

General Frameworks for Conditional Two-Sample Testing

Seongchan Lee^{1*}

Suman Cha^{1*}

Ilmun Kim²

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Abstract

We study the problem of conditional two-sample testing, which aims to determine whether two populations have the same distribution after accounting for confounding factors. This problem commonly arises in various applications, such as domain adaptation and algorithmic fairness, where comparing two groups is essential while controlling for confounding variables. We begin by establishing a hardness result for conditional two-sample testing, demonstrating that no valid test can have significant power against any single alternative without proper assumptions. We then introduce two general frameworks that implicitly or explicitly target specific classes of distributions for their validity and power. Our first framework allows us to convert any conditional independence test into a conditional two-sample test in a black-box manner, while preserving the asymptotic properties of the original conditional independence test. The second framework transforms the problem into comparing marginal distributions with estimated density ratios, which allows us to leverage existing methods for marginal two-sample testing. We demonstrate this idea in a concrete manner with classification and kernel-based methods. Finally, simulation studies are conducted to illustrate the proposed frameworks in finite-sample scenarios.

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*These authors contributed equally to this work.

¹Department of Statistics and Data Science, Yonsei University, Seoul, South Korea.

²Department of Statistics and Data Science, Department of Applied Statistics, Yonsei University, Seoul, South Korea.

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1 Introduction

This paper addresses the problem of testing for equivalence between two conditional distributions, namely conditional two-sample testing. Statistical methods for this problem have important applications across diverse fields such as domain adaptation and algorithmic fairness. In domain adaptation, for instance, this methodology can serve as a formal framework to validate the covariate shift assumption, where the conditional distribution of Y given X remains unchanged, while the marginal distributions of X may differ. By confirming this assumption, practitioners can effectively re-weight the training data according to the marginal density ratio regarding X , which potentially leads to improved predictive performance and better adaptation to new domains (Shimodaira, 2000; Sugiyama et al., 2007a,b). Moreover, in algorithmic fairness, conditional two-sample testing plays a role in detecting and mitigating biases. In particular, it helps identify whether a certain machine learning model unfairly favors or disadvantages specific groups based on demographic characteristics such as age, gender, or ethnicity (Hardt et al., 2016; Barocas et al., 2023). Conditional two-sample testing also finds applications beyond machine learning. In genomics, for example, scientists seek to identify differences in genetic distributions conditional on various factors such as disease status and environmental exposures (Virolainen et al., 2022; Wu et al., 2023). This methodology aids scientists in understanding the genetic basis of diseases and in developing strategies for personalized medicine by providing a rigorous framework for comparing conditional distributions.

1.1 Problem Setup

With the practical motivation in mind, we now formally set up the problem. Given $n_1, n_2 \in \mathbb{N}$, suppose we observe two mutually independent samples

$$\{(X_i^{(1)}, Y_i^{(1)})\}_{i=1}^{n_1} \stackrel{\text{i.i.d.}}{\sim} P_{XY}^{(1)} \quad \text{and} \quad \{(X_i^{(2)}, Y_i^{(2)})\}_{i=1}^{n_2} \stackrel{\text{i.i.d.}}{\sim} P_{XY}^{(2)},$$

where $P_{XY}^{(1)}$ and $P_{XY}^{(2)}$ are joint distributions supported on some generic product space $\mathcal{X} \times \mathcal{Y}$. Let $P_{Y|X}^{(1)}$ and $P_{Y|X}^{(2)}$ denote the conditional distributions of $Y^{(1)} | X^{(1)}$ and $Y^{(2)} | X^{(2)}$, respectively. Similarly, let $P_X^{(1)}$ and $P_X^{(2)}$ denote the marginal distributions of $X^{(1)}$ and $X^{(2)}$, respectively. Given these two samples, our goal is to test the equality of two conditional distributions

$$H_0 : P_X^{(1)} \{P_{Y|X}^{(1)}(\cdot | X) = P_{Y|X}^{(2)}(\cdot | X)\} = 1 \quad \text{versus} \quad H_1 : P_X^{(1)} \{P_{Y|X}^{(1)}(\cdot | X) \neq P_{Y|X}^{(2)}(\cdot | X)\} > 0, \quad (1)$$

where $P_X^{(j)}(\cdot | x)$ denotes the conditional distribution of $Y^{(j)}$ given $X^{(j)} = x$ for $j \in \{1, 2\}$. In other words, we are interested in determining whether two populations have the same distribution after controlling for potential confounding variables. Throughout this paper, we assume that $P_X^{(1)}$ and $P_X^{(2)}$ have the same support, satisfying $P_X^{(1)} \ll P_X^{(2)}$ and $P_X^{(2)} \ll P_X^{(1)}$ where the symbol \ll denotes absolute continuity. Since $P_X^{(1)}$ and $P_X^{(2)}$ have the same support, the above hypotheses (1) for conditional two-sample testing can be equivalently defined using $P_X^{(2)}$ instead of $P_X^{(1)}$.

As pointed out by Boeken and Mooij (2021) and Yan and Zhang (2022), conditional two-sample testing is closely connected to conditional independence testing. To illustrate this connection, we introduce a binary variable $Z \in \{1, 2\}$, and see that the conditional independence between Y and Z given X is equivalently expressed as

$$Y \perp\!\!\!\perp Z | X \iff Y | X, Z = 1 \stackrel{d}{=} Y | X, Z = 2, \quad (2)$$

where the symbol $\stackrel{d}{=}$ denotes equality in distribution. This equivalence enables us to convert the problem of conditional two-sample testing to that of conditional independence testing based on the datasets $\{(Y_i, X_i) : Z_i = 1\}$ and $\{(Y_i, X_i) : Z_i = 2\}$. Consequently, we can leverage various existing methods for conditional independence testing to tackle conditional two-sample testing. However, prior work has not rigorously explored this approach, and indeed Yan and Zhang (2022) claim that it is not a sensible approach as the variable Z in the conditional two-sample problem is deterministic. Specifically, letting $n = n_1 + n_2$, $\sum_{i=1}^n \mathbf{1}(Z_i = 1)$ and $\sum_{i=1}^n \mathbf{1}(Z_i = 2)$ correspond to the sample sizes for two populations (i.e., n_1 and n_2), which are fixed in advance for the conditional two-sample problem. Therefore, a gap remains in rigorously connecting these seemingly similar, yet distinct, problems.

1.2 An Overview of Our Results

In this work, we make several contributions to the field of conditional two-sample testing. First, we reaffirm that comparing conditional distributions is intrinsically more difficult than comparing marginal distributions. For marginal two-sample testing, one can design permutation tests that control the type I error, while being powerful against certain alternatives (e.g., Kim et al., 2022a). However, we show that this is not the case for conditional two-sample testing. Our result (Theorem 1) proves that any valid conditional two-sample test has power at most equal to its size against any single alternative if the type of a conditional random vector is continuous. This is reminiscent of the negative result for conditional independence testing proved in Shah and Peters (2020). It is worth highlighting, however, that their negative result does not directly imply our Theorem 1. The proof of Shah and Peters (2020) relies on the assumption that the data $\{(X_i, Y_i, Z_i)\}_{i=1}^n$ are i.i.d., which does not hold in our setup as $\sum_{i=1}^n \mathbf{1}(Z_i = 1)$ and $\sum_{i=1}^n \mathbf{1}(Z_i = 2)$ are deterministic numbers. We handle this distinction through a concentration argument and show that conditional two-sample testing is as difficult as conditional independence testing. This negative result naturally motivates additional assumptions that make the problem feasible.

Our next contribution is to introduce two general frameworks for conditional two-sample testing. The first framework effectively addresses the issue pointed out by Yan and Zhang (2022). In particular, we develop a generic method that converts any conditional independence test into a conditional two-sample test. This general method directly transfers the asymptotic properties of a conditional independence test computed using $\{(X_i, Y_i, Z_i)\}_{i=1}^n \stackrel{\text{i.i.d.}}{\sim} P_{XYZ}$ to the setting of conditional two-sample testing (Theorem 2). At the heart of this approach is the concentration property of a Binomial random variable to its mean, which facilitates the effective construction of i.i.d. samples drawn from P_{XYZ} (see Algorithm 1). This development paves a way to leverage any existing methods for conditional independence testing in the literature, thereby expanding the range of tools available to practitioners for conducting two-sample tests.

The second framework that we introduce is based on density ratio estimation. To elaborate, let us assume that $P_X^{(1)}$ and $P_X^{(2)}$ have density functions $f_X^{(1)}$ and $f_X^{(2)}$ with respect to some base measure, and

similarly $P_{Y|X}^{(1)}(\cdot | x)$ and $P_{Y|X}^{(2)}(\cdot | x)$ have density functions $f_{Y|X}^{(1)}(\cdot | x)$ and $f_{Y|X}^{(2)}(\cdot | x)$, respectively. Then for all $x, y \in \mathcal{X} \times \mathcal{Y}$, we have the identity:

$$f_{Y|X}^{(1)}(y | x) = f_{Y|X}^{(2)}(y | x) \iff f_{YX}^{(1)}(y, x) = \frac{f_X^{(1)}(x)}{f_X^{(2)}(x)} f_{YX}^{(2)}(y, x) := f_{YX}(y, x), \quad (3)$$

where $f_{YX}^{(1)}$ is the joint density function of $(Y^{(1)}, X^{(1)})$ such that $f_{YX}^{(1)}(y, x) = f_{Y|X}^{(1)}(y | x) f_X^{(1)}(x)$, and $f_{YX}^{(2)}$ is similarly defined for $(Y^{(2)}, X^{(2)})$. The above equivalence (3) allows us to transform the problem of testing for conditional distributions into the one that compares marginal distributions with densities $f_{YX}^{(1)}$ and $f_{YX}^{(2)}$. The latter problem has been extensively studied with various methods, ranging from classical approaches such as Hotelling's test to modern methods such as kernel maximum mean discrepancy (Gretton et al., 2012; Liu et al., 2020; Schrab et al., 2023) and machine learning-based approaches (e.g., Lopez-Paz and Oquab, 2017; Kim et al., 2019, 2021; Hediger et al., 2022). The issue, however, is that we do not observe samples from f_{YX} but from $f_{YX}^{(2)}$. Therefore, the success of this framework relies on how accurate one can estimate the density ratio

$$r_X(x) := \frac{f_X^{(1)}(x)}{f_X^{(2)}(x)}, \quad (4)$$

and incorporate it into the procedure to fill the gap between f_{YX} and $f_{YX}^{(2)}$. We demonstrate this methodology focusing on a classification-based test in Section 4.1 and a kernel-based test in Section 4.2.

1.3 Literature Review

As mentioned earlier, conditional two-sample testing has a wide range of applications in various fields, including machine learning, genetics and economics, where it is important to compare two samples controlling for confounding variables. Despite its broad range of applications and significance, there has been limited research dedicated to tackling this fundamental problem. Similar problems, on the other hand, have been explored in the literature such as testing for the equality of conditional moments (Hall and Hart, 1990; Kulasekera, 1995; Kulasekera and Wang, 1997; Fan and Lin, 1998; Neumeyer and Dette, 2003; Pardo-Fernández et al., 2015) and goodness-of-fit testing for pre-specified conditional distributions (Andrews, 1997; Zheng, 2000; Fan et al., 2006). These methods aim to facilitate the comparison of specific aspects of a distribution such as the conditional mean or second moments, rather than the entire distribution. Our research, however, is centered on nonparametric comparisons of two conditional distributions. This approach is of great importance as it enables a more comprehensive comparison of distributions, capturing differences that may not be evident when only specific moments or pre-specified models are compared.

It is only in recent years that conditional two-sample testing has gained attention, with several novel methods being proposed. Yan and Zhang (2022), for instance, proposed a method that extends unconditional energy distance to its conditional counterpart. They demonstrated that many key properties of the unconditional energy distance are retained in the conditional version. Moreover, they proposed a bootstrap procedure to calibrate their test statistic. To the best of our knowledge, however, the validity of their test remains unexplored, and the $O(n^4)$ time complexity of their algorithm poses a bottleneck to its practical application.

As another example, Hu and Lei (2024) built on the idea of conformal prediction and introduced a nonparametric conditional two-sample test using a weighted rank-sum statistic. This approach involves estimating both marginal and conditional density ratios, and the validity of their method depends on the quality of these ratio estimators. As explained in Example 4, their test statistic can be viewed as an example of our general framework based on density ratio estimation. A more recent work by Chen and Lei (2024) extended the idea of Hu and Lei (2024), leveraging Neyman orthogonality to reduce the first-order bias for the asymptotic normality. As another closely related work, Chatterjee et al. (2024) introduced a kernel-based

conditional two-sample test using nearest neighbors. They considered the setting where a random sample $\{(X_i, Y_i^{(1)}, Y_i^{(2)})\}_{i=1}^n$ is generated from a joint distribution, i.e., the response variables $Y^{(1)}$ and $Y^{(2)}$ are conditioned on the same set of covariates X . This setting is notably different from that considered in the prior work (Yan and Zhang, 2022; Hu and Lei, 2024; Chen and Lei, 2024) as well as in our study, which consider potentially different covariates. Hence, the methods proposed by Chatterjee et al. (2024) are not directly comparable to ours.

As explained before, the first framework that we propose can be constructed based on essentially any conditional independence tests from the literature. The problem of testing for conditional independence has been extensively studied, resulting in a variety of methods to handle different scenarios and challenges. Shah and Peters (2020) proposed the Generalized Covariance Measure (GCM) whose validity depends on the performance of regression methods. Recent improvements to this method include the strategies such as weighting (Scheidegger et al., 2022) and applying GCM to a projected random vector (Lundborg et al., 2022; Chakraborty et al., 2024). Other notable methodologies for conditional independence testing include kernel-based tests (Zhang et al., 2011; Doran et al., 2014; Strobl et al., 2019; Pogodin et al., 2024), binning-based tests (Neykov et al., 2021; Kim et al., 2022b; Neykov et al., 2023), regression-based tests (Dai et al., 2022; Williamson et al., 2023) and tests under the model-X framework (Candes et al., 2018; Berrett et al., 2020; Liu et al., 2022; Tansey et al., 2022). Our method can leverage these developments to effectively solve the problem of conditional two-sample testing.

