Kernel-based Conditional Independence Test and Application in Causal Discovery

Kun Zhang  Jonas Peters  Dominik Janzing  Bernhard Schölkopf
Max Planck Institute for Intelligent Systems
Spemannstr. 38, 72076 Tübingen
Germany

Abstract

Conditional independence testing is an important problem, especially in Bayesian network learning and causal discovery. Due to the curse of dimensionality, testing for conditional independence of continuous variables is particularly challenging. We propose a Kernel-based Conditional Independence test (KCI-test), by constructing an appropriate test statistic and deriving its asymptotic distribution under the null hypothesis of conditional independence. The proposed method is computationally efficient and easy to implement. Experimental results show that it outperforms other methods, especially when the conditioning set is large or the sample size is not very large, in which case other methods encounter difficulties.

1 Introduction

Statistical independence and conditional independence (CI) are important concepts in statistics, artificial intelligence, and related fields (Dawid 1979). Let $X$, $Y$, and $Z$ denote sets of random variables. The CI between $X$ and $Y$ given $Z$, denoted by $X \perp \perp Y | Z$, reflects the fact that given the values of $Z$, further knowing the values of $X$ (or $Y$) does not provide any additional information about $Y$ (or $X$). Independence and CI play a central role in causal discovery and Bayesian network learning (Pearl 2000; Spirtes et al. 2001; Koller and Friedman 2009). Generally speaking, the CI relationship $X \perp Y | Z$ allows us to drop $Y$ when constructing a probabilistic model for $X$ with $(Y, Z)$, which results in a parsimonious representation.

Testing for CI is much more difficult than that for unconditional independence (Bergsma 2004). For CI tests, traditional methods either focus on the discrete case, or impose simplifying assumptions to deal with the continuous case – in particular, the variables are often assumed to have linear relations with additive Gaussian errors. In that case, $X \perp Y | Z$ reduces to zero partial correlation or zero conditional correlation between $X$ and $Y$ given $Z$, which can be easily tested (for the links between partial correlation, conditional correlation, and CI, see Lawrance (1976)). However, nonlinearity and non-Gaussian noise are frequently encountered in practice, and hence this assumption can lead to incorrect conclusions.

Recently, practical methods have been proposed for testing CI for continuous variables without assuming a functional form between the variables as well as the data distributions, which is the case we are concerned with in this paper. To our knowledge, the existing methods fall into four categories. The first category is based on explicit estimation of the conditional densities or their variants. For example, Su and White (2008) define the test statistic as some distance between the estimated conditional densities $p(X|Y, Z)$ and $p(X|Y)$, and Su and White (2007) exploit the difference between the characteristic functions of these conditional densities. The estimation of the conditional densities or related quantities is difficult, which deteriorates the testing performance especially when the conditioning set $Z$ is not small enough. Methods in the second category, such as Margaritis (2005) and Huang (2010), discretize the conditioning set $Z$ to a set of bins, and transform CI to the unconditional one in each bin. Inevitably, due to the curse of dimensionality, as the conditioning set becomes larger, the required sample size increases dramatically. Methods in the third category, including Linton and Gozalo (1997) and Song (2009), provide slightly weaker tests than that for CI. For instance, the method proposed by Song (2009) tests whether one can find some (nonlinear) function $h$ and parameters $\theta_0$, such that $X$ and $Y$ are conditionally independent given a single index function $\lambda_{\theta_0}(Z) = h(Z^T \theta_0)$ of $Z$. In general, this is different from the test for $X \perp Y | Z$: to see this, consider the case where $X$ and $Y$ depend on two different
but overlapping subsets of $Z$; even if $X \perp Y | Z$, it is impossible to find $\lambda_{0i}(Z)$ given which $X$ and $Y$ are conditionally independent. 

[Fukumizu et al. (2004)] give a general nonparametric characterization of CI using covariance operators in the reproducing kernel Hilbert spaces (RKHS), which inspired a kernel-based measure of conditional dependence (see also [Fukumizu et al. (2008)]). However, the distribution of this measure under the CI hypothesis is unknown, and consequently, it could not directly serve as a CI test. To get a test, one has to combine this conditional dependence measure with local bootstrap, or local permutation, which is used to determine the rejector region [Fukumizu et al. 2008; Tillman et al. 2009]. This leads to the method in the fourth category. We denote it by $\text{CI}_{\text{PERM}}$. Like the methods in the second category, this approach would require a large sample size and tends to be unreliable when the number of conditioning variables increases.

In this paper we aim to develop a CI testing method which avoids the above drawbacks. In particular, based on appropriate characterizations of CI, we define a simple test statistic which can be easily calculated from the kernel matrices associated with $X$, $Y$, and $Z$, and we further derive its asymptotic distribution under the null hypothesis. We also provide ways to estimate such a distribution, and finally CI can be tested conveniently. This results in a Kernel-based Conditional Independence test (KCI-test). In this procedure we do not explicitly estimate the conditional or joint densities, nor discretize the conditioning variables. Our method is computationally appealing and is less sensitive to the dimensionality of $Z$ compared to other methods. Our results contain unconditional independence testing (similar to [Gretton et al. 2008]) as a special case.

2 Characterization of Independence and Conditional Independence

We introduce the following notational convention. Throughout this paper, $X$, $Y$, and $Z$ are continuous random variables or sets of continuous random variables, with domains $\mathcal{X}$, $\mathcal{Y}$, and $\mathcal{Z}$, respectively. Define a measurable, positive definite kernel $k_X$ on $\mathcal{X}$ and denote the corresponding RKHS by $\mathcal{H}_X$. Similarly we define $k_Y$, $\mathcal{H}_Y$, $k_Z$, and $\mathcal{H}_Z$. In this paper we assume that all involved RKHS’s are separable and square integrable. The probability law of $X$ is denoted by $P_X$, and similarly for the joint probability laws such as $P_{XZ}$. The spaces of square integrable functions of $X$ and $(X, Z)$ are denoted by $L^2_X$ and $L^2_{XZ}$, respectively. E.g., $L^2_{XZ} = \{g(X, Z) | \mathbb{E}[g^2] < \infty\}$. $x = \{x_1, ..., x_n\}$ denotes the i.i.d. sample of $X$ of size $n$. $K_X$ is the kernel matrix of the sample $x$, and the corresponding centralized kernel matrix is $\tilde{K}_X = \mathbb{H}K_X\mathbb{H}$, where $\mathbb{H} = I - \frac{1}{n}11^T$ with $I$ and $1$ being the $n \times n$ identity matrix and the vector of 1’s, respectively. By default we use the Gaussian kernel, i.e., the $(i, j)$th entry of $K_X$ is $k(x_i, x_j) = \exp(-\frac{||x_i-x_j||^2}{2\sigma_X^2})$, where $\sigma_X$ is the kernel width. Similar notations are used for $Y$ and $Z$.

