$ may be


Figure 4.1: Seven-day Eurodollar deposit rate, June 1, 1973 to February 25, 1995.


Figure 4.2: Nonparametric kernel density estimator of the Eurodollar rate.
modeled by a nonlinear time series model of the form

$$
\begin{equation*}
Y_{i}=\mu\left(X_{i}\right)+\sigma\left(X_{i}\right) \epsilon_{i}, 1 \leq i \leq n, \tag{4.10}
\end{equation*}
$$

where $Y_{i}=\frac{X_{i+1}-X_{i}}{\Lambda}, \sigma(\cdot)>0$ is unknown nonparametrically, and $\epsilon_{i} \sim N\left(0, \Lambda^{-1}\right)$, in
which $\Lambda$ is the time between successive observations. Since we consider a daily data set, this gives $\Lambda=\frac{1}{250}$.

On the question of whether there is any nonlinearity in the drift function $\mu(\cdot)$, existing studies have provided no definitive answer. For example, Aït-Sahalia (1996), and Arapis and Gao (2006) show that there is some evidence of supporting nonlinearity in the drift on the one hand. On the other hand, existing studies, such as Chapman and Pearson (2000), and Fan and Zheng (2003), suggest that nonlinearity may just be caused by estimation biases when using nonparametric kernel estimation.

To further discuss whether the assumption on linearity in the drift is appropriate for the given set of data, we apply our test to propose testing

$$
\begin{equation*}
\mathcal{H}_{01}: \mu(x)=\mu\left(x ; \theta_{0}\right)=\beta_{0}\left(\alpha_{0}-x\right) \text { versus } \mathcal{H}_{11}: \mu(x)=\beta_{1}\left(\alpha_{1}-x\right)+c_{n} \Delta(x) \tag{4.11}
\end{equation*}
$$

for some $\theta_{j}=\left(\alpha_{j}, \beta_{j}\right) \in \Theta$ for $j=0,1$ and $c_{n}=\sqrt{n^{-1} \log \log (n)}$, where $\Theta$ is a parameter space in $\mathbb{R}^{2}$ and $\Delta(x)$ is a continuous function.

It can be shown that the proposed test in Section 2 has an asymptotically equivalent version of the form:

$$
\begin{equation*}
\widetilde{T}_{n}(h)=\frac{\sum_{j=1}^{n} \sum_{i=1, i \neq j}^{n} \widehat{e}_{j} K\left(\frac{X_{i}-X_{j}}{h}\right) \widehat{e}_{i}}{\sqrt{2 \sum_{j=1}^{n} \sum_{i=1}^{n} \widehat{e}_{j}^{2} K^{2}\left(\frac{X_{i}-X_{j}}{h}\right) \widehat{e}_{i}^{2}}}, \tag{4.12}
\end{equation*}
$$

where $\widehat{e}_{i}=Y_{i}-\widehat{\beta}\left(\widehat{\alpha}-X_{i}\right)$, in which $(\widehat{\alpha}, \widehat{\beta})$ is the pair of the conventional least squares estimators minimizing $\sum_{i=1}^{n}\left(Y_{i}-\widehat{\beta}\left(\widehat{\alpha}-X_{i}\right)\right)^{2}$.

As pointed out in the literature (Arapis and Gao 2006), $\widetilde{T}_{n}(h)$ is independent of the structure of the conditional variance $\sigma^{2}(\cdot)$. The kernel function used is the standard normal density function given by $K(x)=\frac{1}{\sqrt{2 \pi}} e^{-\frac{x^{2}}{2}}$.