Our second framework can benefit from extensive research done on density ratio estimation in the literature. A straightforward way of estimating density ratio is to first estimate individual density functions, and take their ratio as an estimate. However, this method tends to become unstable, especially in high-dimensional settings. To overcome this issue, Sugiyama et al. (2007b) and Tsuboi et al. (2009) developed methods that directly estimate density ratio without involving density estimation. Kanamori et al. (2010) compared different methods of density ratio estimation, and discussed their theoretical properties. Kanamori et al. (2009) reformulated the problem as a least-squares problem to provide a closed-form solution, whereas Liu et al. (2017) proposed trimmed density ratio estimation to improve stability and robustness by trimming extreme values. More recent advancements in density ratio estimation include Choi et al. (2021); Rhodes et al. (2020); Choi et al. (2022). As explained in Section 4, our approach uses density ratio estimation to deal with discrepancies between $f_{Y|X}^{(1)}$ and $f_{Y|X}$, and transforms the problem of comparing conditional distributions into that of comparing marginal distributions.

1.4 Organization

The rest of this paper is organized as follows. We begin with a hardness result for conditional two-sample testing in Section 2, which shows that no test can have power greater than its size against any alternative without additional assumptions. Section 3 presents our framework that converts tests for conditional independence into those for the equality of conditional distributions. Section 4 introduces another framework based on density ratio estimation. Numerical results illustrating the finite-sample performance of our methods are presented in Section 5, followed by the conclusion in Section 6. The proofs of the results omitted in the main text can be found in the appendix.

2 Hardness Result

Before introducing our frameworks, we present a fundamental hardness result for conditional two-sample testing. Specifically, for a continuous random vector X , our result demonstrates that any valid conditional two-sample test has no power against any alternative. This finding parallels the negative result established by Shah and Peters (2020) for conditional independence testing, and our proof builds crucially on their work. Given the connection established in (2), one might argue that their negative result directly applies to the two-sample problem. However, additional effort is required to make this connection concrete since the

sample sizes n_1 and n_2 are deterministic in our setting, which violates the i.i.d. assumption required in [Shah and Peters \(2020\)](#).

To state the result, let (X, Y, Z) be a random vector in $\mathbb{R}^{d_X+d_Y} \times \{1, 2\}$, and \mathcal{E} be the set of all distributions such that X and Y have marginal distributions, which are absolutely continuous with respect to the Lebesgue measure. In addition, Z has a discrete distribution supported on $\{1, 2\}$ with $\mathbb{P}(Z = 1) = \lambda_n \in (0, 1)$ where $\lambda_n := n_1/n$. Let $\mathcal{P}_0 \subset \mathcal{E}$ be the set of null distributions for conditional independence, i.e., distributions such that $Y \perp\!\!\!\perp Z | X$, and let $\mathcal{P}_1 = \mathcal{E} \setminus \mathcal{P}_0$ be the set of alternative distributions. For $M \in (0, \infty]$, let $\mathcal{E}_M \subseteq \mathcal{E}$ be the subset of all distributions where X and Y are supported within an ℓ_∞ ball of radius M . Set $\mathcal{P}_{0,M} = \mathcal{P}_0 \cap \mathcal{E}_M$ and $\mathcal{P}_{1,M} = \mathcal{P}_1 \cap \mathcal{E}_M$. Further denoting $N_1 := \sum_{i=1}^n \mathbf{1}(Z_i = 1)$ and $N_2 := \sum_{i=1}^n \mathbf{1}(Z_i = 2)$, the following theorem shows that no valid test ϕ for conditional two-sample testing has power greater than its size.

Theorem 1. *Let $n_1, n_2 \in \mathbb{N}$ with $n = n_1 + n_2$, $\alpha \in (0, 1)$ and $M \in (0, \infty]$. For $\{(X_i, Y_i, Z_i)\}_{i=1}^n \stackrel{\text{i.i.d.}}{\sim} P_{XYZ} := P$, consider a test $\phi : \{(X_i, Y_i, Z_i)\}_{i=1}^n \mapsto \{0, 1\}$. Suppose that ϕ controls the type I error at level α conditional on $N_1 = n_1$ and $N_2 = n_2$ as*

$$\sup_{P \in \mathcal{P}_{0,M}} \mathbb{E}_P[\phi | N_1 = n_1, N_2 = n_2] \leq \alpha.$$

Then the power of ϕ conditional on $N_1 = n_1$ and $N_2 = n_2$ is at most α for any $P \in \mathcal{P}_{1,M}$ as

$$\mathbb{E}_P[\phi | N_1 = n_1, N_2 = n_2] \leq \alpha.$$

Remark. We note that the type I error and the power of ϕ conditional on $N_1 = n_1$ and $N_2 = n_2$ in Theorem 1 are equivalent to those computed under the two-sample setting with deterministic sample sizes. In other words, the testing errors for conditional two-sample testing can be understood as the testing errors for conditional independence testing conditional on $N_1 = n_1$ and $N_2 = n_2$; thereby Theorem 1 implying the negative result for conditional two-sample testing. We also note that Theorem 1 only focuses on non-randomized tests for simplicity of presentation, but our proof also holds for randomized tests. The proof of Theorem 1 is provided below.

Proof of Theorem 1. For a constant $N \in \mathbb{N}$ greater than $2n$, we work with N i.i.d. random samples $\{(X_i, Y_i, Z_i)\}_{i=1}^N \stackrel{\text{i.i.d.}}{\sim} P$ and define $N'_1 = \sum_{i=1}^N \mathbf{1}(Z_i = 1)$ and $N'_2 = \sum_{i=1}^N \mathbf{1}(Z_i = 2)$, which follow $N'_1 \sim \text{Binomial}(N, \lambda_n)$ and $N'_2 \sim \text{Binomial}(N, 1 - \lambda_n)$, respectively. For $n_1, n_2 \in \mathbb{N}$ given in the theorem statement, define a good event $\mathcal{A} := \{N'_1 \geq n_1, N'_2 \geq n_2\}$, whose probability satisfies $\mathbb{P}(\mathcal{A}) \geq 1 - \mathbb{P}(N'_1 < n_1) - \mathbb{P}(N'_2 < n_2)$ by the union bound. Since $N \geq 2n$, we can ensure that $n_1 - N\lambda_n \leq -\frac{1}{2}N\lambda_n$ and thus

$$\begin{aligned} \mathbb{P}(N'_1 < n_1) &= \mathbb{P}(N'_1 - N\lambda_n < n_1 - N\lambda_n) \leq \mathbb{P}\left(N'_1 - N\lambda_n < -\frac{1}{2}N\lambda_n\right) \\ &\leq \mathbb{P}\left(|N'_1 - N\lambda_n| > \frac{1}{2}N\lambda_n\right) \leq \frac{4(1 - \lambda_n)}{N\lambda_n}, \end{aligned}$$

where the last inequality uses Chebyshev's inequality along with $N'_1 \sim \text{Binomial}(N, \lambda_n)$. We can similarly obtain that $\mathbb{P}(N'_2 < n_2) \leq \frac{4\lambda_n}{N(1 - \lambda_n)}$. Therefore, the probability of the good event \mathcal{A} is lower bounded as

$$\mathbb{P}(\mathcal{A}) \geq 1 - \frac{4(1 - \lambda_n)^2 + 4\lambda_n^2}{N\lambda_n(1 - \lambda_n)} \stackrel{\text{set}}{=} 1 - \varepsilon_N, \quad (5)$$

where $\varepsilon_N \rightarrow 0$ as $N \rightarrow \infty$.

Now consider a test ϕ that only uses $n_1 + n_2$ data points out of N samples. Importantly, this sample consists of n_1 observations from $\{(Y_i, X_i) : Z_i = 1\}$ and n_2 observations from $\{(Y_i, X_i) : Z_i = 2\}$, whenever $\mathbf{1}(\mathcal{A}) = 1$. If $\mathbf{1}(\mathcal{A}) = 0$, this test simply returns 0 (i.e., accept H_0). Moreover, we assume that this test satisfies $\sup_{P \in \mathcal{P}_{0,M}} \mathbb{E}_P[\phi | \mathcal{A}] \leq \alpha$. In fact, since ϕ only uses n_1 observations with $Z = 1$ and n_2 observations with $Z = 2$, the previous inequality implies $\sup_{P \in \mathcal{P}_{0,M}} \mathbb{E}_P[\phi | N_1 = n_1, N_2 = n_2] \leq \alpha$, i.e., ϕ is a valid level

α test for conditional two-sample testing. Based on the previous results, the type I error of ϕ constructed based on $\{(X_i, Y_i, Z_i)\}_{i=1}^N \stackrel{\text{i.i.d.}}{\sim} P$ fulfills

$$\sup_{P \in \mathcal{P}_{0,M}} \mathbb{E}_P[\phi] = \sup_{P \in \mathcal{P}_{0,M}} \mathbb{E}_P[\phi \mathbf{1}(\mathcal{A})] \leq \sup_{P \in \mathcal{P}_{0,M}} \mathbb{E}_P[\phi | \mathcal{A}] \leq \alpha.$$

This implies that ϕ is a valid test for conditional independence with size α . Therefore, for any $P \in \mathcal{P}_{1,M}$,

$$\begin{aligned} \mathbb{E}_P[\phi | \mathcal{A}](1 - \varepsilon_N) &\stackrel{\text{(i)}}{\leq} \mathbb{E}_P[\phi | \mathcal{A}] \mathbb{E}_P[\mathbf{1}(\mathcal{A})] \leq \mathbb{E}_P[\phi] \stackrel{\text{(ii)}}{\leq} \alpha \\ \iff \mathbb{E}_P[\phi | \mathcal{A}] &\stackrel{\text{(iii)}}{=} \mathbb{E}_P[\phi | N_1 = n_1, N_2 = n_2] \leq \frac{\alpha}{1 - \varepsilon_N}, \end{aligned}$$

where step (i) uses the inequality in (5), step (ii) holds by [Shah and Peters \(2020\)](#), Theorem 2 and Remark 4) and step (iii) uses the fact that ϕ only uses $n_1 + n_2$ observations as described before. Since N was an arbitrary number greater than or equal to $2n$ and $\varepsilon_N \rightarrow 0$ as $N \rightarrow \infty$, we can conclude that $\mathbb{E}_P[\phi | N_1 = n_1, N_2 = n_2] \leq \alpha$ for any $P \in \mathcal{P}_{1,M}$. \square

Theorem 1 clearly explains that it is necessary to impose additional assumptions (e.g., smoothness for distributions) in order to make the conditional two-sample problem feasible. In the next two sections, we explore two general frameworks, which implicitly or explicitly incorporate reasonable assumptions to address this problem. The first framework utilizes any conditional independence test and considers scenarios where this test performs well for verifying conditional independence. Conversely, the second framework assumes that the marginal density ratio r_X is well-behaved and can be estimated with high accuracy.

3 Approach via Conditional Independence Testing

In this section, we introduce our first framework that converts a conditional independence test to a conditional two-sample test, while maintaining the same asymptotic guarantees. The key idea is to construct a dataset $\mathcal{D}_{\tilde{n}}$ consisting of i.i.d. random vectors (Y, Z, X) of size \tilde{n} based on the given two samples $\{(Y_i^{(1)}, X_i^{(1)})\}_{i=1}^{n_1}$ and $\{(Y_i^{(2)}, X_i^{(2)})\}_{i=1}^{n_2}$. To achieve this, letting $n = n_1 + n_2$, we first draw a random variable \tilde{n}_1 from $\text{Binomial}(\tilde{n}, n_1/n)$ where \tilde{n} is set to be smaller than n and $\tilde{n}/n \rightarrow 1$. Since a Binomial random variable is highly concentrated around its mean, we can guarantee that $\tilde{n}_1 \leq n_1$ and $\tilde{n}_2 := \tilde{n} - \tilde{n}_1 \leq n_2$ with high probability. If a bad event happens where either $\tilde{n}_1 > n_1$ or $\tilde{n}_2 > n_2$, making the construction of $\mathcal{D}_{\tilde{n}}$ infeasible, we simply accept the null hypothesis. This slightly inflates the type II error in finite-sample scenarios, but it is asymptotically negligible. A similar idea has been utilized in [Neykov et al. \(2021\)](#) in a different context to eliminate Poissonization for conditional independence testing.

Having constructed $\mathcal{D}_{\tilde{n}}$ consisting of i.i.d. random samples drawn from the joint distribution of (Y, Z, X) , we can now implement a conditional independence test based on $\mathcal{D}_{\tilde{n}}$, while retaining the same theoretical guarantees for conditional two-sample testing. Algorithm 1 summarizes this procedure, and the following theorem formally establishes its theoretical guarantees.