The problem we consider here is to test for CI between sets of continuous variables $X$ and $Y$ given $Z$ from their observed i.i.d. samples, without making specific assumptions on their distributions or the functional forms between them. $X$ and $Y$ are said to be conditionally independent given $Z$ if and only if $p_{X\mid Y, Z} = p_X|Z$ (or equivalently, $p_{Y\mid X, Z} = p_Y|Z$, or $p_{XY\mid Z} = p_{X}|p_{Y}|Z$). Therefore, a direct way to assess if $X \perp Y | Z$ is to estimate certain probability densities and then evaluate if the above equation is plausible. However, density estimation in high dimensions is a difficult problem: it is well known that in nonparametric joint or conditional density estimation, due to the curse of dimensionality, to achieve the same accuracy the number of required data points is exponentially increasing in the data dimension. Fortunately, conditional (in)dependence is just one particular property associated with the distributions; to test for it, it is possible to avoid explicitly estimating the densities.

There are other ways to characterize the CI relation that do not explicitly involve the densities or the variates, and they may result in more efficient methods for CI testing. Recently, a characterization of CI is given in terms of the cross-covariance operator $\Sigma_{XY}$ on RKHS ([Fukumizu et al. 2004]). For the random vector $(X, Y)$ on $\mathcal{X} \times \mathcal{Y}$, the cross-covariance operator from $\mathcal{H}_X$ to $\mathcal{H}_Y$ is defined by the relation:

\[ (g, \Sigma_{XY} f) = \mathbb{E}_X[f(X)g(Y)] - \mathbb{E}_X[f(X)]\mathbb{E}_Y[g(Y)] \]

for all $f \in \mathcal{H}_X$ and $g \in \mathcal{H}_Y$. The conditional cross-covariance operator of $(X, Y)$ given $Z$ is further defined by:

\[ \Sigma_{X\mid Y, Z} = \Sigma_{XY} - \Sigma_{YZ}\Sigma_{ZZ}^{-1}\Sigma_{ZX}. \]

Intuitively, one can interpret it as the partial covariance between $\{f(X), \forall f \in \mathcal{H}_X\}$ and $\{g(Y), \forall g \in \mathcal{H}_Y\}$ given $\{h(Z), \forall h \in \mathcal{H}_Z\}$.

If characteristic kernels are used, the conditional...
cross-covariance operator is related to the CI relation, as seen from the following lemma.

**Lemma 1** [Characterization based on conditional cross-covariance operators (Fukumizu et al. 2008)]

Denote $\tilde{X} \triangleq (X, Z)$, $k_{\tilde{X}} \triangleq k_X k_Z$, and $\mathcal{H}_{\tilde{X}}$ the RKHS corresponding to $k_{\tilde{X}}$. Assume $\mathcal{H}_X \subset L^2_X$, $\mathcal{H}_Y \subset L^2_Y$, and $\mathcal{H}_Z \subset L^2_Z$. Further assume that $k_X k_Y$ is a characteristic kernel on $(X \times Y) \times Z$, and that $\mathcal{H}_Z + \mathbb{R}$ (the direct sum of the two RKHSs) is dense in $L^2(P_Z)$. Then

$$\Sigma_{XY|Z} = 0 \iff X \perp Y|Z. \quad (2)$$

Note that one can replace $\Sigma_{XY|Z}$ with $\Sigma_{\tilde{X}Y|Z}$, where $\tilde{Y} \triangleq (Y, Z)$, in the above lemma. Alternatively, Daudin (1980) gives the characterization of CI by explicitly enforcing the uncorrelatedness of functions in suitable spaces, which may be intuitively more appealing. In particular, consider the constrained $L^2$ spaces

$$\mathcal{E}_{XZ} \triangleq \{ f \in L^2_{XZ} \mid \mathbb{E}(f|Z) = 0 \},$$
$$\mathcal{E}_X \triangleq \{ g \in L^2_X \mid \mathbb{E}(g|Z) = 0 \},$$
$$\mathcal{E}_Y \triangleq \{ g' \mid g' = g'(Y) - \mathbb{E}(g'|Z), \ g' \in L^2_Y \}. \quad (3)$$

They can be constructed from the corresponding $L^2$ spaces via nonlinear regression. For instance, for any function $f \in L^2_{XZ}$, the corresponding function $\tilde{f}$ is given by

$$\tilde{f}(\tilde{X}) = f(\tilde{X}) - \mathbb{E}(f|Z) = f(\tilde{X}) - h^*_f(Z), \quad (4)$$

where $h^*_f(Z) \in L^2_Z$ is the regression function of $f(\tilde{X})$ on $Z$. One can then relate CI to uncorrelatedness in the following way.

**Lemma 2** [Characterization based on partial association (Daudin 1980)]

The following conditions are equivalent to each other. (i.) $X \perp Y|Z$; (ii.) $\mathbb{E}(\tilde{g}) = 0$, $\forall \tilde{f} \in \mathcal{E}_{XZ}$ and $\tilde{g} \in \mathcal{E}_{YZ}$; (iii.) $\mathbb{E}((\tilde{f} g) = 0$, $\forall \tilde{f} \in \mathcal{E}_{XZ}$ and $g \in L^2_Y$; (iv.) $\mathbb{E}(\tilde{f}g') = 0$, $\forall \tilde{f} \in \mathcal{E}_{XZ}$ and $g' \in \mathcal{E}'_{Y}$; (v.) $\mathbb{E}(\tilde{f}g') = 0$, $\forall \tilde{f} \in \mathcal{E}_{XZ}$ and $g' \in \mathcal{E}'_Y$.

The above result can be considered as a generalization of the partial correlation based characterization of CI for Gaussian variables. Suppose that $(X, Y, Z)$ is jointly Gaussian; then $X \perp Y|Z$ is equivalent to the vanishing of the partial correlation coefficient $\rho_{XY|Z}$. Here, intuitively speaking, condition (ii) means that any “residual” function of $(X, Z)$ given $Z$ is uncorrelated from that of $(Y, Z)$ given $Z$. Note that $\mathcal{E}_{XZ}$ (resp. $\mathcal{E}_{YZ}$) contains all functions of $X$ and $Z$ (resp. of $Y$ and $Z$) that cannot be “explained” by $Z$, in the sense that any function $\tilde{f} \in \mathcal{E}_{XZ}$ (resp. $\tilde{g} \in \mathcal{E}_{YZ}$) is uncorrelated with any function of $Z$ (Daudin, 1980).