Let $\widetilde{h}_{\text {test }}$ be the corresponding version of (3.11). It has been shown in the supplementary document ${ }^{5}$ that

$$
\begin{equation*}
\widetilde{h}_{\text {test }}=\widehat{a}_{1}^{-\frac{1}{2}} \widehat{t}_{n}^{\frac{3}{2}} \tag{4.13}
\end{equation*}
$$

[^3]where $\widehat{t}_{n}$ and $\widehat{a}_{1}$ are the same as in (3.11), in which $\widehat{c}(\pi)$ becomes
\[

$$
\begin{equation*}
\widehat{c}(\pi)=\frac{1}{n} \sum_{i=1}^{n} \widehat{\pi}^{2}\left(X_{i}\right) \widehat{\sigma}^{6}\left(X_{i}\right) \cdot\left(\frac{1}{n} \sum_{i=1}^{n} \widehat{\pi}\left(X_{i}\right) \widehat{\sigma}^{4}\left(X_{i}\right)\right)^{-\frac{3}{2}} \tag{4.14}
\end{equation*}
$$

\]

with

$$
\widehat{\sigma}^{2}\left(X_{i}\right)=\frac{\sum_{u=1}^{n} K\left(\frac{X_{i}-X_{u}}{\widehat{h}_{\mathrm{cv}}}\right) \widehat{e}_{u}^{2}}{\sum_{v=1}^{n} K\left(\frac{X_{i}-X_{v}}{\widehat{h}_{\mathrm{cv}}}\right)} .
$$

Let $L_{1}=\widetilde{T}_{n}\left(\widetilde{h}_{\text {test }}\right)$ and $L_{2}=\widetilde{T}_{n}\left(\widetilde{h}_{\mathrm{cv}}\right)$. To apply the test $L_{j}$ for each $j=1,2$ to test $\mathcal{H}_{01}$, we propose the following procedure for computing the $p$-value of $L_{j}$ :

- Compute $\widehat{e}_{i}=Y_{i}-\widehat{\beta}\left(\widehat{\alpha}-X_{i}\right)$ and then generate a sequence of bootstrap resamples $\left\{\widehat{e}_{i}^{*}\right\}$ given by $\widehat{e}_{i}^{*}=\widehat{\sigma}\left(X_{i}\right) \epsilon_{i}^{*}$, where $\left\{\epsilon_{i}^{*}\right\}$ is a sequence of i.i.d. bootstrap resamples generated from $N\left(0, \Lambda^{-1}\right)$ and $\widehat{\sigma}^{2}(\cdot)$ is defined as above.
- Generate $\widehat{Y}_{i}^{*}=\widehat{\beta}\left(\widehat{\alpha}-X_{i}\right)+\widehat{e}_{i}^{*}$. Compute the corresponding version $L_{j}^{*}$ of $L_{j}$ for each $j=1,2$ based on $\left\{\widehat{Y}_{i}^{*}\right\}$.
- Repeat the above steps $M=1000$ times to find the bootstrap distribution of $L_{j}^{*}$ and then compute the proportion that $L_{j}<L_{j}^{*}$ for each $j=1,2$. This proportion is a simulated $p$-value of $L_{j}$.

Our simulation results return the simulated $p$-values of $\widehat{p}_{1}=0.102$ for $L_{1}$ and $\widehat{p}_{2}=0.072$ for $L_{2}$. While both of the simulated $p$-values suggest that there is no enough evidence of rejecting the linearity in the drift at the $5 \%$ significance level, the evidence of accepting the linearity based on $L_{1}$ is stronger than that based on $L_{2}$.

As our test $\widetilde{T}_{n}\left(\widetilde{h}_{\text {test }}\right)$ involves no estimation biases, the process of computing the simulated $p$-values is quite robust. We therefore believe that this improved test further reinforces the findings of Chapman and Pearson (2000) and Fan and Zhang (2003) that there is no definitive answer to the question whether the short rate drift is actually nonlinear.

## 5. Conclusion

This paper has addressed the issue of how to appropriately choose the bandwidth parameter when using a nonparametric kernel-based test. Both the size and power properties of the proposed test have been studied systematically. The established theory and methodology has shown that a suitable bandwidth can be optimally chosen after appropriately balancing the size and power functions. Furthermore, the new methodology has resulted in a closed-form representation for the leading term of such an optimal bandwidth in the finite-sample case.