Theorem 2. Consider a class of distributions \mathcal{P} of (X, Y, Z) where Z takes a value among $\{1, 2\}$ with probability n_1/n and n_2/n , respectively, and let $\{(X_i, Y_i, Z_i)\}_{i=1}^n \stackrel{\text{i.i.d.}}{\sim} P \in \mathcal{P}$. For disjoint subclasses $\mathcal{P}_0 \subset \mathcal{P}$ and $\mathcal{P}_1 \subset \mathcal{P}$ and $\alpha \in (0, 1)$, assume that a test $\phi : \{(X_i, Y_i, Z_i)\}_{i=1}^n \mapsto \{0, 1\}$ satisfies

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_0} \mathbb{E}_P[\phi] \leq \alpha \quad \text{and} \quad \lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_1} \mathbb{E}_P[1 - \phi] = 0.$$

Denote the output of Algorithm 1 with $\varepsilon = o(1)$ as $\tilde{\phi} \in \{0, 1\}$ where $\tilde{\phi} = 1$ if and only if H_0 is rejected. Then it holds that

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_0} \mathbb{E}_P[\tilde{\phi} | N_1 = n_1, N_2 = n_2] \leq \alpha \quad \text{and} \quad \lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_1} \mathbb{E}_P[1 - \tilde{\phi} | N_1 = n_1, N_2 = n_2] = 0,$$

where we recall $N_1 = \sum_{i=1}^n \mathbf{1}(Z_i = 1)$ and $N_2 = \sum_{i=1}^n \mathbf{1}(Z_i = 2)$.

Algorithm 1 Converting a Conditional Independence Test into a Conditional Two-Sample Test

Require: Data $\{(Y_i^{(1)}, X_i^{(1)})\}_{i=1}^{n_1}$ and $\{(Y_i^{(2)}, X_i^{(2)})\}_{i=1}^{n_2}$ of sizes n_1 and n_2 , a conditional independence test ϕ for $H_0 : Y \perp\!\!\!\perp Z | X$ of (asymptotic) size $\alpha \in (0, 1)$, adjustment parameter $\varepsilon \in (0, 1)$

- 1: Draw $\tilde{n}_1 \sim \text{Binomial}(\tilde{n}, n_1/n)$ where $\tilde{n} = k^*n$ and $k^* = 1 - 3\log(\varepsilon)/(2n_1) - \sqrt{(1 - 3\log(\varepsilon)/(2n_1))^2 - 1}$ and set $\tilde{n}_2 = \tilde{n} - \tilde{n}_1$.
- 2: **if** $\tilde{n}_1 > n_1$ or $\tilde{n}_2 > n_2$ **then** Accept H_0 .
- 3: **else**
- 4: Merge $\{(X_i^{(1)}, Y_i^{(1)}, Z_i = 1)\}_{i=1}^{\tilde{n}_1}$ and $\{(X_i^{(2)}, Y_i^{(2)}, Z_i = 2)\}_{i=1}^{\tilde{n}_2}$, yielding $\mathcal{D}_{\tilde{n}} := \{(X_i, Y_i, Z_i)\}_{i=1}^{\tilde{n}}$.
- 5: Run a conditional independence test ϕ using $\mathcal{D}_{\tilde{n}}$ at level α , and Denote the resulting test as $\phi_{\tilde{n}}$.
- 6: **if** $\phi_{\tilde{n}} = 1$ **then** Reject H_0 **else** Accept H_0 .
- 7: **end if**
- 8: **end if**

Proof of Theorem 2. We may write $\tilde{\phi} = \mathbb{1}(\tilde{n}_1 \leq n_1)\mathbb{1}(\tilde{n}_2 \leq n_2)\phi_{\tilde{n}}$ where $\phi_{\tilde{n}}$ is defined as ϕ based on $\mathcal{D}_{\tilde{n}} = \{(X_i, Y_i, Z_i)\}_{i=1}^{\tilde{n}}$ in Algorithm 1. Now generate new i.i.d. samples $\tilde{\mathcal{D}}_{\tilde{n}} = \{(\tilde{X}_i, \tilde{Y}_i, \tilde{Z}_i)\}_{i=1}^{\tilde{n}} \stackrel{\text{i.i.d.}}{\sim} P$ independent of $\mathcal{D}_{\tilde{n}}$ and define as

$$\tilde{\phi}^\dagger = \mathbb{1}\left\{\sum_{i=1}^{\tilde{n}} \mathbb{1}(\tilde{Z}_i = 1) \leq n_1\right\} \mathbb{1}\left\{\sum_{i=1}^{\tilde{n}} \mathbb{1}(\tilde{Z}_i = 2) \leq n_2\right\} \phi_{\tilde{n}}^\dagger,$$

where $\phi_{\tilde{n}}^\dagger$ is defined as ϕ but based on $\tilde{\mathcal{D}}_{\tilde{n}}$. Observe that the conditional distribution of $\tilde{\phi}$ given $N_1 = n_1, N_2 = n_2$ is the identical to the marginal distribution of $\tilde{\phi}^\dagger$. Thus it can be seen that

$$\mathbb{E}_P[\tilde{\phi} | N_1 = n_1, N_2 = n_2] = \mathbb{E}_P[\tilde{\phi}^\dagger] \leq \mathbb{E}_P[\phi_{\tilde{n}}^\dagger],$$

for all P . In other words, we can effectively remove conditioning on $N_1 = n_1, N_2 = n_2$ by working with $\tilde{\phi}^\dagger$. Therefore, the first claim on type I error control follows.

Moving to the type II error, observe that

$$\begin{aligned} \mathbb{1}(\tilde{n}_1 \leq n_1)\mathbb{1}(\tilde{n}_2 \leq n_2)\phi_{\tilde{n}} &= \phi_{\tilde{n}} - \mathbb{1}(\tilde{n}_1 > n_1 \text{ or } \tilde{n}_2 > n_2)\phi_{\tilde{n}} \\ &\geq \phi_{\tilde{n}} - \mathbb{1}(\tilde{n}_1 > n_1) - \mathbb{1}(\tilde{n}_2 > n_2), \end{aligned}$$

by the union bound, which yields

$$\mathbb{E}_P[1 - \tilde{\phi} | N_1 = n_1, N_2 = n_2] = \mathbb{E}_P[1 - \tilde{\phi}^\dagger] \leq \mathbb{E}_P[1 - \phi_{\tilde{n}}^\dagger] + \mathbb{E}_P[\mathbb{1}(\tilde{n}_1 > n_1)] + \mathbb{E}_P[\mathbb{1}(\tilde{n}_2 > n_2)]. \quad (6)$$

In addition, letting $p = n_1/n$ in Lemma 1 of Appendix B, take $(1 + \delta)\tilde{n}n_1/n = k(1 + \delta)n_1 = n_1$, which gives $1 + \delta = k^{-1} \iff \delta = k^{-1} - 1$. Thus, by Lemma 1, we have

$$\mathbb{E}_P[\mathbb{1}(\tilde{n}_1 > n_1)] \leq \exp\left(-\frac{n_1 k(k^{-1} - 1)^2}{3}\right).$$

Letting the right-hand side equal ε and solving for $k \in (0, 1)$, we derive the form of k^* as presented in Algorithm 1, which shows that $\mathbb{E}_P[\mathbb{1}(\tilde{n}_1 > n_1)] \leq \varepsilon$. By symmetry, the same analysis holds for the inequality $\mathbb{E}_P[\mathbb{1}(\tilde{n}_2 > n_2)] \leq \varepsilon$. As a result, continuing from the inequality (6), we can upper bound the type II error of $\tilde{\phi}_\alpha$ as

$$\mathbb{E}_P[1 - \tilde{\phi} | N_1 = n_1, N_2 = n_2] \leq \mathbb{E}_P[1 - \phi_{\tilde{n}}^\dagger] + 2\varepsilon.$$

Since $\varepsilon = o(1)$, the above display proves the second claim on type II error control. This completes the proof of Theorem 2. \square

Our analysis in Theorem 2 is not limited to conditional two-sample testing, and it can be applied to marginal two-sample testing as well. Indeed, the problem of conditional two-sample testing becomes equivalent to the unconditional counterpart when X is degenerate (e.g., $X = 0$ with probability one). Thus, our algorithm serves as a generic method to convert unconditional independence tests to unconditional two-sample tests as well. We also mention that a specific form of k^* in Algorithm 1 is derived from the multiplicative Chernoff bound for a Binomial random variable (Lemma 1), which can be refined by numerically computing the tail probability of a Binomial random variable.

Despite its generality, one obvious drawback of Algorithm 1 is that it does not take the datasets $\{(Y_i^{(1)}, X_i^{(1)})\}_{i=\tilde{n}_1+1}^{n_1}$ and $\{(Y_i^{(2)}, X_i^{(2)})\}_{i=\tilde{n}_2+1}^{n_2}$ into account in the procedure when $\tilde{n}_1 < n_1$ and $\tilde{n}_2 < n_2$. It can be seen that the expected number of discarded samples, i.e., $\mathbb{E}[n - \tilde{n}_1 - \tilde{n}_2]$, is $O(\sqrt{n \log(1/\varepsilon)})$. This loss might degrade the performance in small-sample size regimes, but it can be negligible when n is sufficiently large and ε decreases slowly (see Appendices C.4 and C.5 for empirical support). Nevertheless, when a test statistic is sufficiently stable, the conclusion of Theorem 2 may hold without further modification of ϕ , meaning the conditional testing errors of ϕ are asymptotically equivalent to its marginal errors.

To illustrate this, we build upon the coupling argument presented by Chung and Romano (2013). First, draw $\bar{n}_1 \sim \text{Binomial}(n, n_1/n)$ and set $\bar{n}_2 = n - \bar{n}_1$. If $\bar{n}_1 > n_1$, draw $\bar{n}_1 - n_1$ additional samples $\{(Y_i^{(1)}, X_i^{(1)})\}_{i=n_1+1}^{\bar{n}_1}$ from $P_{XY}^{(1)}$. Otherwise, draw $\bar{n}_2 - n_2$ additional samples $\{(Y_i^{(2)}, X_i^{(2)})\}_{i=n_2+1}^{\bar{n}_2}$ from $P_{XY}^{(2)}$. In either case, set $\{(\tilde{X}_i, \tilde{Y}_i, \tilde{Z}_i)\}_{i=1}^n = \{X_i^{(1)}, Y_i^{(1)}, 1\}_{i=1}^{\bar{n}_1} \cup \{X_i^{(2)}, Y_i^{(2)}, 2\}_{i=1}^{\bar{n}_2}$, which can be viewed as i.i.d. draws from a joint distribution of (X, Y, Z) , after randomly permuting indices. When this newly constructed dataset is compared with the original dataset $\{(X_i, Y_i, Z_i)\}_{i=1}^n = \{X_i^{(1)}, Y_i^{(1)}, 1\}_{i=1}^{n_1} \cup \{X_i^{(2)}, Y_i^{(2)}, 2\}_{i=1}^{n_2}$, there are $|\bar{n}_1 - n_1|$ distinct data points with the expectation $\mathbb{E}[|\bar{n}_1 - n_1|] \leq \sqrt{n/4}$. This suggests that if a test statistic T is asymptotically invariant to \sqrt{n} -data perturbations, the asymptotic behavior of T remains consistent across both the original and the newly constructed datasets. However, the stability of T needs to be evaluated on a case-by-case basis. Below, we provide examples illustrating both stable and unstable cases.

Example 1 (Stable case). To simplify our presentation, consider a univariate case of $Y \in \mathbb{R}$, and assume $f(x) := \mathbb{E}[Y | X = x]$ and $g(x) := \mathbb{E}[Z | X = x]$ are known. Letting $R_i := \{Y_i - f(X_i)\}\{Z_i - g(X_i)\}$, the generalized covariance measure introduced by Shah and Peters (2020) is

$$T = \frac{\frac{1}{\sqrt{n}} \sum_{i=1}^n R_i}{\left\{ \frac{1}{n} \sum_{i=1}^n R_i^2 - \left(\frac{1}{n} \sum_{r=1}^n R_r \right)^2 \right\}^{1/2}},$$

and let \tilde{T} be similarly defined as T based on $\{(\tilde{X}_i, \tilde{Y}_i, \tilde{Z}_i)\}_{i=1}^n$. Focusing on the numerators of T and \tilde{T} , it can be seen that their difference is

$$\frac{1}{\sqrt{n}} \sum_{i=1}^n R_i - \frac{1}{\sqrt{n}} \sum_{i=1}^n \tilde{R}_i = \frac{1}{\sqrt{n}} \sum_{i=n_1+1}^{\bar{n}_1} (R_i - \tilde{R}_i) \cdot \mathbf{1}(\bar{n}_1 > n_1) + \frac{1}{\sqrt{n}} \sum_{i=\bar{n}_1+1}^{n_1} (R_i - \tilde{R}_i) \cdot \mathbf{1}(\bar{n}_1 \leq n_1).$$

Under the null hypothesis, the expectation of the difference is zero and the variance is bounded above by $1/\sqrt{n}$ up to a constant, provided that each Y_i has a finite second moment. Therefore, the difference of the numerators is asymptotically negligible. We can show similarly that the difference of the denominators is also asymptotically negligible as detailed in Appendix A.5. Putting things together concludes that T and \tilde{T} are asymptotically equivalent.

In the above example, we assumed that the conditional expectations f and g are known. In practice, f and g are estimated from the data, and the generalized covariance measure can become highly unstable when the estimators of f and g are themselves unstable.

Example 2 (Unstable case). Consider an extreme case where the estimators \hat{f} and \hat{g} are defined as follows: $\hat{f}(X_i) := Y_i \mathbf{1}(\sum_{i=1}^n Z_i = n_1) + f(X_i) \mathbf{1}(\sum_{i=1}^n Z_i \neq n_1)$ and $\hat{g}(X_i) := Z_i \mathbf{1}(\sum_{i=1}^n Z_i = n_1) + g(X_i) \mathbf{1}(\sum_{i=1}^n Z_i \neq n_1)$.

n_1). In this case, T is not well-defined as it takes the form 0/0 deterministically. On the other hand, when n_1 and n_2 are well-balanced (e.g., $n_1/n = 1/2$), the probability of the event $\sum_{i=1}^n \tilde{Z}_i = n_1$ converges to zero as n increases. Under such condition and assuming suitable moment conditions, the test statistic \tilde{T} based on $\{(\tilde{X}_i, \tilde{Y}_i, \tilde{Z}_i)\}_{i=1}^n$ can still converge to a Gaussian limit. This example illustrates that the limiting behavior of T and \tilde{T} can differ significantly, which is attributed to the instability of the estimators \hat{f} and \hat{g} .