From the definition of the conditional cross-covariance operator $[1]$, one can see the close relationship between the conditions in Lemma 1 and those in Lemma 2. However, Lemma 1 has practical advantages: in Lemma 2 one has to consider all functions in $L^2$, while Lemma 1 exploits the spaces corresponding to some characteristic kernels, which might be much smaller. In fact, if we restrict the functions $f$ and $g'$ to the spaces $\mathcal{H}_X$ and $\mathcal{H}_Y$, respectively, Lemma 2 is then reduced to Lemma 1. The above characterizations of CI motivated our statistic for testing $X \perp Y|Z$, as presented below.

### 3 A Kernel-Based Conditional Independence Test

#### 3.1 General results

As seen above, independence and CI can be characterized by uncorrelatedness between functions in certain spaces. We first give some general results on the asymptotic distributions of some statistics defined in terms of the kernel matrices under the condition of such uncorrelatedness. Later those results will be used for testing for CI as well as unconditional independence.

Suppose that we are given the i.i.d. samples $x \triangleq (x_1, ..., x_t, ..., x_n)$ and $y \triangleq (y_1, ..., y_t, ..., y_n)$ for $X$ and $Y$, respectively. Suppose further that we have the eigenvalue decompositions (EVD) of the centralized kernel matrices $\overline{K}_X$ and $\overline{K}_Y$, i.e., $\overline{K}_X = V_X \Lambda_X V_X^T$ and $\overline{K}_Y = V_Y \Lambda_Y V_Y^T$, where $\Lambda_X$ and $\Lambda_Y$ are the diagonal matrices containing the non-negative eigenvalues $\lambda_{X,i}$ and $\lambda_{Y,i}$, respectively. Here, the eigenvalues are sorted in descending order, i.e., $\lambda_{X,1} \geq \lambda_{X,2} \geq \ldots \geq \lambda_{X,i} \geq 0$ and $\lambda_{Y,1} \geq \lambda_{Y,2} \geq \ldots \geq \lambda_{Y,i} \geq 0$.

Let $\psi = [\psi_1(x), ..., \psi_n(x)] \triangleq V_X \Lambda_X^{1/2}$ and $\phi_y = [\phi_1(y), ..., \phi_n(y)] \triangleq V_Y \Lambda_y^{1/2}$. I.e., $\psi_i(x) = \sqrt{\lambda_{X,i}} \psi_{X,i}$, where $\psi_{X,i}$ denotes the $i$th eigenvector of $\overline{K}_X$.

On the other hand, consider the eigenvalues $\lambda_{X,i}$ and eigenfunctions $u_{X,i}$ of the kernel $k_X$ w.r.t. the probability measure with the density $p(x)$, i.e., $\lambda_{X,i}$ and $u_{X,i}$ satisfy

$$\int k_X(x, x') u_{X,i}(x)p(x)dx = \lambda_{X,i} u_{X,i}(x'). \quad (6)$$

We assume that $u_{X,i}$ have unit variance, i.e., $\mathbb{E}[u_{X,i}(x)] = 1$. We also sort $\lambda_{X,i}$ in descending order. Similarly, we define $\lambda_{Y,i}$ and $u_{Y,i}$ of $k_Y$. Define

$$S_{ij} \triangleq \frac{1}{\sqrt{n}} \psi_i(x)^T \phi_j(y) = \frac{\sum_{t=1}^n \psi_i(x_t) \phi_j(y_t)}{\sqrt{n}}$$

with $\psi_i(x_t)$ being the $t$-th component of the vector $\psi_i(x)$. We then have the following results.
Theorem 3 Suppose that we are given arbitrary centered kernels \( k_X \) and \( k_Y \) with discrete eigenvalues and the corresponding RKHS’s \( H_X \) and \( H_Y \) for sets of random variables \( X \) and \( Y \). We have the following three statements.

1) Under the condition that \( f(X) \) and \( g(Y) \) are uncorrelated for all \( f \in H_X \) and \( g \in H_Y \), for any \( L \) such that \( \lambda_{X,L+1}^* \neq \lambda_{X,L}^* \) and \( \lambda_{Y,L+1}^* \neq \lambda_{Y,L}^* \), we have

\[
\sum_{i,j=1}^{L} S_{ij}^2 \overset{d}{\to} \sum_{k=1}^{L} \lambda_{k}^* z_{k}^2, \quad \text{as } n \to \infty,
\]

where \( z_k \) are i.i.d. standard Gaussian variables (i.e., \( z_k^2 \) are i.i.d. \( \chi_1^2 \)-distributed variables), \( \lambda_k^* \) are the eigenvalues of \( E(ww^T) \) and \( w \) is the random vector obtained by stacking the \( L \times L \) matrix \( N \) whose \((i,j)\)th entry is \( N_{ij} = \sqrt{\lambda_{X,i}^* \lambda_{Y,j}^*} u_{X,i}(X) u_{Y,j}(Y) \).

2) In particular, if \( X \) and \( Y \) are further independent, we have

\[
\sum_{i,j=1}^{L} S_{ij}^2 \overset{d}{\to} \sum_{i,j=1}^{L} \lambda_{X,i}^* \lambda_{Y,j}^* \cdot z_{ij}^2, \quad \text{as } n \to \infty,
\]

where \( z_{ij}^2 \) are i.i.d. \( \chi_1^2 \)-distributed variables.

3) The results (3) and (6) also hold if \( L = n \to \infty \).

All proofs are sketched in Appendix. We note that

\[
\text{Tr}(\tilde{K}_X \tilde{K}_Y) = \text{Tr}(\psi_x^T \psi_x^T \psi_y^T \psi_y^T)
= \text{Tr}(\psi_x^T \psi_y^T \psi_y^T \psi_x^T) = n \sum_{i,j=1}^{n} S_{ij}^2,
\]

Hence, the above theorem gives the asymptotic distribution of \( T \) under the condition that \( f(X) \) and \( g(Y) \) are always uncorrelated. Combined with the characterizations of (conditional) independence, this inspires the corresponding testing method. We further give the following remarks on Theorem 3.

First, \( X \) and \( Y \) are not necessarily disjoint, and \( H_X \) and \( H_Y \) can be any RKHS’s. Second, in practice the eigenvalues \( \lambda_k^* \) are not known, and one needs to use the empirical ones instead, as discussed below.