Existing results (see, for example, Li and Wang 1998; Li 1999; Fan and Linton 2003) show that this kind of nonparametric kernel test associated with a large sample critical value may not have good size and power properties. Our small and mediumsample studies in both the simulated and real-data examples have shown that the performance of such a test can be significantly improved when it is coupled with a power-based optimal bandwidth as well as a bootstrap simulated critical value.

It is pointed out that the established theory and methodology has various applications in providing solutions to some other related testing problems, in which nonparametric methods are involved. Future extensions also include dealing with cases where both $X_{i}$ and $e_{i}$ may be strictly stationary time series.

## Appendix

This appendix lists the necessary assumptions for the establishment and the proofs of the main results given in Section 2.

## A.1. Assumptions

Assumption A.1. (i) Assume that $\left\{e_{i}\right\}$ is a sequence of i.i.d. continuous random errors with $E\left[e_{1}\right]=0,0<\sigma^{2}=E\left[e_{1}^{2}\right]=\sigma^{2}<\infty$ and $E\left[e_{1}^{6}\right]<\infty$.
(ii) We assume that $\left\{X_{i}\right\}$ is strictly stationary and $\alpha$-mixing with mixing coefficient $\alpha(t)$ being defined by

$$
\begin{equation*}
\alpha(t)=\sup \left\{|P(A \cap B)-P(A) P(B)|: A \in \Omega_{1}^{s}, B \in \Omega_{s+t}^{\infty}\right\} \leq C_{\alpha} \alpha^{t} \tag{A.1}
\end{equation*}
$$

for all $s, t \geq 1$, where $0<C_{\alpha}<\infty$ and $0<\alpha<1$ are constants, and $\Omega_{i}^{j}$ denotes the $\sigma$-field generated by $\left\{X_{k}: i \leq k \leq j\right\}$.
(iii) We also assume that $\left\{X_{s}\right\}$ and $\left\{e_{t}\right\}$ are independent for all $1 \leq s \leq t \leq n$. Let $\pi(\cdot)$ be the marginal density such that $0<\int \pi^{3}(x) d x<\infty$, and $\pi_{\tau_{1}, \tau_{2}, \cdots, \tau_{l}}(\cdot)$ be the joint probability density of $\left(X_{1+\tau_{1}}, \ldots, X_{1+\tau_{l}}\right)(1 \leq l \leq 4)$. Assume that $\pi_{\tau_{1}, \tau_{2}, \cdots, \tau_{l}}(\cdot)$ for all $1 \leq l \leq 4$ do exist and are continuous and bounded.
(iv) Assume that the univariate kernel function $K(\cdot)$ is a symmetric and bounded probability density function. In addition, we assume the existence of both $K^{(3)}(\cdot)$, the three-time convolution of $K(\cdot)$ with itself, and $K_{2}^{(2)}(\cdot)$, the two-time convolution of $K^{2}(\cdot)$ with itself.
(v) The bandwidth parameter $h$ satisfies both $\lim _{n \rightarrow \infty} h=0$ and $\lim _{n \rightarrow \infty} n h^{d}=\infty$.

Assumption A.2. (i) Let $\mathcal{H}_{0}$ be true. Then for any sufficiently small $\varepsilon_{1}>0$ and some $B_{1 L}>0$

$$
\lim _{n \rightarrow \infty} P\left(\sqrt{n}\left\|\hat{\theta}-\theta_{0}\right\|>B_{1 L}\right)<\varepsilon_{1}
$$

where $\theta_{0}$ is the same as defined in (1.2).
(ii) Let $\mathcal{H}_{1}$ be true. Then for any sufficiently small $\varepsilon_{2}>0$ and some $B_{2 L}>0$

$$
\lim _{n \rightarrow \infty} P\left(\sqrt{n}\left\|\widehat{\theta}-\theta_{1}\right\|>B_{2 L}\right)<\varepsilon_{2}
$$

where $\theta_{1}$ is the same as defined in (1.2).
(iii) There exist some absolute constants $\varepsilon_{3}>0$ and $0<B_{3 L}<\infty$ such that the following