The previous examples highlight the need for caution when converting a conditional independence test to a conditional two-sample test, and also justify our generic approach to converting a conditional independence test to a conditional two-sample test in Algorithm 1. The next section introduces another general framework for conditional two-sample testing based on density ratio estimation.

4 Approach via Density Ratio Estimation

In this section, we present our second framework, which transforms the problem of conditional two-sample testing into one that involves comparing marginal distributions via density ratio estimation. Concretely, we recall from (3) that the null hypothesis of equality of two conditional distributions holds if and only if $f_{Y|X}^{(1)} = f_{Y|X}$ where $f_{Y|X} = r_X \cdot f_{Y|X}^{(2)}$ and r_X is the density ratio defined in (4). A challenge when applying this approach is that we only have samples from $f_{Y|X}^{(2)}$, and not from $f_{Y|X}$, which makes it impossible to directly compare samples from $f_{Y|X}^{(1)}$ with those from $f_{Y|X}$. However, once the density ratio is known or accurately estimated, we can effectively correct the bias arising from the difference between $f_{Y|X}^{(2)}$ and $f_{Y|X}$ in various test statistics, frequently used for marginal two-sample testing. To facilitate our discussion, we first assume that the density ratio r_X is known and provide a detailed analysis on how to deal with the unknown case by focusing on a few cases.

At the core of our idea is importance weighting (Kimura and Hino, 2024, for a survey), a technique that assigns different levels of importance to data points to correct biases and prioritize relevant data. For instance, suppose we would like to estimate the expectation of X under the distribution P with density p , while we only observe data X_1, \dots, X_n from another distribution Q with density q . Then by re-weighting data points using the density ratio p/q , we can obtain an unbiased estimator of the expectation under P as

$$\frac{1}{n} \sum_{i=1}^n \frac{p(X_i)}{q(X_i)} X_i \quad \text{such that} \quad \mathbb{E}_Q \left[\frac{1}{n} \sum_{i=1}^n \frac{p(X_i)}{q(X_i)} X_i \right] = \mathbb{E}_P[X].$$

This idea can be applied to a range of marginal two-sample test statistics as we demonstrate below. Throughout this section, we use the shorthand $V^{(1)} := (X^{(1)}, Y^{(1)})$ and $V^{(2)} := (X^{(2)}, Y^{(2)})$ to simplify the notation.

Example 3. (Mean comparison) We start with a simple case of comparing the mean of transformed samples. Given a feature map $\psi : \mathcal{X} \times \mathcal{Y} \mapsto \mathbb{R}$, one can consider

$$\frac{1}{n_1} \sum_{i=1}^{n_1} \psi(V_i^{(1)}) - \frac{1}{n_2} \sum_{i=1}^{n_2} r_X(X_i^{(2)}) \psi(V_i^{(2)})$$

as a test statistic for the hypotheses in (1). The expectation of this statistic is equal to zero under the null hypothesis. Moreover, since the test statistic is simply a linear combination of independent random variables, it can be calibrated using the Gaussian approximation.

Example 4. (Rank sum statistic) Instead of comparing the mean, one can compare the stochastic order of two distributions using ranks. Specifically, given a feature map $\psi : \mathcal{X} \times \mathcal{Y} \mapsto \mathbb{R}$, a rank sum statistic based on the transformed samples can be computed as

$$\frac{1}{n_1 n_2} \sum_{i=1}^{n_1} \sum_{j=1}^{n_2} r_X(X_j^{(2)}) \mathbb{1}\{\psi(V_j^{(2)}) < \psi(V_i^{(1)})\}.$$

Under the null hypothesis, and assuming no ties among transformed samples, it can be seen that the expectation is equal to zero. As in [Hu and Lei \(2024\)](#), the test statistic can be shown to be asymptotically Gaussian using the asymptotic theory of U-statistics under conditions. Therefore, the critical value can be determined based on this Gaussian approximation. The power, however, changes depending on ψ . [Hu and Lei \(2024\)](#) takes ψ as an estimate of $f_{Y|X}^{(1)}(\cdot|\cdot)/f_{Y|X}^{(2)}(\cdot|\cdot)$.

Example 5. (Classifier-based approach) Let \mathcal{H} be a class of classifiers. Given a binary classifier $h : \mathcal{X} \times \mathcal{Y} \mapsto \{1, 2\}$ where $h \in \mathcal{H}$, $\ell : \mathbb{R} \times \{1, 2\} \mapsto \mathbb{R}$ is a loss function that measures the difference between the predicted value and the true output. The core idea behind classifier-based two-sample tests ([Lopez-Paz and Oquab, 2017](#); [Kim et al., 2021](#); [Hediger et al., 2022](#)) is that when the null hypothesis of equality of distributions is true, any classifier will return a random guess. On the other hand, when two distributions are significantly different, the accuracy of a reasonable classifier would be greater than chance level. Therefore, empirical classification accuracy can serve as an effective test statistic. However, since we do not observe a sample from f_{XY} but a sample from $f_{YX}^{(2)}$, we need to take this into consideration when we train a classifier. Specifically, we compute a classifier

$$\hat{h} = \arg \min_{h \in \mathcal{H}} \left\{ \frac{1}{n_1} \sum_{i=1}^{n_1} \ell(h(V_i^{(1)}), 2) + \frac{1}{n_2} \sum_{i=1}^{n_2} r_X(X_i^{(2)}) \ell(h(V_i^{(2)}), 1) \right\}, \quad (7)$$

and use the empirical classification accuracy of \hat{h} , again corrected by the density ratio, as our test statistic. When training and testing are performed on independent datasets, the asymptotic null distribution of the classification accuracy is approximately Gaussian ([Kim et al., 2021](#); [Hediger et al., 2022](#)); thereby the critical value can be determined based on this Gaussian approximation. We provide a detailed analysis of this approach in Section 4.1, and present numerical results in Section 5.

Example 6. (Kernel MMD) The last example is a kernel MMD statistic ([Gretton et al., 2012](#)). Given a kernel k , the population MMD compares the kernel mean embeddings of two distributions with density functions $f_{YX}^{(1)}$ and f_{YX} , respectively. In a kernel form, the squared MMD can be written as

$$\text{MMD}^2 = \mathbb{E}[k(V_1^{(1)}, V_2^{(1)})] + \mathbb{E}[r_X(X_1^{(2)})r_X(X_2^{(2)})k(V_1^{(2)}, V_2^{(2)})] - 2\mathbb{E}[r_X(X_2^{(2)})k(V_1^{(1)}, V_2^{(2)})],$$

where the bias is corrected via importance weighting. The squared MMD can be estimated as

$$\begin{aligned} & \frac{1}{n_1(n_1-1)} \sum_{1 \leq i \neq j \leq n_1} k(V_i^{(1)}, V_j^{(1)}) + \frac{1}{n_2(n_2-1)} \sum_{1 \leq i \neq j \leq n_2} r_X(X_i^{(2)})r_X(X_j^{(2)})k(V_i^{(2)}, V_j^{(2)}) \\ & - \frac{2}{n_1n_2} \sum_{i=1}^{n_1} \sum_{j=1}^{n_2} r_X(X_j^{(2)})k(V_i^{(1)}, V_j^{(2)}), \end{aligned}$$

which is an unbiased estimator of the population MMD. Unlike the other three test statistics mentioned earlier, this estimator converges to an infinite sum of weighted chi-squared distributions, whose weights are unknown in practice. This was not a major issue for marginal two-sample testing, where permutation tests could calibrate any test statistic in a nonparametric way. However, the standard permutation approach is no longer valid for conditional two-sample testing, which presents a challenge. Therefore, we focus on another estimator, a linear-time MMD statistic, in Section 4.2, which offers advantages in both the tractability of the asymptotic distribution and computational efficiency.

The preceding discussion assumes that the marginal density ratio r_X is known. As mentioned before, the success of this approach hinges on accurately estimating the density ratio r_X . To this end, one can draw upon a wide range of existing techniques in the literature for density ratio estimation (e.g., [Sugiyama et al., 2010, 2012](#)) to obtain a reliable testing result. Using the same dataset for both density ratio estimation and other parts of a statistic often results in plug-in bias. Hence, we recommend using auxiliary dataset

obtained through, e.g., sample splitting, to estimate the density ratio. We concretely illustrate this approach in Section 4.1 and Section 4.2, using a classifier-based test statistic and a linear-time MMD statistic.

While various existing tools for density ratio estimation offer flexibility, an inherent drawback of this approach is that the behavior of the test statistic could be erratic when the density ratio is irregular and potentially unbounded. This issue can be mitigated by clipping the density ratio estimate at a certain value or shrinking it to zero (e.g., Shimodaira, 2000). Nevertheless, when prior knowledge indicates that the density ratio behaves poorly or is difficult to estimate, the testing performance may degrade (see Section 5 for numerical results) and thus this approach should be used with caution.

4.1 Classifier-based Approach

This subsection illustrates a classifier-based test for conditional two-sample testing. To simplify the presentation, we assume that $n_1 = n_2 = 2n$ and split the dataset into two: $D_a := \{V_i^{(1)}\}_{i=1}^n \cup \{V_i^{(2)}\}_{i=1}^n$ and $D_b := \{V_i^{(1)}\}_{i=n+1}^{2n} \cup \{V_i^{(2)}\}_{i=n+1}^{2n}$. For some positive integer $m < n$, we further divide D_a as $D_a^* := \{V_i^{(1)}\}_{i=1}^m \cup \{V_i^{(2)}\}_{i=1}^m$ and $D_a^{**} := D_a \setminus D_a^{(1)}$, and let \hat{r}_X denote an estimator of r_X formed on D_a^{**} . Additionally, let \hat{h} be a classifier trained as in (7) based on D_b . Let us write $\hat{A}_{1,i} := \mathbb{1}\{\hat{h}(V_i^{(1)}) = 1\}$ and $\hat{A}_{2,i} := \hat{r}_X(X_i^{(2)})\mathbb{1}\{\hat{h}(V_i^{(2)}) = 2\}$ for $i \in \{1, \dots, m\}$, and define $\bar{A}_1 := m^{-1} \sum_{i=1}^m \hat{A}_{1,i}$ and $\bar{A}_2 := m^{-1} \sum_{i=1}^m \hat{A}_{2,i}$. The population-level classification accuracy of \hat{h} is $\mathbb{P}\{\hat{h}(V^{(1)}) = 1\}/2 + \mathbb{E}[r_X(X^{(2)})\mathbb{1}\{\hat{h}(V^{(2)}) = 2\}]/2$, which is 1/2 for any classifier \hat{h} under the null hypothesis. This observation leads to a classifier-based test statistic for conditional two-sample testing given as

$$\widehat{\text{Acc}} := \frac{\sqrt{m}(\bar{A}_1 + \bar{A}_2 - 1)}{\sqrt{\hat{\sigma}_1^2 + \hat{\sigma}_2^2}}, \quad (8)$$

where $\hat{\sigma}_1^2 := (m-1)^{-1} \sum_{i=1}^m (\hat{A}_{1,i} - \bar{A}_1)^2$ and $\hat{\sigma}_2^2 := (m-1)^{-1} \sum_{i=1}^m (\hat{A}_{2,i} - \bar{A}_2)^2$. To formally establish the limiting distribution of $\widehat{\text{Acc}}$, we consider the following assumptions.

Assumption 1. Let $m_n := m$ be an increasing sequence of positive integers with $\lim_{n \rightarrow \infty} m_n = \infty$. Consider a class of null distributions \mathcal{P}_0 such that

- (a) There are constants $c_1, c_2 \in (0, 1)$ such that $c_1 \leq \inf_{P \in \mathcal{P}_0} \mathbb{P}_P\{\hat{h}(V^{(1)}) = 1 | \hat{h}\} \leq \sup_{P \in \mathcal{P}_0} \mathbb{P}_P\{\hat{h}(V^{(1)}) = 1 | \hat{h}\} \leq c_2$ for all sufficiently large n . Moreover, assume that there exist constants $C, \delta > 0$ such that $\sup_{P \in \mathcal{P}_0} \mathbb{E}_P[\{\hat{r}_X(X^{(2)})\}^{2+\delta}] \leq C$ for all sufficiently large n .
- (b) For any $\epsilon > 0$, the density ratio estimator satisfies

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_0} \mathbb{P}_P(m \mathbb{E}_P[\{\hat{r}_X(X^{(2)}) - r_X(X^{(2)})\}^2 | \hat{r}_X] \geq \epsilon) = 0.$$

Assumption 1(a) is imposed to establish the (conditional) central limit theorem for the test statistic with the true density ratio, and it excludes a deterministic classifier, which would return the same prediction value regardless of inputs, under the null hypothesis. Assumption 1(b) ensures that the approximation error from \hat{r}_X is asymptotically negligible, which is similarly assumed in Hu and Lei (2024). In order to theoretically justify this condition, one needs to take m much smaller than the sample size used for training \hat{r}_X . However, our empirical results in Section 5 illustrate that $\widehat{\text{Acc}}$ approximates $N(0, 1)$ closely even under balanced splitting. Therefore, echoing Hu and Lei (2024), we suggest taking $m = \lfloor n/2 \rfloor$ in practice.

Under Assumption 1, the classifier-based test statistic in (8) converges to $N(0, 1)$ uniformly over \mathcal{P}_0 .