3.2 Unconditional independence testing

As a direct consequence of the above theorem, we have the following result which allows for kernel-based unconditional independence testing.

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\( ^3 \)Equivalently, one can consider \( \lambda_k^* \) as the eigenvalues of the tensor \( T \) with \( T_{ijkl} = E(N_{ij,t} N_{kl,t}) \).

Theorem 4 [Independence test]

Under the null hypothesis that \( X \) and \( Y \) are statistically independent, the statistic

\[
T_{UI} \triangleq \frac{1}{n} \text{Tr}(\tilde{K}_X \tilde{K}_Y)
\]

has the same asymptotic distribution as

\[
\tilde{T}_{UI} \triangleq \frac{1}{n^2} \sum_{i,j=1}^{n} \lambda_{x,i}^* \lambda_{y,j}^* z_{ij}^2,
\]

i.e., \( T_{UI} \overset{d}{\to} \tilde{T}_{UI} \) as \( n \to \infty \).

This theorem inspires the following unconditional independence testing procedure. Given the samples \( x \) and \( y \), one first calculates the centralized kernel matrices \( \tilde{K}_X \) and \( \tilde{K}_Y \) and their eigenvalues \( \lambda_{x,i}^* \) and \( \lambda_{y,j}^* \), and then evaluates the statistic \( \tilde{T}_{UI} \) according to (8).

Next, the empirical null distribution of \( T_{UI} \) under the null hypothesis can be simulated in the following way: one draws i.i.d. random samples from the \( \chi_1^2 \)-distributed variables \( z_{ij}^2 \), and then generates samples for \( \tilde{T}_{UI} \) according to (9). (Later we will give another way to approximate the asymptotic null distribution.) Finally the \( p \)-value can be found by locating \( T_{UI} \) in the null distribution.

This unconditional independence testing method is closely related to the one based on the Hilbert-Schmidt independence criterion (HSIC) proposed by [Gretton et al. (2008)]. Actually the defined statistics (our statistic \( T_{UI} \) and HSIC in [Gretton et al. (2008)]) are the same, but the asymptotic distributions are in different forms. In our results the asymptotic distribution only involves the eigenvalues of the two regular kernel matrices \( \tilde{K}_X \) and \( \tilde{K}_Y \), while in their results it depends on the eigenvalues of an order-four tensor, which are more difficult to calculate.

3.3 Conditional independence testing

Here we would like to make use of condition (ii) in Lemma 2 to test for CI, but the considered functional spaces are \( f(\tilde{X}) \in H_{\tilde{X}}, g'(Y) \in H_Y, h^*_y(Z) \in H_Z \), and \( h^*_x(Z) \in H_Z \), as in Lemma 1. The functions \( f(\tilde{X}) \) and \( g'(Y,Z) \) appearing in condition (ii), whose spaces are denoted by \( H_{\tilde{X}|Z} \) and \( H_{Y|Z} \), respectively, can be constructed from the functions \( f, g', h^*_y, \) and \( h^*_y \).

Suppose that we already have the centralized kernel matrices \( \tilde{K}_{\tilde{X}} \) (which is the centralized kernel matrix of \( \tilde{X} = (X,Z) \)), \( \tilde{K}_Y \), and \( \tilde{K}_Z \) on the samples \( x, y, \) and \( z \). We use kernel ridge regression to estimate the regression function \( h^*_y(Z) \) in (1), and one can easily see that \( h^*_y(z) = \tilde{K}_Z(\tilde{K}_Z + \varepsilon I)^{-1} \cdot f(\tilde{x}) \), where \( \varepsilon \) is a
small positive regularization parameter (Schölkopf and Smola 2002). Consequently, \( f(\tilde{x}) \) can be constructed as \( f(\tilde{x}) = f(\tilde{x}) - \hat{h}_f^*(z) = R_Z \cdot f(\tilde{x}) \), where

\[
R_Z = I - \tilde{K}_Z(K_Z + \epsilon I)^{-1} = \epsilon(\tilde{K}_Z + \epsilon I)^{-1}.
\]  

(10)

Based on the EVD decomposition \( \tilde{K}_X = V_{X}^2 \Lambda_{X} V_{X} \), we can construct \( \varphi_{X} = [\varphi_{1}(\tilde{x}), \ldots, \varphi_{n}(\tilde{x})] \cong V_{X} \Lambda_{X}^{1/2} \) as an empirical kernel map for \( \tilde{x} \). Correspondingly, an empirical kernel map for the space \( H_{X|Z} \) is given by \( \tilde{\varphi}_{X} = R_Z \varphi_{X}(\tilde{x}) \). Consequently, the centralized kernel matrix corresponding to the functions \( \tilde{f}(\tilde{x}) \) is

\[
\tilde{K}_{X|Z} = \tilde{\varphi}_{X} \tilde{\varphi}_{X}^T = R_Z \tilde{K}_X R_Z.
\]  

(11)

Similarly, that corresponding to \( \tilde{g} \) is

\[
\tilde{K}_{Y|Z} = R_Z \tilde{K}_Y R_Z.
\]  

(12)

Furthermore, we let the EVD decompositions of \( \tilde{K}_{X|Z} \) and \( \tilde{K}_{Y|Z} \) be \( \tilde{K}_{X|Z} = E_{X|Z} \Lambda_{X|Z} E_{X|Z}^T \) and \( \tilde{K}_{Y|Z} = E_{Y|Z} \Lambda_{Y|Z} E_{Y|Z}^T \), respectively. \( \Lambda_{X|Z} \) (resp. \( \Lambda_{Y|Z} \)) is the diagonal matrix containing non-negative eigenvalues \( \lambda_{X|Z,i} \) (resp. \( \lambda_{Y|Z,i} \)). Let \( \psi_{X|Z} = [\psi_{X|Z,1}(\tilde{x}), \ldots, \psi_{X|Z,n}(\tilde{x})] \cong V_{X|Z} \Lambda_{X|Z}^{1/2} \) and \( \phi_{Y|Z} = [\phi_{Y|Z,1}(\tilde{y}), \ldots, \phi_{Y|Z,n}(\tilde{y})] \cong V_{Y|Z} \Lambda_{Y|Z}^{1/2} \). We then have the following result which the proposed KCI-test is based on.