$$
\lim _{n \rightarrow \infty} P\left(\sqrt{n}\left\|\widehat{\theta}^{*}-\widehat{\theta}\right\|>B_{3 L} \mid \mathcal{W}_{n}\right)<\varepsilon_{3}
$$

holds in probability, where $\widehat{\theta}^{*}$ is as defined in the Simulation Procedure above Theorem 2.1.
(iv) Let $m_{\theta}(x)$ be differentiable with respect to $\theta$ and $\frac{\partial m_{\theta}(x)}{\partial \theta}$ be continuous in both $x$ and $\theta$. In addition, $E\left[\left(\frac{m_{\theta_{0}}\left(X_{1}\right)}{\partial \theta}\right)\left(\frac{m_{\theta_{0}}\left(X_{1}\right)}{\partial \theta}\right)^{\tau}\right]$ is a positive definite matrix, and

$$
0<\int\left\|\left.\frac{\partial m_{\theta}(x)}{\partial \theta}\right|_{\theta=\theta_{0}}\right\|^{2} \pi^{2}(x) d x<\infty
$$

AsSumption A.3. (i) Let $\left\{\Delta_{n}(x)\right\}$ be a sequence of continuous functions such that $0<\int \Delta_{n}^{2}(x) \pi^{2}(x) d x<\infty$.
(ii) Let $C_{n}^{2}$ satisfy $\lim _{n \rightarrow \infty} n \sqrt{h^{d}} C_{n}^{2}=\infty$ and $\lim _{n \rightarrow \infty} n C_{n}^{6}=0$, where

$$
C_{n}^{2}=\frac{\int \Delta_{n}^{2}(x) \pi^{2}(x) d x}{\sigma^{2} \sqrt{2 \nu_{2} \int K^{2}(v) d v}}
$$

in which $\nu_{2}=E\left[\pi^{2}\left(X_{1}\right)\right]<\infty$.
Assumptions A.1-A. 3 are standard and justifiable conditions. Some detailed justifications are given in the supplemental material ${ }^{6}$.

## A.2. Technical lemmas

Recall that using $\lim _{n \rightarrow \infty} n h^{d}=\infty$

$$
\begin{align*}
\kappa_{n} & =\frac{\sqrt{h^{d}}\left(\frac{\mu_{3}^{2} K^{2}(0)}{n h^{d}}+\frac{4 \mu_{2}^{3} \nu_{3}}{3} K^{(3)}(0)\right)}{\sigma_{n}^{3}} \equiv c_{1} \sqrt{h^{d}}+c_{2} \frac{1}{n \sqrt{h^{d}}} \\
& =c_{1} \sqrt{h^{d}}\left(1+c_{2} c_{1}^{-1} \frac{1}{n h^{d}}\right) \approx c_{1} \sqrt{h^{d}} . \tag{A.2}
\end{align*}
$$

In order to establish some useful lemmas without including non-essential technicality, we introduce the following simplified notation:

$$
\begin{align*}
a_{i j} & =\frac{1}{n \sqrt{h^{d}} \sigma_{n}} K\left(\frac{X_{i}-X_{j}}{h}\right), \quad L_{n}(h)=\sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} a_{i j} e_{i} e_{j}, \\
\rho(h) & =\frac{\sqrt{2} K^{(3)}(0) \int \pi^{3}(u) d u}{3}\left(\sqrt{\int \pi^{2}(u) d u \int K^{2}(v) d v}\right)^{-3} \sqrt{h^{d}} . \tag{A.3}
\end{align*}
$$

We need the following lemmas; their proofs are given in the supplemental material ${ }^{7}$.
Lemma A.1. Suppose that the conditions of Theorem 2.2(i) hold. Then for any $h$

$$
\begin{equation*}
\sup _{x \in \mathbb{R}^{1}}\left|P\left(L_{n}(h) \leq x\right)-\Phi(x)+\rho(h)\left(x^{2}-1\right) \phi(x)\right|=O\left(h^{d}\right) . \tag{A.4}
\end{equation*}
$$