Theorem 3. For the class of null distributions \mathcal{P}_0 satisfying Assumption 1, $\widehat{\text{Acc}}$ in (8) converges to $N(0, 1)$:

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_0} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(\widehat{\text{Acc}} \leq t) - \Phi(t)| = 0.$$

The proof of Theorem 3 can be found in Appendix A.1. According to Theorem 3, the classifier-based test rejects the null hypothesis when $\widehat{\text{Acc}} > \Phi^{-1}(1 - \alpha)$, which has asymptotic validity over \mathcal{P}_0 satisfying Assumption 1. We next improve the efficiency of this procedure via K -fold cross-validation. To describe the procedure, we begin by considering K disjoint subsets of D_a , denoted as $D_{a,1}, D_{a,2}, \dots, D_{a,K}$, of equal size $m := \lfloor n/K \rfloor$ for simplicity. For $j \in \{1, \dots, K\}$, let $\bar{A}_{1,j} + \bar{A}_{2,j} - 1$ and $\widehat{\sigma}_{1,j}^2 + \widehat{\sigma}_{2,j}^2$ denote the quantities analogous to $\bar{A}_1 + \bar{A}_2 - 1$ and $\widehat{\sigma}_1^2 + \widehat{\sigma}_2^2$, respectively, by letting $D_a^* = D_{a,j}$ and $D_a^{**} = D_a \setminus D_{a,j}$. We then define the cross-validated classification accuracy statistic as

$$\widehat{\text{Acc}}_{\text{cv}} := \frac{1}{\sqrt{K}} \sum_{j=1}^K \frac{\sqrt{m}(\bar{A}_{1,j} + \bar{A}_{2,j} - 1)}{\sqrt{\widehat{\sigma}_{1,j}^2 + \widehat{\sigma}_{2,j}^2}}. \quad (9)$$

The next corollary proves that the cross-validated accuracy statistic is asymptotically normally distributed under the null hypothesis.

Corollary 1. *Consider the same setting as in Theorem 3. For any fixed $K \geq 2$, it holds that*

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_0} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(\widehat{\text{Acc}}_{\text{cv}} \leq t) - \Phi(t)| = 0.$$

According to Corollary 1, the test that rejects the null when $\widehat{\text{Acc}}_{\text{cv}} > \Phi^{-1}(1 - \alpha)$ maintains the asymptotic type I error under control. In terms of power, the cross-validated version is generally more powerful than the accuracy test without cross-validation as it uses the sample more efficiently. We numerically demonstrate this point in Section 5.

Many practical classifiers attempt to mimic the Bayes optimal classifier. For the balanced-sample setting, the Bayes optimal classifier is given as $h^*(x, y) := \mathbf{1}(f_{Y|X}^{(1)}(y, x)/\{f_{Y|X}^{(1)}(y, x) + f_{Y|X}^{(2)}(y, x)\} > 1/2)$ whose classification accuracy can be explicitly computed in terms of the total variation (TV) distance. Specifically, twice the classification accuracy can be computed as

$$\mathbb{P}\{h^*(V^{(1)}) = 1\} + \mathbb{E}[r_X(X^{(2)})\mathbf{1}\{h^*(V^{(2)}) = 2\}] = 1 + \text{TV}(f_{Y|X}^{(1)}, f_{Y|X}^{(2)}),$$

where $\text{TV}(f_{Y|X}^{(1)}, f_{Y|X}^{(2)})$ denotes the TV distance between two distributions with densities $f_{Y|X}^{(1)}$ and $f_{Y|X}^{(2)}$, respectively. Since the TV distance becomes zero if and only if two distributions are identical, our classifier-based test can be powerful against general alternatives when the classifier in use approximates the Bayes classifier.

The next subsection develops parallel results using a linear-time MMD statistic.

4.2 Linear-time MMD

In this subsection, we provide a detailed treatment of our second framework by focusing on a linear-time MMD statistic (Gretton et al., 2012, Lemma 14) with a kernel k . As in Section 4.1, we assume that $n_1 = n_2 = 2n$ and split the dataset into two: $D_a := \{V_i^{(1)}\}_{i=1}^n \cup \{V_i^{(2)}\}_{i=1}^n$ and $D_b := \{V_i^{(1)}\}_{i=n+1}^{2n} \cup \{V_i^{(2)}\}_{i=n+1}^{2n}$. Letting \widehat{r}_X be an estimator of r_X formed on D_b and $m = \lfloor n/2 \rfloor$, define

$$\begin{aligned} \widehat{S}_i &:= k(V_i^{(1)}, V_{i+m}^{(1)}) + \widehat{r}_X(X_i^{(2)})\widehat{r}_X(X_{i+m}^{(2)})k(V_i^{(2)}, V_{i+m}^{(2)}) \\ &\quad - \widehat{r}_X(X_i^{(2)})k(V_i^{(2)}, V_{i+m}^{(1)}) - \widehat{r}_X(X_{i+m}^{(2)})k(V_i^{(1)}, V_{i+m}^{(2)}). \end{aligned}$$

The test statistic that we analyze is a t -statistic based on $\widehat{S}_1, \dots, \widehat{S}_m$. Specifically, letting $\overline{S} := m^{-1} \sum_{i=1}^m \widehat{S}_i$ and $\widehat{\sigma}^2 := (m-1)^{-1} \sum_{i=1}^m (\widehat{S}_i - \overline{S})^2$, the (studentized) linear-time MMD statistic is given as

$$\widehat{\text{MMD}}_\ell^2 := \frac{\sqrt{m\overline{S}}}{\widehat{\sigma}}. \quad (10)$$

In order to establish the asymptotic normality of $\widehat{\text{MMD}}_\ell^2$, we make the following assumptions. Below, let S_i denote the quantity defined similarly as \widehat{S}_i by replacing \widehat{r}_X with the population counterpart r_X .

Assumption 2. Consider a class of null distributions \mathcal{P}_0 and assume that

- (a) There exist constants $c, C > 0$ such that $\inf_{P \in \mathcal{P}_0} \mathbb{E}_P[S_1^2] \geq c$ and $\sup_{P \in \mathcal{P}_0} \mathbb{E}_P[S_1^{2+\delta}] \leq C$ for some $\delta > 0$.
- (b) $\sup_{P \in \mathcal{P}_0} \mathbb{E}_P[\{r_X(X^{(2)})\}^2] < \infty$ and $\sup_{P \in \mathcal{P}_0} \mathbb{E}_P[\{\widehat{S}_X(X^{(2)}) - r_X(X^{(2)})\}^2] = o(m^{-1/2})$.
- (c) The kernel is uniformly bounded as $\|k\|_\infty \leq K$.

Assumption 2(a) is about the moment condition for the population counterpart of \widehat{S}_i . This assumption is required to apply the uniform central limit theorem. Assumption 2(b) is, on the other hand, required to prove that the difference between $\widehat{\text{MMD}}_\ell^2$ using $\{\widehat{S}_i\}_{i=1}^m$ and that using $\{S_i\}_{i=1}^m$ are asymptotically negligible. This is similar to Assumption 1(b) and the assumption made in [Hu and Lei \(2024\)](#), but this condition is considerably weaker in terms of convergence rate. Assumption 2(c) assumes that the kernel k is uniformly bounded. While this assumption is met for many practical kernels (e.g., Gaussian kernel), it can be relaxed by adopting more complex moment or convergence assumptions.

Having stated the assumptions, we now present the asymptotic normality of $\widehat{\text{MMD}}_\ell^2$ under the null hypothesis in (1).

Theorem 4. For the class of null distributions \mathcal{P}_0 satisfying Assumption 2, $\widehat{\text{MMD}}_\ell^2$ converges to $N(0, 1)$ as

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_0} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(\widehat{\text{MMD}}_\ell^2 \leq t) - \Phi(t)| = 0.$$

We defer the proof of Theorem 4 to Appendix A.3. Based on Theorem 4, the test that rejects the null when $\widehat{\text{MMD}}_\ell^2 > \Phi^{-1}(1 - \alpha)$ controls the size uniformly over the class of distributions that satisfy Assumption 2.

Similarly to the classification-based test, we can improve the efficiency of $\widehat{\text{MMD}}_\ell^2$ via K -fold cross-validation. To describe this process, we begin by partitioning the dataset of size $2n$ into K -folds, denoted as D_1, D_2, \dots, D_K , of equal size for simplicity. For $j \in \{1, \dots, K\}$, let \bar{S}_j and $\widehat{\sigma}_j^2$ denote the quantities similarly defined as \bar{S} and $\widehat{\sigma}^2$, respectively, by letting $D_a = D_j$ and $D_b = \cup_{i=1}^K D_i \setminus D_j$. We then define the cross-validated MMD statistic as

$$\widehat{\text{MMD}}_{\text{cv}}^2 := \frac{1}{K} \sum_{j=1}^K \frac{\sqrt{n}\bar{S}_j}{\widehat{\sigma}_j}. \quad (11)$$

The next corollary shows that $\widehat{\text{MMD}}_{\text{cv}}^2$ converges to $N(0, 1)$ under the same conditions for Theorem 4.

Corollary 2. Consider the same setting as in Theorem 4. Then for a fixed $K \geq 2$, it holds that

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_0} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(\widehat{\text{MMD}}_{\text{cv}}^2 \leq t) - \Phi(t)| = 0.$$

Again, the test that rejects the null when $\widehat{\text{MMD}}_{\text{cv}}^2 > \Phi^{-1}(1 - \alpha)$ is asymptotically level α based on Corollary 2. When the kernel k is a characteristic kernel (e.g., [Fukumizu et al., 2007](#)), the population MMD becomes equal to zero if and only if two distributions coincide. Thus as for the classifier-based tests in Section 4.1, the MMD-based tests can be powerful against general alternatives, provided that the density ratio r_X can be accurately estimated.

The next section illustrates the numerical performance of the proposed tests in comparison to existing methods.

5 Numerical Experiments

In this section, we evaluate the empirical performance of the proposed tests, alongside existing methods from the literature, across various scenarios. Within each scenario, we compare the conditional independence testing (CIT) approach described in Section 3 with the density ratio-based testing (DRT) approach described in Section 4. A brief overview of these methods is provided in Section 5.1, with additional implementation details available in Appendix C.

We empirically evaluate the type I error and power of these methods using both synthetic datasets in Section 5.2 and two real-world datasets in Section 5.3. Our simulation studies with synthetic datasets cover three distinct scenarios, each featuring both bounded and unbounded marginal density ratios. This setup allows us to explore a range of situations from relatively simple cases (with bounded density ratios) to more complex and challenging ones (with unbounded density ratios).

In all simulation studies, the dimension of the covariates X is fixed at $p = 10$. The simulation results are averaged over 500 repetitions with a significance level of $\alpha = 0.05$. For DRT methods, we employ a probabilistic classification approach using linear logistic regression for density ratio estimation (Sugiyama et al., 2010, Section 3). For CIT methods relying on regression estimation, such as the Generalized Covariance Measure (Shah and Peters, 2020) and the Projected Covariance Measure (Lundborg et al., 2022), we use a Random Forest model as the underlying regression method. On the other hand, for the WGSC (Williamson et al., 2023), we utilize XGBoost as it appears to perform best in our simulation scenarios.

It is crucial to note that the efficacy and validity of these methods depend on specific assumptions, which vary across different approaches. For example, DRT methods rely heavily on accurate density ratio estimation, while CIT methods depend on reliable estimation of conditional operators, such as regression functions. We aim to emphasize this distinction by providing a comprehensive evaluation of these approaches for conditional two-sample testing. This analysis sheds light on how each method performs under different conditions, guiding practitioners in choosing the most suitable approach for their specific applications.

5.1 Overview of Testing Methods

This subsection outlines the testing methods employed in our experiments. Further details on the implementation of these methods can be found in Appendix C. We denote the single-split classifier-based test in Section 4.1 as CLF and its cross-fit version as ${}^t\text{CLF}$. Both classifier-based test statistics are built on a specific form of classifiers detailed in Appendix C.1. Moreover, we denote the linear-time MMD in Section 4.2 as MMD- ℓ and its cross-fit version as ${}^t\text{MMD-}\ell$. Additional conditional two-sample testing methods included in our experiments are as follows:

CP: The conformal prediction (CP) test utilizes a conformity score to produce a weighted rank sum test statistic. This statistic is constructed by estimating both marginal and conditional density ratios, which can be approached using various density ratio estimation methods. For further details, please refer to Hu and Lei (2024).

DCP: The debiased conformal prediction (DCP) test refines the CP test by reducing bias through the use of Neyman orthogonality and using cross-fitting to improve efficiency. This enhancement guarantees asymptotic normality under certain conditions. Further technical details and theoretical guarantees are described in Chen and Lei (2024).

For CIT methods, we employ one kernel-based and three regression-based testing approaches. All of these CIT methods are implemented via Algorithm 1. We empirically observe in Appendix C.4 that the performance of the CIT methods remains largely consistent regardless of whether Algorithm 1 is applied or not, especially when the sample size is large. The following methods are included in our experiments:

RCIT: The randomized conditional independence test (RCIT) approximates the kernel conditional independence test by leveraging random Fourier features, allowing it to scale linearly with sample size. We use

the default options in its implementation, and further details are provided in [Strobl et al. \(2019\)](#).

GCM: The generalized covariance measure (GCM) by [Shah and Peters \(2020\)](#) utilizes the normalized covariance between residuals from regression models as a test statistic. This approach provides a flexible framework that can be adapted to various settings by selecting appropriate regression techniques.

PCM: The projected covariance measure (PCM) is a variation of the GCM applied to a transformed version of X . For our simulations, we follow Algorithm 1 from [Lundborg et al. \(2022\)](#). This method retains the general structure of the GCM while introducing a projection step, which enhances power, particularly when the conditional covariance is zero or near zero.