**Proposition 5 [Conditional independence test]**

Under the null hypothesis \( H_0 \) (\( X \) and \( Y \) are conditionally independent given \( Z \)), we have that the statistic

\[
T_{CI} \triangleq \frac{1}{n} \text{Tr}(\tilde{K}_{X|Z} \tilde{K}_{Y|Z})
\]  

has the same asymptotic distribution as

\[
T_{CI} \triangleq \frac{1}{n} n \sum_{k=1}^{n} \hat{\lambda}_k \cdot z_k^2,
\]  

(13)

where \( \hat{\lambda}_k \) are eigenvalues of \( \tilde{w} \tilde{w}^T \) and \( \tilde{w} = [\tilde{w}_1, \ldots, \tilde{w}_n] \) with the vector \( \tilde{w}_t \) obtained by stacking \( M_t = [\tilde{\psi}_{X|Z,1}(\tilde{x}_t), \ldots, \tilde{\psi}_{X|Z,n}(\tilde{x}_t)]^T \) and \( [\tilde{\phi}_{Y|Z,1}(\tilde{y}_t), \ldots, \tilde{\phi}_{Y|Z,n}(\tilde{y}_t)] \). \( z_k \) are i.i.d. \( \chi^2 \) samples and summing them up with weights \( \hat{\lambda}_k \). (For computational efficiency, in practice we drop all \( \lambda_{X|Z,i} \), \( \lambda_{Y|Z,i} \), and \( \hat{\lambda}_k \) which are smaller than \( 10^{-5} \).)

Finally the \( p \)-value is calculated as the probability of \( T_{CI} \) exceeding \( T_{CI} \). Approximating the null distribution with a Gamma distribution, which is given next, avoids calculating \( \lambda_k \) and simulating the null distribution, and is computationally more efficient.

### 3.4 Approximating the null distribution by a Gamma distribution

In addition to the simulation-based method to find the null distribution, as in Gretton et al. (2008), we provide approximations to the null distributions with a two-parameter Gamma distribution. The two parameters of the Gamma distribution are related to the mean and variance. In particular, under the null hypothesis that \( X \) and \( Y \) are independent (resp. conditionally independent given \( Z \)), the distribution of \( T_{CI} \) given by \( \theta \) (resp. of \( T_{CI} \) given by \( \theta \)) can be approximated by the \( \Gamma(k, \theta) \) distribution:

\[
p(t) = \frac{t^{k-1} e^{-t/\theta}}{\theta^k \Gamma(k)},
\]

where \( k = \text{E}^2(T_{UI}) / \text{Var}(T_{UI}) \) and \( \theta = \text{Var}(T_{UI}) / \text{E}(T_{UI}) \) in the unconditional case, and \( k = \text{E}^2(T_{CI}) / \text{Var}(T_{CI}) \) and \( \theta = \text{Var}(T_{CI}) / \text{E}(T_{CI}) \) in the conditional case. The means and variances of \( T_{UI} \) and \( T_{CI} \) on the given sample \( D \) are given in the following proposition.

**Proposition 6**

i. Under the null hypothesis that \( X \) and \( Y \) are independent, on the given sample \( D \), we have

\[
\text{E}(T_{UI}|D) = \frac{1}{n^2} \text{Tr}(\tilde{K}_X) \cdot \text{Tr}(\tilde{K}_Y),
\]

\[
\text{Var}(T_{UI}|D) = \frac{2}{n^4} \text{Tr}(\tilde{K}_X^2) \cdot \text{Tr}(\tilde{K}_Y^2).
\]

ii. Under the null hypothesis of \( X \perp Y|Z \), we have

\[
\text{E}(T_{CI}|D) = \frac{1}{n} \text{Tr}(\tilde{w} \tilde{w}^T),
\]

\[
\text{Var}(T_{CI}|D) = \frac{2}{n^2} \text{Tr}(\tilde{w} \tilde{w}^T)^2.
\]

### 3.5 Practical issues: On determination of the hyperparameters

In unconditional independence testing, the hyperparameters are the kernel widths for constructing the kernel matrices \( \tilde{K}_X \) and \( \tilde{K}_Y \). We found that the performance is robust to these parameters within a certain
range; in particular, as in [Gretton et al., 2008], we use the median of the pairwise distances of the points.

In KCI-test, the hyperparameters fall into three categories. The first one includes the kernel widths used to construct the kernel matrices $\tilde{K}_X$ and $\tilde{K}_Y$, which are used later to form the matrices $\tilde{K}_{X|Z}$ and $\tilde{K}_{Y|Z}$ according to (11) and (12). The values of the kernel widths should be such that one can use the information in the data effectively. In our experiments we normalize all variables to unit variance, and use some empirical values for those kernel widths: they are set to 0.8 if the sample size $n \leq 200$, to 0.3 if $n > 1200$, or to 0.5 otherwise. We found that this simple setting always works well in all our experiments.

The other two categories contain the kernel widths for constructing $\tilde{K}_Z$ and the regularization parameter $\varepsilon$, which are needed in calculating $\tilde{R}_Z$ in (10). These parameters should be selected carefully, especially when the conditioning set $Z$ is large. If $\varepsilon$ is too large, the corresponding regression functions $h_f^*$ and $h_g^*$ may underfit, resulting in a large probability of Type I errors (where the CI hypothesis is incorrectly rejected).

On the contrary, if they are too small, these regression function may overfit, which may increase the probability of Type II errors (where the CI hypothesis is incorrectly rejected). Moreover, if $\varepsilon$ is too small, $\tilde{R}_Z$ in (10) tends to vanish, resulting in very small values in $\tilde{K}_{X|Z}$ and $\tilde{K}_{Y|Z}$, and then the performance may be deteriorated by rounding errors.

We found that when the dimensionality of $Z$ is small (say, one or two variables), the proposed method works well even with some simple empirical settings (say, $\varepsilon = 10^{-3}$, and the kernel width for constructing $\tilde{K}_Z$ equals half of that for constructing $\tilde{K}_X$ and $\tilde{K}_Y$). When $Z$ contains many variables, to make the regression functions $h_f^*$ and $h_g^*$ more flexible, such that they could properly capture the information about $Z$ in $f(\tilde{X}) \in \mathcal{H}_X$ and $g'(Y) \in \mathcal{H}_Y$, respectively, we use separate regularization parameters and kernel widths for them, denoted by $\{\varepsilon_f, \sigma_f\}$ and $\{\varepsilon_g', \sigma_g'\}$. To avoid both overfitting and underfitting, we extend the Gaussian process (GP) regression framework to the multi-output case, and learn these hyperparameters by maximizing the total marginal likelihood. Details are skipped. The MATLAB source code is available at http://people.tuebingen.mpg.de/kzhang/KCI-test.zip.

4 Experiments

We apply the proposed method KCI-test to both synthetic and real data to evaluate its practical performance and compare it with CI\_PERM [Fukumizu et al., 2008]. It is also used for causal discovery.