Recall $L_{n}(h)=\sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} e_{i} a_{i j} e_{j}$ as defined in (A.3) and let

$$
\bar{T}_{n}(h)=\frac{h^{\frac{d}{2}}}{n \sigma_{n}} \sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} \widehat{e}_{i} K_{h}\left(X_{i}-X_{j}\right) \widehat{e}_{j}=\frac{h^{\frac{d}{2}}}{n \sigma_{n}} \sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} e_{i} K_{h}\left(X_{i}-X_{j}\right) e_{j}
$$

[^4]\[

$$
\begin{align*}
& +\frac{h^{\frac{d}{2}}}{n \sigma_{n}} \sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} K_{h}\left(X_{i}-X_{j}\right)\left[m\left(X_{i}\right)-m_{\widehat{\theta}}\left(X_{i}\right)\right]\left[m\left(X_{j}\right)-m_{\widehat{\theta}}\left(X_{j}\right)\right] \\
& +\frac{2 h^{\frac{d}{2}}}{n \sigma_{n}} \sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} e_{i} K_{h}\left(X_{i}-X_{j}\right)\left[m\left(X_{j}\right)-m_{\widehat{\theta}}\left(X_{j}\right)\right] \\
\equiv & L_{n}(h)+S_{n}(h)+D_{n}(h), \tag{A.5}
\end{align*}
$$
\]

where $S_{n}(h)=\frac{h^{\frac{d}{2}}}{n \sigma_{n}} \sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} K_{h}\left(X_{i}-X_{j}\right)\left[m\left(X_{i}\right)-m_{\widehat{\theta}}\left(X_{i}\right)\right]\left[m\left(X_{j}\right)-m_{\widehat{\theta}}\left(X_{j}\right)\right]$ and

$$
\begin{equation*}
D_{n}(h)=\frac{2 h^{\frac{d}{2}}}{n \sigma_{n}} \sum_{i=1}^{n} \sum_{j=1, \neq i}^{n} e_{i} K_{h}\left(X_{i}-X_{j}\right)\left[m\left(X_{j}\right)-m_{\widehat{\theta}}\left(X_{j}\right)\right] . \tag{A.6}
\end{equation*}
$$

Define $L_{n}^{*}(h), S_{n}^{*}(h)$ and $D_{n}^{*}(h)$ as the corresponding versions of $L_{n}(h), S_{n}(h)$ and $D_{n}(h)$ involved in (A.5) with $\left(X_{i}, Y_{i}\right)$ and $\widehat{\theta}$ being replaced by $\left(X_{i}, Y_{i}^{*}\right)$ and $\widehat{\theta^{*}}$ respectively.

Lemma A.2. Suppose that the conditions of Theorem 2.2(i) hold. Then

$$
\begin{equation*}
\sup _{x \in \mathbb{R}^{1}}\left|P^{*}\left(L_{n}^{*}(h) \leq x\right)-\Phi(x)+\rho(h)\left(x^{2}-1\right) \phi(x)\right|=O_{P}\left(h^{d}\right) . \tag{A.7}
\end{equation*}
$$

Lemma A.3. Suppose that the conditions of Theorem 2.2(i) hold. Then under $\mathcal{H}_{0}$

$$
\begin{align*}
E\left[S_{n}(h)\right] & =O\left(\sqrt{h^{d}}\right) \quad \text { and } \quad E\left[D_{n}(h)\right]=o\left(\sqrt{h^{d}}\right),  \tag{A.8}\\
E^{*}\left[S_{n}^{*}(h)\right] & =O_{P}\left(\sqrt{h^{d}}\right) \quad \text { and } \quad E^{*}\left[D_{n}^{*}(h)\right]=o_{P}\left(\sqrt{h^{d}}\right),  \tag{A.9}\\
E\left[S_{n}(h)\right]-E^{*}\left[S_{n}^{*}(h)\right] & =O_{P}\left(\sqrt{h^{d}}\right) \text { and } E\left[D_{n}(h)\right]-E^{*}\left[D_{n}^{*}(h)\right]=o_{P}\left(\sqrt{h^{d}}\right) \tag{A.10}
\end{align*}
$$