WGSC: This testing procedure proposes a general framework for nonparametric inference on interpretable, algorithm-agnostic variable importance. In our simulation, we follow the approach outlined in [Williamson et al. \(2023, Algorithm 3\)](#), which utilizes sample splitting and cross-fitting.

The code that reproduces all simulation results is available at: <https://github.com/suman-cha/Cond2ST>.

5.2 Synthetic Data Examples

We design three synthetic data scenarios to evaluate the performance of conditional two-sample testing methods under different conditions. Each scenario is implemented with both unbounded (U) and bounded (B) density ratios to assess how the difficulty of density ratio estimation affects the performance of each method. The marginal distributions of X remain consistent across scenarios with unbounded density ratios, and similarly, across scenarios with bounded density ratios.

For the unbounded case (U), we employ Gaussian distributions for marginal distributions of X . Specifically, for $j = 1$, samples are drawn from a standard Gaussian distribution, i.e., $x^{(1)} \sim N(0, I_p)$, where I_p is the identity matrix of dimension p . For $j = 2$, we introduce a covariate shift by sampling from a Gaussian distribution with mean vector $\mu = (1, 1, -1, -1, 0, \dots, 0)^\top$ and the same covariance structure, i.e., $x^{(2)} \sim N(\mu, I_p)$. In the bounded case (B), we truncate the support of both distributions to $[-0.5, 0.5]$ in each dimension, resulting in truncated Gaussian distributions, $x^{(1)} \sim TN(0, I_p)$ and $x^{(2)} \sim TN(\mu, I_p)$, where μ is the same as in the unbounded case.

Scenario 1: Linear Model with Mean Shift. Inspired by the work of [Hu and Lei \(2024\)](#), this scenario investigates the efficacy of testing methods in detecting the mean difference between two linear models. For each $j \in \{1, 2\}$, we set $y^{(j)} | x^{(j)} = \delta^{(j)} + x^{(j)\top} \beta + \epsilon^{(j)}$, where $\epsilon^{(j)}$ follows a t -distribution with 2 degrees of freedom. The regression coefficient β is set to $(1, -1, -1, 1, 0, \dots, 0)^\top$. Under the null hypothesis, we set $\delta^{(1)} = \delta^{(2)} = 0$, while for the alternative hypothesis, we introduce a mean shift by setting $\delta^{(1)} = 0$ and $\delta^{(2)} = 0.5$, thereby creating a difference in the two conditional distributions.

Scenario 2: High Variability in Conditional Distribution. We also investigate the effect of high variability in the conditional distribution, slightly modifying the example outlined in [Chatterjee et al. \(2024, Section 6.2\)](#). Under the null hypothesis, we model the conditional distributions as $y^{(j)} | x^{(j)} \sim N(x^{(j)\top} \beta^{(j)}, (\sigma^{(j)})^2)$, where $\beta^{(j)} = \mathbf{1}_p$ defined as a p -dimensional vector of ones and $(\sigma^{(j)})^2 = 10^2$ for both $j \in \{1, 2\}$. This implies that $\beta^{(1)}$ and $\beta^{(2)}$ are identical under the null hypothesis. For the alternative hypothesis, we modify $\beta^{(2)}$ to $(1, \dots, 1, 0)^\top$ and introduce heteroscedasticity by varying the variance for $j = 2$ as $(\sigma^{(2)})^2 = 10(1 + \exp(-\|x^{(2)} - 0.5\mathbf{1}_p\|_2^2/64))$.

Scenario 3: Post-Nonlinear Model. Our final scenario considers a post-nonlinear (PNL) model, which is widely used in causal predictive inference ([Zhang et al., 2017; Li et al., 2023](#)). It tests the capability of the methods to detect differences in non-linear relationships between variables. We model the conditional distributions as $y^{(j)} | x^{(j)} = f^{(j)}(x^{(j)\top} \mathbf{1}_p + 2\epsilon)$, where $\epsilon \sim N(0, 1)$ and $j \in \{1, 2\}$. Under the null hypothesis, we set $f^{(j)}(x) = \cos(x)$ for both $j \in \{1, 2\}$, while for the alternative hypothesis, $f^{(2)}(x)$ is randomly sampled from the set $\{x, x^2, x^3, \sin(x), \tanh(x)\}$.

Rejection Rate for Scenario 1

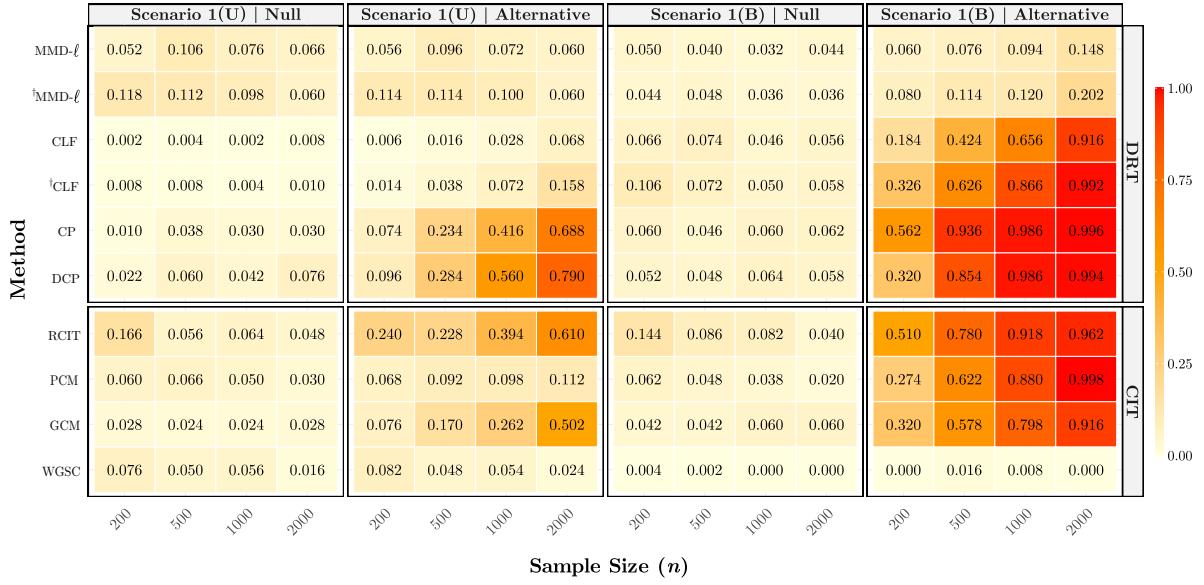


Figure 1: Rejection rates for Scenario 1 under null and alternative hypotheses, shown for both unbounded (U) and bounded (B) settings. Results are averaged over 500 repetitions with significance level $\alpha = 0.05$.

Rejection Rate for Scenario 2

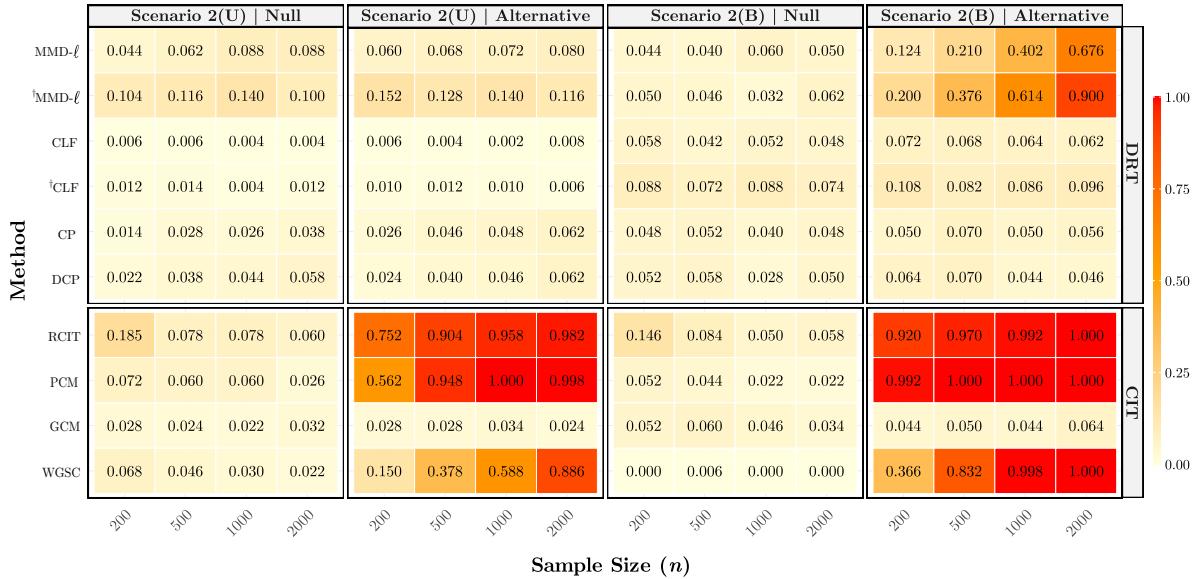


Figure 2: Rejection rates for Scenario 2 under null and alternative hypotheses, shown for both unbounded (U) and bounded (B) settings. Results are averaged over 500 repetitions with significance level $\alpha = 0.05$.

Our experimental results provide several key insights into the performance of conditional two-sample testing methods across diverse scenarios. A consistent pattern observed throughout all scenarios is the superior performance of DRT methods in bounded settings compared to unbounded settings. This improvement can be attributed to the relative ease of density ratio estimation when the density ratio is bounded, leading to more stable results. In contrast, CIT methods exhibit relatively consistent performance regardless of

Rejection Rate for Scenario 3

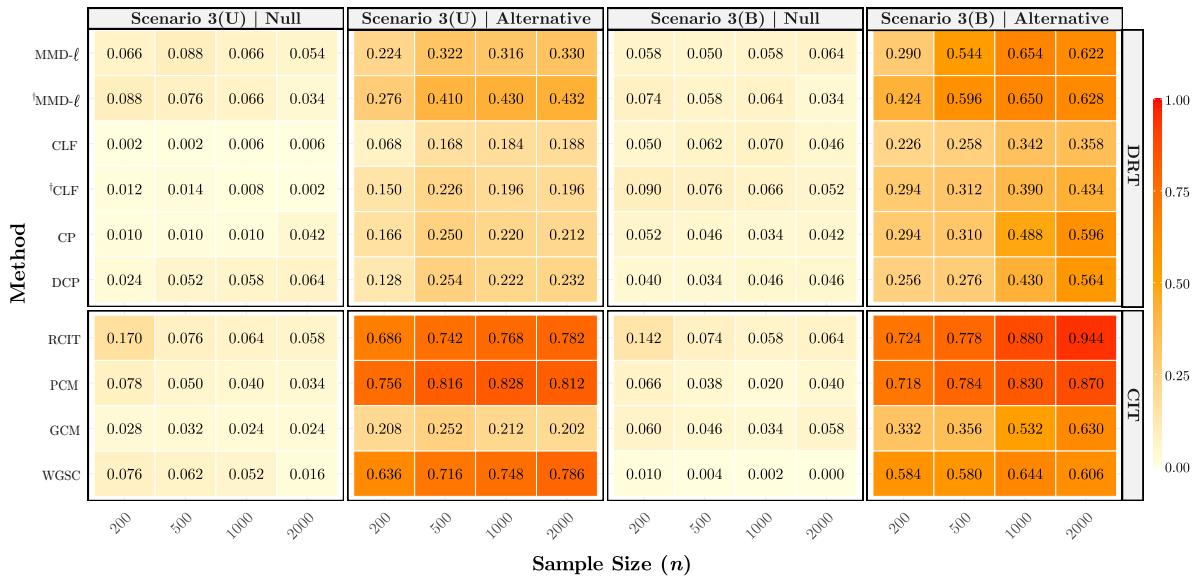


Figure 3: Rejection rates for Scenario 3 under null and alternative hypotheses, shown for both unbounded (U) and bounded (B) settings. Results are averaged over 500 repetitions with significance level $\alpha = 0.05$.

whether the density ratios are bounded or unbounded.

In Scenario 1, the classifier-based test shows the most significant improvement in performance when transitioning from unbounded to bounded cases. However, MMD- ℓ shows lower sensitivity in detecting mean shifts compared to other DRT methods. Among CIT methods, RCIT and GCM exhibit the best performance in this scenario. Scenario 2 highlights the strengths of MMD- ℓ , which only considers the marginal density ratio of X , in comparison to other DRT methods that account for conditional density ratios. MMD- ℓ shows a distinct advantage in this scenario. CIT methods also generally perform well under these conditions, showing their robustness to complex distributional changes. In Scenario 3, which tests the ability of methods to detect non-linear relationships, all DRT methods improve performance in the bounded case. Among CIT methods, there is no significant difference in performance, except for GCM and WGSC. GCM shows improved performance in the bounded case, while WGSC shows degraded performance.

These results underscore the critical role of accurate density ratio estimation in determining the performance of DRT methods. While CIT methods demonstrate consistent performance across both bounded and unbounded cases, suggesting their utility in a wide range of practical scenarios, they also have limitations. CIT methods, particularly regression-based approaches like GCM, PCM, and WGSC, can be sensitive to the choice of regression model, as we demonstrate in Appendix C. Notably, RCIT exhibits high type I error rates in all scenarios when sample sizes are relatively small, suggesting that caution is needed when applying RCIT to limited datasets. On the other hand, some methods show overly conservative behavior in certain scenarios. The cross-validated versions of DRT methods (${}^{\dagger}\text{MMD-}\ell$ and ${}^{\dagger}\text{CLF}$) consistently show power gains compared to their non-cross-validated counterparts as discussed in Section 4.1 and Section 4.2. Overall, our findings offer important insights into the strengths and limitations of different conditional two-sample testing methods.