4.1 On the effect of the dimensionality of $Z$ and the sample size

We examine how the probabilities of Type I and II errors of KCI-test change along with the size of the conditioning set $Z$ ($D = 1, 2, ..., 5$) and the sample size ($n = 200$ and $400$) in particular situations with simulations. We consider two cases as follows.

In Case I, only one variable in $Z$, namely, $Z_1$, is effective, i.e., other conditioning variables are independent from $X$, $Y$, and $Z_1$. To see how well the derived asymptotic null distribution approximates the true one, we examined if the probability of Type I errors is consistent with the significance level $\alpha$ that is specified in advance. We generated $X$ and $Y$ from $Z_1$ according to the post-nonlinear data generating procedure [Zhang and Hyvärinen, 2009]: they were constructed as $G(F(Z_1) + E)$ where $G$ and $F$ are random mixtures of linear, cubic, and tanh functions and are different for $X$ and $Y$, and $E$ is independent across $X$ and $Y$ and has various distributions. Hence $X \perp Y|Z$ holds. In our simulations $Z_i$ were i.i.d. Gaussian.

A good test is expected to have small probability of Type II errors. To see how large it is for KCI-test, we also generated the data which do not follow $X \perp Y|Z$ by adding the same variable to $X$ and $Y$ produced above. We increased the dimensionality of $Z$ and the sample size $n$, and repeated the CI tests 1000 random replications. Figure 1 (a,b) plots the resulting probability of Type I errors and that of Type II errors at the significance level $\alpha = 0.01$ and $\alpha = 0.05$, respectively.

In Case II, all variables in the conditioning set $Z$ are effective in generating $X$ and $Y$. We first generated the independent variables $Z_i$, and then similarly to Case I, to examine Type I errors, $X$ and $Y$ were generated as $G(\sum_i F_i(Z_i) + E)$. To examine Type II errors, we further added the same variable to $X$ and $Y$ such that they became conditionally dependent given $Z$. Figure 1 (c,d) gives the probabilities of Type I and Type II errors of KCI-test obtained on 1000 random replications.

One can see that with the derived null distribution, the 1% and 5% quantiles are approximated very well for both sample sizes, since the resulting probabilities of Type I errors are very close to the significance levels. Gamma approximation tends to produce slightly larger probabilities of Type I errors, meaning that the two-parameter Gamma approximation may have a slightly lighter tail than the true null distribution. With a fixed significance level $\alpha$, as $D$ increases, the probability of Type II errors always increases. This is intuitively reasonable: due to the finite sample size effect, as the conditioning set becomes larger and larger, X and Y tend to be considered as conditionally inde-
it is far less sensitive to $D$.

Figure 2: Results of CIPERM shown for comparison. (a) Probabilities of both Type I and II errors in Case II (dashed line: $n = 200$; solid line: $n = 400$). (b) Average CPU time taken by KCI-test and CIPERM in log scale. Note that for KCI-test, we include both the time for simulating the null distribution and that for Gamma approximation.

4.2 Application in causal discovery

CI tests are frequently used in problems of causal inference: in those problems, one assumes that the true causal structure of $n$ random variables $X_1, \ldots, X_n$ can be represented by a directed acyclic graph (DAG) $\mathcal{G}$. More specifically, the causal Markov condition assumes that the joint distribution satisfies all CIs that are imposed by the true causal graph (note that this is an assumption about the physical generating process of the data, not only about their distribution). So-called constraint-based methods like the PC algorithm (Spirtes et al., 2001) make the additional assumption of faithfulness (i.e., the joint distribution does not allow any CIs that are not entailed by the Markov condition) and recover the graph structure by exploiting the (conditional) independences that can be found in the data. Obviously, this is only possible up to Markov equivalence classes, which are sets of graphs that impose exactly the same independences and CIs. It is well known that small mistakes at the beginning of the algorithm (e.g., missing an independence relation) may lead to significant errors in the resulting DAG. Therefore the performance of those methods relies heavily on (conditional) independence testing methods.

For continuous data, the PC algorithm can be applied son, we always used GP for hyperparameter learning.

As claimed in Sec. 3.3, when $D$ is not high, say, $D = 1$ or 2, even fixed empirical values of the hyperparameters work very well. However, for consistency of the compar-
using partial correlation or mutual information. The
former assumes linear relationships and Gaussian dis-
tributions, and the latter does not lead to a significance
test. Sun et al. (2007) and Tillman et al. (2009) pro-
pose to use ClPERM for CI testing in PC. Based on
the promising results of KCI-test, we propose to also
apply it to causal inference.

### 4.2.1 Simulated data

We generated data from a random DAG $\mathcal{G}$. In partic-
ular, we randomly chose whether an edge exists and
sampled the functions that relate the variables from
a Gaussian Process prior. We sampled four random
variables $X_1, \ldots, X_4$ and allowed arrows from $X_i$ to
$X_j$ only for $i < j$. With probability 0.5 each possi-
ble arrow is either present or absent. If arrows exist,
from $X_1$ and $X_3$ to $X_4$, say, we sample $X_4$ from
a Gaussian Process with mean function $U_1 \cdot X_1 + U_3 \cdot X_3$
(with $U_1, U_3 \sim U[-2; 2]$) and a Gaussian kernel (with
each dimension randomly weighted between 0.1 and
0.6) plus a noise kernel as the covariance function. For
significance level 0.01 and sample sizes between 100
and 700 we simulated 100 DAGs, and checked how of-
ten the different methods infer the correct Markov
equivalence class. Figure 3 shows how often PC based
on KCI-test, ClPERM, or partial correlation recovered
the correct Markov equivalence class. PC based on
KCI-test gives clearly the best results.

4.2.2 Real data

We applied our method to continuous variables in the
Boston Housing data set, which is available at the UCI
Repository [Asuncion and Newman 2007]. Due to the
large number of variables we choose the significance
level to be 0.001, as a rough way to correct for mul-
tiple testing. Figure 4 shows the results for PC using
ClPERM (PC_{ClPERM}) and KCI-test (PC_{KCI-test}). For
conciseness, we report them in the same figure: the red
arrows are the ones inferred by PC_{ClPERM} and all solid
lines show the result by PC_{KCI-test}. Ergo, red solid
lines were found by both methods. Please refer to the
data set for the explanation of the variables. Although
one can argue about the ground truth for this data set,
we regard it as promising that our method finds links
between number of rooms (RM) and median value of
houses (MED) and between non-retail business (IND)
and nitric oxides concentration (NOX). The latter is
also missing in the result on these data given by Mar-
garitis (2005); instead their method gives some dubi-
ous links like crime rate (CRI) to nitric oxides (NOX),
for example.