in probability with respect to the joint distribution of $\mathcal{W}_{n}$, where $\left.E^{*} \cdot \cdot\right]=E\left[\cdot \mid \mathcal{W}_{n}\right]$.
Lemma A.4. Suppose that the conditions of Theorem 2.2(ii) hold. Then under $\mathcal{H}_{1}$

$$
\begin{equation*}
\lim _{n \rightarrow \infty} E\left[S_{n}(h)\right]=\infty \quad \text { and } \quad \lim _{n \rightarrow \infty} \frac{E\left[D_{n}(h)\right]}{E\left[S_{n}(h)\right]}=0 . \tag{A.11}
\end{equation*}
$$

## A.3. Proof of Theorem 2.1:

A.3.1. Proof of (2.10): Recall from (2.8) and (A.5)-(A.6) that

$$
\begin{align*}
& \widehat{T}_{n}(h)=\left(L_{n}(h)+S_{n}(h)+D_{n}(h)\right) \cdot \frac{\sigma_{n}}{\widehat{\sigma}_{n}}+o_{P}\left(\sqrt{h^{d}}\right),  \tag{A.12}\\
& \widehat{T}_{n}^{*}(h)=\left(L_{n}^{*}(h)+S_{n}^{*}(h)+D_{n}^{*}(h)\right) \cdot \frac{\sigma_{n}}{\widehat{\sigma}_{n}^{*}}+o_{P}\left(\sqrt{h^{d}}\right), \tag{A.13}
\end{align*}
$$

where $\sigma_{n}^{2}, \widehat{\sigma}_{n}^{2}$ and $\widehat{\sigma}_{n}^{* 2}$ are as defined in (2.4), (2.7) and (2.9) respectively.
In view of Assumption A. 2 and Lemmas A.1-A.3, we may ignore any terms with orders higher than $\sqrt{h^{d}}$ and then consider the following approximations:

$$
\begin{align*}
& \widehat{T}_{n}(h)=L_{n}(h)+E\left[S_{n}(h)\right]+o_{P}\left(\sqrt{h^{d}}\right) \text { and } \\
& \widehat{T}_{n}^{*}(h)=L_{n}^{*}(h)+E^{*}\left[S_{n}^{*}(h)\right]+o_{P}\left(\sqrt{h^{d}}\right) . \tag{A.14}
\end{align*}
$$

Let $s(h)=E\left[S_{n}(h)\right]$ and $s^{*}(h)=E^{*}\left[S_{n}^{*}(h)\right]$. We then apply Lemmas A. 1 and A. 2 to obtain that uniformly over $x \in \mathbb{R}^{1}$,

$$
\begin{align*}
P\left(\widehat{T}_{n}(h) \leq x\right) & =P\left(L_{n}(h) \leq x-s(h)+o_{P}\left(\sqrt{h^{d}}\right)\right)  \tag{A.15}\\
& =\Phi(x-s(h))-\rho(h)\left((x-s(h))^{2}-1\right) \phi(x-s(h))+o\left(\sqrt{h^{d}}\right) \text { and } \\
P^{*}\left(\widehat{T}_{n}^{*}(h) \leq x\right) & =P^{*}\left(L_{n}^{*}(h) \leq x-s^{*}(h)+o_{P}\left(\sqrt{h^{d}}\right)\right) \\
& =\Phi\left(x-s^{*}(h)\right)-\rho(h)\left(\left(x-s^{*}(h)\right)^{2}-1\right) \phi\left(x-s^{*}(h)\right)+o_{P}\left(\sqrt{h^{d}}\right) .
\end{align*}
$$

Theorem 2.2(i) follows consequently from (A.10) and (A.15).
A.3.2. Proof of (2.11): In view of the definition that $P^{*}\left(\widehat{T}_{n}^{*}(h) \geq l_{\alpha}^{*}\right)=\alpha$ and the conclusion from Theorem 2.1(i) that

$$
\begin{equation*}
P\left(\widehat{T}_{n}(h) \geq l_{\alpha}^{*}\right)-P^{*}\left(\widehat{T}_{n}^{*}(h) \geq l_{\alpha}^{*}\right)=O_{P}\left(\sqrt{h^{d}}\right), \tag{A.16}
\end{equation*}
$$

the proof of $P\left(\widehat{T}_{n}(h) \geq l_{\alpha}^{*}\right)=\alpha+O\left(\sqrt{h^{d}}\right)$ follows unconditionally from the dominated convergence theorem.