5.3 Real Data Analysis

We further evaluate the performance of our proposed approaches on two real-world datasets: the diamonds dataset and the superconductivity dataset. Following Kim et al. (2023), we treat each dataset as a population

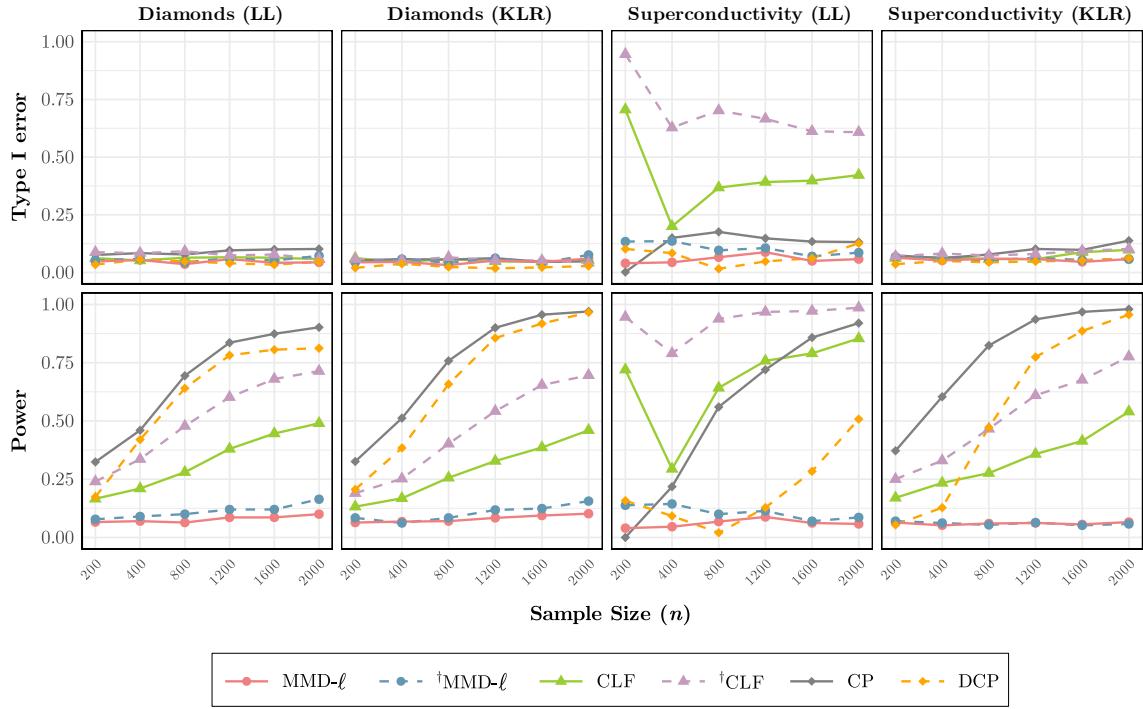


Figure 4: Performance comparison of DRT methods on diamonds and superconductivity datasets using LL and KLR for density ratio estimation. Rejection rates are averaged over 500 repetitions with $\alpha = 0.05$, under null (*top*) and alternative (*bottom*) hypotheses.

from which we draw samples, allowing for controlled experiments with known ground truth. Prior to analysis, we apply standard scaling to both X and Y variables. To introduce covariate shift, we implement biased sampling procedures. Specifically, we sample $X^{(1)}$ uniformly from the original feature space, while $X^{(2)}$ is sampled with probability proportional to $\exp(-x_1^2)$, where x_1 denotes the first feature of X . For the response variable Y , under the null hypothesis, we employ uniform sampling for both $Y^{(1)}$ and $Y^{(2)}$. Under the alternative hypothesis, $Y^{(1)}$ is sampled uniformly, while $Y^{(2)}$ is sampled with probability proportional to $\exp(-y)$, where y represents the values of Y in the dataset. Figure 4 illustrates the performance of the DRT methods on both datasets, using linear logistic (LL) and kernel logistic regression (KLR) for density ratio estimation.

Diamonds dataset. The diamonds dataset, available in the R package `ggplot2`, consists of 53,490 observations and 10 features, including price, carat, clarity and color. In our analysis, we set the price variable as Y , and use the 6 numerical variables (`carat`, `depth`, `table`, `x`, `y`, `z`) as X . As illustrated in Figure 4, most DRT methods exhibit good type I error control under both LL and KLR, with rejection rates generally close to the significance level α . Under the alternative hypothesis, we observe a clear trend of increasing power with sample size for all methods. Particularly, the cross-validated versions ($\dagger\text{MMD-}\ell$ and $\dagger\text{CLF}$) exhibit improved power, consistent with our observations in the synthetic data examples.

Superconductivity dataset. The superconductivity dataset, obtained from the UCI Machine Learning Repository and compiled by Hamidieh (2018), presents a more complex and high-dimensional challenge compared to the diamonds dataset. It comprises 81 features extracted from 21,263 superconductors, with the critical temperature at which the material transitions to a superconducting state serving as the response variable Y . The results reveal a significant contrast between density ratio estimation methods based on LL

and KLR. Under LL, several DRT methods, especially the classifier-based tests, struggle to control the type I error, with rejection rates far exceeding the significance level. Conversely, when using KLR for density ratio estimation, DRT methods show improved type I error control.

These empirical findings emphasize the importance of carefully considering the nature of the data and the choice of density ratio estimation techniques when applying DRT methods for conditional two-sample testing. The performance of different methods can vary significantly, indicating the need for careful method selection and, potentially, more advanced approaches when handling complex and high-dimensional data. While we focus on DRT methods in this section, experimental results for CIT methods are presented in Appendix C.3 for completeness.

6 Conclusion

In this paper, we shed new light on the relatively underexplored problem of conditional two-sample testing. We begin by characterizing the fundamental difficulty of the problem and highlighting the importance of assumptions to make it feasible. We then introduce two general frameworks: (1) converting conditional independence tests into conditional two-sample tests and (2) transforming the problem of comparing conditional distributions into marginal distributions based on density ratio estimation. Both approaches offer significant flexibility, allowing one to leverage well-developed tools to effectively tackle the problem.

Our work opens up several interesting directions for future work. One promising avenue is to extend our framework to conditional K -sample testing with a general $K \geq 2$. Such an extension would expand the applicability of our framework beyond the comparison of just two groups. This setting is related to conditional independence testing where Z is a categorical random variable taking values in $\{1, 2, \dots, K\}$. We expect our results established in Section 3 to serve as a cornerstone for this extension. Another direction worth exploring is establishing a framework for conditional two-sample testing based on resampling methods. One promising approach is the Sampling Importance Resampling (SIR) algorithm (Givens and Hoeting, 2012, Chapter 6.3), which allows us to obtain an approximate sample from the distribution with density $f_{Y|X}$. Future work can focus on methods that compare the sample from $P_{XY}^{(1)}$ with the approximate sample obtained from the SIR algorithm. Finally, one can explore other two-sample test statistics beyond those listed in Section 4. Of particular interest is the block-wise MMD statistic (Zaremba et al., 2013). This statistic has a tractable limiting distribution, while achieving lower variance with a slight increase in computational cost compared to the linear-time MMD statistic. We leave all these interesting topics for future work.

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Overview of Appendices. In Appendix A, we present the proofs omitted in the main paper. Appendix B gathers several lemmas that support these proofs. Finally, Appendix C provides implementation details of numerical experiments and additional simulation results.

A Proofs

Notation. For real sequences (a_n) and (b_n) , we say that $a_n \lesssim b_n$ if there exists a constant $C > 0$ such that $a_n \leq Cb_n$ for all n . Let $(X_{P,n})_{n \in \mathbb{N}, P \in \mathcal{P}}$ be a family of sequences of random variables determined by $P \in \mathcal{P}$. We say that $X_{P,n} = X_n = o_P(n^{-a})$ if for all $\epsilon > 0$,

$$\sup_{P \in \mathcal{P}} \mathbb{P}_P(n^a |X_{P,n}| > \epsilon) \rightarrow 0.$$

For a positive integer n , we use the shorthand $[n]$ to denote the set $\{1, \dots, n\}$.

A.1 Proof of Theorem 3

We analyze the numerator and the denominator of $\widehat{\text{Acc}}$, separately. In particular, we first show that the numerator converges to a Gaussian distribution and the denominator is ratio-consistent to the population-level standard deviation under Assumption 1.

Analysis of the numerator. Starting with the numerator, let us rewrite

$$\begin{aligned}\overline{A}_1 + \overline{A}_2 - 1 &= \frac{1}{m} \sum_{i=1}^m [\mathbb{1}\{\widehat{h}(V_i^{(1)}) = 1\} + r_X(X_i^{(2)})\mathbb{1}\{\widehat{h}(V_i^{(2)}) = 2\}] - 1 \\ &\quad + \frac{1}{m} \sum_{i=1}^m \{\widehat{r}_X(X_i^{(2)}) - r_X(X_i^{(2)})\}\mathbb{1}\{\widehat{h}(V_i^{(2)}) = 2\} \\ &= \frac{1}{m} \sum_{i=1}^m \underbrace{[\mathbb{1}\{\widehat{h}(V_i^{(1)}) = 1\} + r_X(X_i^{(2)})\mathbb{1}\{\widehat{h}(V_i^{(2)}) = 2\}] - 1}_{:=L_i(\widehat{h})} + o_{\mathcal{P}_0}(m^{-1/2}),\end{aligned}$$

where the last approximation holds since

$$\left| \frac{1}{m} \sum_{i=1}^m \{\widehat{r}_X(X_i^{(2)}) - r_X(X_i^{(2)})\}\mathbb{1}\{\widehat{h}(V_i^{(2)}) = 2\} \right| \leq \sqrt{\frac{1}{m} \sum_{i=1}^m \{\widehat{r}_X(X_i^{(2)}) - r_X(X_i^{(2)})\}^2},$$

and the upper bound is $o_{\mathcal{P}_0}(m^{-1/2})$ due to Assumption 1(b). Thus $\overline{A}_1 + \overline{A}_2 - 1$ is dominated by the average of $L_i(\widehat{h})$ values. Given this and Slutsky's theorem, it suffices to study the limiting distribution of the sample average of $L_i(\widehat{h})$. Indeed, under Assumption 1(a), the conditional central limit theorem (Lemma 3) yields that

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}_0} \sup_{t \in \mathbb{R}} \left| \mathbb{P}_P \left(\frac{\frac{1}{\sqrt{m}} \sum_{i=1}^m L_i(\widehat{h})}{\{\text{Var}[L(\widehat{h}) | \widehat{h}]\}^{1/2}} \leq t \right) - \Phi(t) \right| = 0.$$

Consistency of the variance estimate. We next show the ratio-consistency of the variance estimator. Observe that

$$\text{Var}[L(\widehat{h}) | \widehat{h}] = \underbrace{\text{Var}[\mathbb{1}\{\widehat{h}(V_i^{(1)}) = 1\} | \widehat{h}]}_{:=\sigma_1^2} + \underbrace{\text{Var}[r_X(X_i^{(2)})\mathbb{1}\{\widehat{h}(V_i^{(2)}) = 2\}]}_{:=\sigma_2^2},$$

and

$$\left| \frac{\widehat{\sigma}_1^2 + \widehat{\sigma}_2^2}{\sigma_1^2 + \sigma_2^2} - 1 \right| \leq \left| \frac{\widehat{\sigma}_1^2 - \sigma_1^2}{\sigma_1^2} \right| + \left| \frac{\widehat{\sigma}_2^2 - \sigma_2^2}{\sigma_2^2} \right|.$$

Therefore, in order to show the ratio consistency of $\widehat{\sigma}_1^2 + \widehat{\sigma}_2^2$, it suffices to show the ratio consistency of $\widehat{\sigma}_1^2$ and $\widehat{\sigma}_2^2$, individually. To this end, we use conditional Chebyshev's inequality and show

$$\sup_{P \in \mathcal{P}_0} \mathbb{P}_P(|\widehat{\sigma}_1^2/\sigma_1^2 - 1| \geq t | \widehat{h}) \leq \frac{1}{t^2} \sup_{P \in \mathcal{P}_0} \text{Var}_P(\widehat{\sigma}_1^2/\sigma_1^2 | \widehat{h}) \leq \frac{1}{t^2 c_1 (1 - c_2)m},$$

for sufficiently large n and for all $t > 0$, under Assumption 1(a). Hence $\widehat{\sigma}_1^2/\sigma_1^2$ converges to one in probability uniformly over \mathcal{P}_0 . On the other hand, letting $A_{2,i} := \widehat{r}_X(X_i^{(2)})\mathbb{1}\{\widehat{h}(V_i^{(2)}) = 2\}$, we have

$$\begin{aligned}\widehat{\sigma}_2^2 &= \frac{1}{m-1} \sum_{i=1}^m \left[A_{2,i} - \frac{1}{m} \sum_{j=1}^m A_{2,j} \right]^2 \\ &\quad + \frac{1}{m-1} \sum_{i=1}^m \left[(\widehat{A}_{2,i} - A_{2,i}) - \frac{1}{m} \sum_{j=1}^m (\widehat{A}_{2,j} - A_{2,j}) \right]^2 \\ &\quad + \frac{2}{m-1} \sum_{i=1}^m \left[A_{2,i} - \frac{1}{m} \sum_{j=1}^m A_{2,j} \right] \cdot \left[(\widehat{A}_{2,i} - A_{2,i}) - \frac{1}{m} \sum_{j=1}^m (\widehat{A}_{2,j} - A_{2,j}) \right] \\ &:= (\text{I}) + (\text{II}) + (\text{III}).\end{aligned}$$

Similarly as before, the term $(\mathbb{I})/\sigma_2^2$ converges to one in probability uniformly over \mathcal{P}_0 under Assumption 1(a). It can be further shown that the term $(\mathbb{III})/\sigma_2^2$ is $o_{\mathcal{P}_0}(1)$ by Markov's inequality combined with Assumption 1(b). Lastly, the term (\mathbb{IIII}) satisfies $(\mathbb{IIII}) \leq \sqrt{(\mathbb{I})} \times (\mathbb{II})$, which is again $o_{\mathcal{P}_0}(1)$. Therefore, $\widehat{\sigma}_2^2/\sigma_2^2$ converges to one in probability uniformly over \mathcal{P}_0 . This further proves that $\sqrt{(\sigma_1^2 + \sigma_2^2)/(\widehat{\sigma}_1^2 + \widehat{\sigma}_2^2)} = 1 + o_{\mathcal{P}_0}(1)$ by Lundborg et al. (2022, Lemma S7).