Figure 4: Outcome of the PC algorithm applied to the
continuous variables of the Boston Housing Data Set
(red lines: PC_{ClPERM}, solid lines: PC_{KCI-test}).

### 5 Conclusion

We proposed a novel method for conditional indepen-
dence testing. It makes use of the characterization of
conditional independence in terms of uncorrelated-
ness of functions in suitable reproducing kernel Hilbert
spaces, and the proposed test statistic can be easily
 calculated from the kernel matrices. We derived its
distribution under the null hypothesis of conditional
independence. This distribution can either be gener-
ated by Monte Carlo simulation or approximated by
a two-parameter Gamma distribution. Compared to
discretization-based conditional independence testing
methods, the proposed one exploits more complete in-
formation of the data, and compared to those methods
defining the test statistic in terms of the estimated con-
ditional densities, calculation of the proposed statistic
involves less random errors. We applied the method
on both simulated and real world data, and the re-
sults suggested that the method outperforms existing
techniques in both accuracy and speed.

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### Appendix: Proofs

The following lemmas will be used in the proof of The-
orem 3. Consider the eigenvalues $\lambda_{X,i}$ and eigenvectors
$V_{x,i}$ of the kernel matrix $K_X$ and the eigenvalues $\lambda_{X,i}^*$
and normalized eigenfunctions $u_{X,i}$ (which have unit variance) of the corresponding kernel $k_X$. Let $\mathbf{x}$ be a fixed-size subsample of $\mathbf{x}$. The following lemma is due to Theorems 3.4 and 3.5 of Baker (1977).

**Lemma 7** i. $\frac{1}{n}\lambda_{X,i}$ converge in probability to $\lambda_{X,i}$.

ii. For the single eigenvalue $\lambda_{X,i}$, whose algebraic multiplicity is 1, $\sqrt{n}V_{X,i}(\mathbf{x})$ converges in probability to $u_{X,i}(\mathbf{x})$, where $V_{X,i}(\mathbf{x})$ denotes the values of $V_{x,i}$ corresponding to the subsample $\mathbf{x}$.

Suppose that the eigenvalues $\lambda_{X,k}, \lambda_{X,k+1}, \ldots, \lambda_{X,k+r}$ ($r \geq 1$) are the same and different from any other eigenvalue. The corresponding eigenfunctions are then non-unique. Denote by $\tilde{u}_{X,k+k+r} \triangleq (u_{X,k}, u_{X,k+1}, \ldots, u_{X,k+r})^T$ an arbitrary vector of the eigenfunctions corresponding to this eigenvalue with multiplicity $r + 1$. Let $S_{X,k}$ be the space of $\tilde{u}_{X,k+k+r}$.

**Lemma 8** i. The distance from $\sqrt{n}V_{X,k+1}(\mathbf{x})$ ($0 \leq q \leq r$) to the corresponding points in the space $S_{X,k}$ converges to zero in the following sense:

$$\inf \{ ||\sqrt{n}V_{X,k+q}(\mathbf{x}) - u^e(\mathbf{x})|| : u^e(\mathbf{x}) \in S_{X,k} \} \to 0, \text{ as } n \to \infty.$$  

ii. There exists an $(r + 1) \times (r + 1)$ orthogonal matrix $P_n$, which may depend on $n$, such that $\sqrt{n}P_n \cdot [V_{X,k}(\mathbf{x}), V_{X,k+1}(\mathbf{x}), \ldots, V_{X,k+r}(\mathbf{x})]^T$ converges in probability to $\tilde{u}_{X,k+k+r}(\mathbf{x})$.

Item (i.) in the above lemma is a reformulation of Theorem 3.6 of Baker (1977), while item (ii.) is its straightforward consequence.

The following lemma is a reformulation of Theorem 4.2 of Billingsley (1999).

**Lemma 9** Let $\{A_{L,n}\}$ be a double sequence of random variables with indices $L$ and $n$, $\{B_L\}$ and $\{C_n\}$ sequences of random variables, and $D$ a random variable. Assume that they are defined in a separable probability space. Suppose that, for each $L$, $A_{L,n} \xrightarrow{d} B_L$ as $n \to \infty$ and that $B_L \xrightarrow{d} D$ as $L \to \infty$. Suppose further that $\lim \sup_l \sum_{n=1}^{\infty} F(\{A_{L,n} - C_n\} \geq \varepsilon) = 0$ for each positive $\varepsilon$. Then $C_n \xrightarrow{d} D$ as $n \to \infty$.

**Sketch of proof of Theorem 3** We first define a $L \times L$ block-diagonal matrix $P_X$ as follows. For simple eigenvalues $\lambda_{X,i}$, $P_{X,ii} = 1$, and all other entries in the $i$th row and column are zero. For the eigenvalues with multiplicity $r + 1$, say, $\lambda_{X,k}, \lambda_{X,k+1}, \ldots, \lambda_{X,k+r}$, the corresponding main diagonal block from the $k$th row to the $(k + r)$th row of $P_X$ is an orthogonal matrix. According to Lemmas 7 and 8 and the continuous mapping theorem (CMT) of Mann and Wald (1943), which states that continuous functions are limit-preserving even if their arguments are sequences of random variables, there exists $P_X$, which may depend on $n$, such that $P_X \cdot [\psi_1(x_1), \ldots, \psi_L(x_1)]^T \to [\sqrt{\lambda_{X,1}u_{X,1}(x_1)}, \ldots, \sqrt{\lambda_{X,L}u_{X,L}(x_1)}]^T$ in probability as $n \to \infty$. Similarly, we can define the orthogonal matrix $P_Y$ such that $P_Y \cdot [\phi_1(y_1), \ldots, \phi_L(y_1)]^T \to [\sqrt{\lambda_{Y,1}u_{Y,1}(y_1)}, \ldots, \sqrt{\lambda_{Y,L}u_{Y,L}(y_1)}]^T$ in probability as $n \to \infty$.