## A.4. Proof of Theorem 2.2

It follows from Lemmas A.1-A. 4 that

$$
\begin{align*}
\alpha_{n}(h) & =P\left(\widehat{T}_{n}(h) \geq l_{\alpha} \mid \mathcal{H}_{0}\right)=P\left(L_{n}(h) \geq l_{\alpha}-S_{n}(h)+o_{P}\left(S_{n}(h)\right) \mid \mathcal{H}_{0}\right) \\
& =1-P\left(L_{n}(h) \leq l_{\alpha}-S_{n}(h)+o_{P}\left(S_{n}(h)\right) \mid \mathcal{H}_{0}\right),  \tag{A.17}\\
\alpha_{n}^{*}(h) & =P\left(\widehat{T}_{n}(h) \geq l_{\alpha}^{*} \mid \mathcal{H}_{0}\right)=P\left(L_{n}(h) \geq l_{\alpha}^{*}-S_{n}(h)+o_{P}\left(S_{n}(h)\right) \mid \mathcal{H}_{0}\right) \\
& =1-P\left(L_{n}(h) \leq l_{\alpha}^{*}-S_{n}(h)+o_{P}\left(S_{n}(h)\right) \mid \mathcal{H}_{0}\right),  \tag{A.18}\\
\beta_{n}(h) & =P\left(\widehat{T}_{n}(h) \geq l_{\alpha} \mid \mathcal{H}_{1}\right)=P\left(L_{n}(h) \geq l_{\alpha}-S_{n}(h)+o_{P}\left(S_{n}(h)\right) \mid \mathcal{H}_{1}\right) \\
& =1-P\left(L_{n}(h) \leq l_{\alpha}-S_{n}(h)+o_{P}\left(S_{n}(h)\right) \mid \mathcal{H}_{1}\right), \tag{A.19}
\end{align*}
$$

$$
\begin{align*}
\beta_{n}^{*}(h) & =P\left(\widehat{T}_{n}(h) \geq l_{\alpha}^{*} \mid \mathcal{H}_{1}\right)=P\left(L_{n}(h) \geq l_{\alpha}^{*}-S_{n}(h)+o_{P}\left(S_{n}(h)\right) \mid \mathcal{H}_{1}\right) \\
& =1-P\left(L_{n}(h) \leq l_{\alpha}^{*}-S_{n}(h)+o_{P}\left(S_{n}(h)\right) \mid \mathcal{H}_{1}\right) \tag{A.20}
\end{align*}
$$

Using Assumptions A.2(iv) and A.3, a Taylor expansion of $m_{\theta}(\cdot)$ at $\theta_{0}$ implies that for sufficiently large $n$

$$
\begin{align*}
& S_{n}(h)=C_{0}(m) \sqrt{h^{d}}\left(1+o_{P}(1)\right) \text { under } \mathcal{H}_{0} \quad \text { and }  \tag{A.21}\\
& S_{n}(h)=n C_{n}^{2} \sqrt{h^{d}}\left(1+o_{P}(1)\right) \text { under } \mathcal{H}_{1} \tag{A.22}
\end{align*}
$$

hold in probability, where $C_{n}^{2}$ is as defined in Theorem 2.2(ii). The proof of Theorem 2.2 then follows from (A.15) and (A.17)-(A.22).
A.5. Proof of Theorem 2.3: The proof follows from that of Theorem 2.2. The details are given in the supplemental material ${ }^{8}$.

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[^1]:    ${ }^{2}$ available at http://www.amstat.org/publications/jasa/supplemental-materials

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[^5]:    ${ }^{8}$ available at http://www.amstat.org/publications/jasa/supplemental-materials