Let $v_i$ be the random vector obtained by stacking the random $L \times L$ matrix $P_X M_i P_Y^T$ with $M_{ij,t} = \psi_i(x_1) \cdot \phi_j(y_1)$ as the $(i,j)$th entry of $M_t$. One can see that

$$\frac{1}{n} \sum_{i=1}^{n} ||v_i||^2 = \frac{1}{n} \sum_{i=1}^{n} \text{Tr}(P_X M_i P_Y^T \cdot P_Y M_i^T P_Y^T) = \frac{1}{n} \sum_{i=1}^{n} \text{Tr}(M_i M_i^T) = \sum_{i,j=1}^{L} S_{ij}^2. \quad (15)$$

Again, according to CMT, one can see that $v_i$ converges in probability to $w_i$. Furthermore, according to the vector-valued CLT (Eicker 1960), as $x_i$ and $y_i$ are i.i.d., the vector $\frac{1}{\sqrt{n}} \sum_{t=1}^{n} w_t$ converges in distribution to a multivariate normal distribution as $n \to \infty$. Because of CMT, $\frac{1}{\sqrt{n}} \sum_{t=1}^{n} v_t$ then converges to the same normal distribution as $n \to \infty$. As $f(X) \in \mathcal{H}_X$ and $g(Y) \in \mathcal{H}_Y$ are uncorrelated, we know that $u_{X,i}$ and $u_{Y,j}$ are uncorrelated, and consequently the mean of this normal distribution is $E(w_i) = 0$. The covariance is $\Sigma = \text{Cov}(\frac{1}{\sqrt{n}} \sum_{t=1}^{n} w_t) = \text{Cov}(w_i) = E(w_i w_i^T)$.

Assume that we have the EVD decomposition $\Sigma = \mathbf{A} \mathbf{W} \mathbf{A}^T \mathbf{V}^T$, where $\mathbf{A}$ is the diagonal matrix containing non-negative eigenvalues $\lambda_i^*$. Let $v_i^\prime = \mathbf{V}^T w_i$. Clearly $\frac{1}{\sqrt{n}} \sum_{t=1}^{n} v_i^\prime$ follows $\mathcal{N}(0, \mathbf{A} \mathbf{W} \mathbf{A}^T \mathbf{V}^T)$ asymptotically. That is,

$$\frac{1}{n} \sum_{t=1}^{n} ||v_i||^2 = \frac{1}{n} \sum_{t=1}^{n} ||v_i^\prime||^2 \xrightarrow{d} \sum_{k=1}^{L} \lambda_{k}^* v_k^2 \quad (16)$$

Combining (15) and (16) gives (5).

If $X$ and $Y$ are independent, for $k \neq i$ or $l \neq j$, one can see that the non-diagonal entries of $\Sigma$ are

$$\mathbb{E}[\sqrt{\lambda_{X,i} \lambda_{Y,j} \lambda_{X,k} \lambda_{Y,l}} u_{X,i}(x_1) u_{Y,j}(y_1) u_{X,k}(x_1) u_{Y,l}(y_1)] = \sqrt{\lambda_{X,i} \lambda_{Y,j} \lambda_{X,k} \lambda_{Y,l}} \mathbb{E}[u_{X,i}(x_1) u_{X,k}(x_1)] \mathbb{E}[u_{Y,j}(y_1) u_{Y,l}(y_1)] = 0.$$  

The diagonal entries of $\Sigma$ are $\lambda_{X,i} \lambda_{Y,j}$, $\mathbb{E}[u_{X,i}(x_1)]^2$, which are also eigenvalues of $\Sigma$. Substituting this result into (5), one obtains (6).

Finally, consider Lemma 9 and let $A_{L,n} \triangleq \sum_{i,j=1}^{L} S_{ij}^2$, $B_L \triangleq \sum_{k=1}^{L} \lambda_{k}^* v_k^2$, $C_n \triangleq \sum_{i,j=1}^{L} S_{ij}^2$, and $D \triangleq \sum_{k=1}^{L} \lambda_{k}^* v_k^2$. One can then see that $\sum_{i,j=1}^{L} S_{ij}^2 \xrightarrow{d}$
\[ \sum_{k=1}^{\infty} \tilde{\lambda}_k z_{ij}^2 \text{ as } n \to \infty. \] That is, (5) also holds as \( L = n \to \infty. \) As a special case of (5), (6) also holds as \( L = n \to \infty. \)

**Sketch of proof of Theorem 4.** On the one hand, due to (7), we have \( T_{UI} = \sum_{i,j=1}^{n} S_{ij}^2 \), which, according to [6], converges in distribution to \( \sum_{i,j=1}^{\infty} \lambda_{X,i}^2 \lambda_{Y,j}^2 z_{ij}^2 \) as \( n \to \infty. \) On the other hand, by extending the proof of Theorem 1 of Gretton et al. (2009), one can show that \( \sum_{i,j=1}^{\infty} \lambda_{X,i}^2 \lambda_{Y,j}^2 z_{ij}^2 \) converges in probability as \( n \to \infty. \) That is, \( T_{UI} \) converges in probability to \( \sum_{i,j=1}^{\infty} \lambda_{X,i}^2 \lambda_{Y,j}^2 z_{ij}^2 \) as \( n \to \infty. \) Consequently, \( T_{UI} \) and \( T_{CI} \) have the same asymptotic distribution.

**Sketch of proof of Proposition 5.** Here we let \( X, \ Y, \ K_X, \) and \( K_Y \) in Theorem 3 be \( X, (Y, Z), \) \( K_X|Z, \) and \( K_Y|Z, \) respectively. According to (7) and Theorem 3, one can see that \( T_{CI} \) follows the distribution with one degree of freedom, we have \( \mathbb{E}(z_{ij}^2) = 1 \) and \( \mathbb{V}(z_{ij}^2) = 2. \) According to (6), we have \( \mathbb{E}(T_{UI}|D) = \frac{1}{n^2} \left( \sum_{i,j} \lambda_{X,i} \lambda_{Y,j} - \frac{1}{n} \sum_{i,j} \lambda_{X,i} \lambda_{Y,j} \right) = \frac{1}{n} \text{Tr}(K_X) \cdot \text{Tr}(K_Y) \). Furthermore, bearing in mind that \( z_{ij}^2 \) are independent variables across \( i \) and \( j, \) and recalling \( \text{Tr}(K_X^2) = \sum_{i} \lambda_{X,i}^2 \), one can see that \( \mathbb{V}(T_{UI}|D) = \frac{1}{n^2} \left( \sum_{i,j} \lambda_{X,i}^2 \lambda_{Y,j} \right) \mathbb{V}(z_{ij}^2) = \frac{1}{n^2} \sum_{i,j} \lambda_{X,i}^2 \lambda_{Y,j} = \frac{2}{n^2} \text{Tr}(K_X^2) \cdot \text{Tr}(K_Y^2). \) Consequently (i) is true.

Similarly, from (14), one can calculate the mean \( \mathbb{E}(T_{CI}|D) \) and variance \( \mathbb{V}(T_{CI}|D), \) as given in (ii).

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