Putting all pieces together with Lemma 4(b) proves the claim.

A.2 Proof of Corollary 1

For each $j \in \{1, \dots, K\}$, the proof of Theorem 3 shows that $\sqrt{\sigma_1^2 + \sigma_2^2}/\sqrt{\widehat{\sigma}_{1,j}^2 + \widehat{\sigma}_{2,j}^2} = 1 + o_{\mathcal{P}_0}(1)$. Thus, by Lemma 4(b), it is enough to show the asymptotic normality of

$$\frac{1}{\sqrt{K}} \sum_{j=1}^K \frac{\sqrt{m}(\overline{A}_{1,j} + \overline{A}_{2,j} - 1)}{\sqrt{\sigma_1^2 + \sigma_2^2}}.$$

Without loss of generality, denote the sample indices of D_1, D_2, \dots, D_K as

$$I_1 = \{1, \dots, m\}, I_2 = \{m+1, \dots, 2m\}, \dots, I_K = \{m(K-1)+1, \dots, mK\}.$$

Then the proof of Theorem 3 establishes that

$$\frac{1}{\sqrt{K}} \sum_{j=1}^K \frac{\sqrt{m}(\overline{A}_{1,j} + \overline{A}_{2,j} - 1)}{\sqrt{\sigma_1^2 + \sigma_2^2}} = \sqrt{\frac{1}{mK}} \sum_{j=1}^K \underbrace{\left\{ \sum_{i \in I_j} \frac{A_{1,i} + A_{2,i} - 1}{\sqrt{\sigma_1^2 + \sigma_2^2}} \right\}}_{:= B_j} + o_{\mathcal{P}_0}(1),$$

where $A_{1,i} + A_{2,i} - 1 := \mathbb{1}\{\widehat{h}(V_i^{(1)}) = 1\} + r_X(X_i^{(2)})\mathbb{1}\{\widehat{h}(V_i^{(2)}) = 2\} - 1$. Note that B_1, \dots, B_K are mutually independent conditional on \widehat{h} . As in the proof of Theorem 3, we apply the conditional central limit theorem (Lemma 3) to the average of B_1, \dots, B_K conditional on \widehat{h} under Assumption 1, which completes the proof of Corollary 1.

A.3 Proof of Theorem 4

The proof consists of two parts as in the proof of Theorem 3. In the first part, we investigate the numerator of $\widehat{\text{MMD}}_\ell^2$, i.e., \overline{S} , whereas in the second part, we show the consistency of the denominator to the population variance under Assumption 2. The proof is then completed by applying Lemma 4(b).

Analysis of the numerator. Using the fact that a kernel can be expressed as an inner product of feature maps, $k(x, y) = \langle \psi(x), \psi(y) \rangle$, we can rewrite \overline{S} as

$$\overline{S} = \frac{1}{m} \sum_{i=1}^m \langle \psi(V_i^{(1)}) - \widehat{r}_X(X_i^{(2)})\psi(V_i^{(2)}), \psi(V_{i+m}^{(1)}) - \widehat{r}_X(X_{i+m}^{(2)})\psi(V_{i+m}^{(2)}) \rangle.$$

By adding and subtracting $r_X(X_i^{(2)})\psi(V_i^{(2)})$ and $r_X(X_{i+m}^{(2)})\psi(V_{i+m}^{(2)})$, \overline{S} can be written as the sum of the

four terms given as:

$$\begin{aligned}
(\mathbb{I}) &:= \frac{1}{m} \sum_{i=1}^m \underbrace{\langle \psi(V_i^{(1)}) - r_X(X_i^{(2)})\psi(V_i^{(2)}), \psi(V_{i+m}^{(1)}) - r_X(X_{i+m}^{(2)})\psi(V_{i+m}^{(2)}) \rangle}_{:= S_i}, \\
(\mathbb{II}) &:= \frac{1}{m} \sum_{i=1}^m \underbrace{\langle \psi(V_i^{(1)}) - r_X(X_i^{(2)})\psi(V_i^{(2)}), \psi(V_{i+m}^{(2)}) \rangle \cdot \{\widehat{r}_X(X_{i+m}^{(2)}) - r_X(X_{i+m}^{(2)})\}}_{:= \widehat{S}_{i,a}}, \\
(\mathbb{III}) &:= \frac{1}{m} \sum_{i=1}^m \underbrace{\langle \psi(V_i^{(2)}), \psi(V_{i+m}^{(1)}) - r_X(X_{i+m}^{(2)})\psi(V_{i+m}^{(2)}) \rangle \cdot \{\widehat{r}_X(X_i^{(2)}) - r_X(X_i^{(2)})\}}_{:= \widehat{S}_{i,b}}, \\
(\mathbb{IV}) &:= \frac{1}{m} \sum_{i=1}^m \underbrace{\langle \psi(V_i^{(2)}), \psi(V_{i+m}^{(2)}) \rangle \{\widehat{r}_X(X_i^{(2)}) - r_X(X_i^{(2)})\} \cdot \{\widehat{r}_X(X_{i+m}^{(2)}) - r_X(X_{i+m}^{(2)})\}}_{:= \widehat{S}_{i,c}}.
\end{aligned}$$

The first term (\mathbb{I}) does not involve an estimate of the density ratio and will be asymptotically Gaussian since it is the sum of i.i.d. random variables under the null hypothesis. The other terms (\mathbb{II}) , (\mathbb{III}) , and (\mathbb{IV}) are asymptotically negligible under the conditions of the theorem. Hence \bar{S} will be dominated by (\mathbb{I}) . Let us analyze each term separately.

1. **Term (\mathbb{I}) .** Define $\text{Var}_P[S_1] = \sigma_P^2$. Then under Assumption 2(a), Lemma 2 yields

$$\sup_{P \in \mathcal{P}_0} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(\sqrt{m}\sigma_P^{-1}(\mathbb{I}) \leq t) - \Phi(t)| \rightarrow 0.$$

2. **Terms (\mathbb{II}) and (\mathbb{III}) .** We only analyze the term (\mathbb{II}) since (\mathbb{III}) can be handled in exactly the same way by symmetry. Under the null hypothesis, by the law of total expectation, it can be seen that the expectation of the summands of (\mathbb{II}) is equal to zero:

$$\mathbb{E}[\langle \psi(V_i^{(1)}) - r_X(X_i^{(2)})\psi(V_i^{(2)}), \psi(V_{i+m}^{(2)}) \rangle \cdot \{\widehat{r}_X(X_{i+m}^{(2)}) - r_X(X_{i+m}^{(2)})\}] = 0,$$

which leads to $\mathbb{E}[(\mathbb{II})] = 0$. On the other hand, the conditional second moment (or the conditional variance) of (\mathbb{II}) given D_b satisfies

$$\begin{aligned}
\mathbb{E}[(\mathbb{II})^2 | D_b] &= \frac{1}{m} \mathbb{E}[\langle \psi(V_1^{(1)}) - r_X(X_1^{(2)})\psi(V_1^{(2)}), \psi(V_{1+m}^{(2)}) \rangle^2 \cdot \{\widehat{r}_X(X_{1+m}^{(2)}) - r_X(X_{1+m}^{(2)})\}^2 | D_b] \\
&\lesssim \frac{K^2}{m} (1 + \mathbb{E}[r_X(X_1^{(2)})^2]) \cdot \mathbb{E}[\{\widehat{r}_X(X_{1+m}^{(2)}) - r_X(X_{1+m}^{(2)})\}^2 | D_b],
\end{aligned}$$

where we use the fact that $\langle \psi(x), \psi(y) \rangle = k(x, y)$, whose ℓ_∞ norm is uniformly bounded by the constant K . Therefore, under the condition that

$$\sup_{P \in \mathcal{P}_0} \mathbb{E}_P[\{\widehat{r}_X(X_{1+m}^{(2)}) - r_X(X_{1+m}^{(2)})\}^2] = o(m^{-1/2}) \quad \text{and} \quad \sup_{P \in \mathcal{P}_0} \mathbb{E}_P[r_X(X_1^{(2)})^2] < \infty,$$

Chebyshev's inequality yields $(\mathbb{II}) = o_{\mathcal{P}_0}(m^{-1/2})$ and similarly $(\mathbb{III}) = o_{\mathcal{P}_0}(m^{-1/2})$.

3. **Term (\mathbb{IV}) .** The fourth term (\mathbb{IV}) can be written as

$$(\mathbb{IV}) = \frac{1}{m} \sum_{i=1}^m k(V_i^{(2)}, V_{i+m}^{(2)}) \{\widehat{r}_X(X_i^{(2)}) - r_X(X_i^{(2)})\} \cdot \{\widehat{r}_X(X_{i+m}^{(2)}) - r_X(X_{i+m}^{(2)})\}.$$

Since the kernel is uniformly bounded and by the Cauchy–Schwarz inequality, we have

$$\begin{aligned} |(\text{IV})| &\leq K \left(\frac{1}{m} \sum_{i=1}^m \{\hat{r}_X(X_i^{(2)}) - r_X(X_i^{(2)})\}^2 \right)^{1/2} \cdot \left(\frac{1}{m} \sum_{i=1}^m \{\hat{r}_X(X_{i+m}^{(2)}) - r_X(X_{i+m}^{(2)})\}^2 \right)^{1/2} \\ &\stackrel{(*)}{=} o_{\mathcal{P}_0}(m^{-1/4})o_{\mathcal{P}_0}(m^{-1/4}) = o_{\mathcal{P}_0}(m^{-1/2}), \end{aligned}$$

which follows by Markov's inequality along with the condition that

$$\sup_{P \in \mathcal{P}_0} \mathbb{E}_P [\{\hat{r}_X(X^{(2)}) - r_X(X^{(2)})\}^2] = o(m^{-1/2}),$$

and step \$(*)\$ holds by [Lundborg et al. \(2022, Lemma S5\)](#). Therefore it holds that \$(\text{IV}) = o_{\mathcal{P}_0}(m^{-1/2})\$.

Now combining the results establishes that

$$\sup_{P \in \mathcal{P}_0} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(\sqrt{m}\sigma_P^{-1}\bar{S} \leq t) - \Phi(t)| \rightarrow 0.$$

Consistency of the variance estimate. Denoting

$$\hat{\sigma}_P^2 := \frac{1}{m-1} \sum_{i=1}^m (\hat{S}_i - \bar{S})^2,$$

we would like to show that \$\hat{\sigma}_P^2/\sigma_P^2\$ converges to one in probability, which further implies \$\sigma_P/\hat{\sigma}_P = 1 + o_{\mathcal{P}_0}(1)\$ by [Lundborg et al. \(2022, Lemma S7\)](#). Since the test statistic \$\widehat{\text{MMD}}_\ell^2\$ is scale-invariant, we may assume that \$\sigma_P^2 = 1\$ without loss of generality. Moreover, the preceding analysis ensures that \$\bar{S} = o_{\mathcal{P}_0}(1)\$. Therefore we only need to show \$\frac{1}{m} \sum_{i=1}^m \hat{S}_i^2\$ converges to one in probability. To this end, observe

$$\begin{aligned} \left| \frac{1}{m} \sum_{i=1}^m \hat{S}_i^2 - 1 \right| &= \left| \frac{1}{m} \sum_{i=1}^m (S_i + \hat{S}_{i,a} + \hat{S}_{i,b} + \hat{S}_{i,c})^2 - 1 \right| \\ &\leq \left| \frac{1}{m} \sum_{i=1}^m S_i^2 - 1 \right| + \left| \frac{1}{m} \sum_{i=1}^m (\hat{S}_{i,a} + \hat{S}_{i,b} + \hat{S}_{i,c})^2 \right| + 2 \left| \frac{1}{m} \sum_{i=1}^m S_i (\hat{S}_{i,a} + \hat{S}_{i,b} + \hat{S}_{i,c}) \right| \\ &\leq \left| \frac{1}{m} \sum_{i=1}^m S_i^2 - 1 \right| + \left| \frac{1}{m} \sum_{i=1}^m (\hat{S}_{i,a} + \hat{S}_{i,b} + \hat{S}_{i,c})^2 \right| \\ &\quad + 2 \sqrt{\frac{1}{m} \sum_{i=1}^m S_i^2} \sqrt{\frac{1}{m} \sum_{i=1}^m (\hat{S}_{i,a} + \hat{S}_{i,b} + \hat{S}_{i,c})^2}, \end{aligned}$$

where the last inequality follows by the Cauchy–Schwarz inequality. By the law of large numbers, \$\frac{1}{m} \sum_{i=1}^m S_i^2\$ converges to one in probability. Thus the proof amounts to showing that \$\frac{1}{m} \sum_{i=1}^m (\hat{S}_{i,a} + \hat{S}_{i,b} + \hat{S}_{i,c})^2 = o_{\mathcal{P}_0}(1)\$, which is implied by

$$\frac{1}{m} \sum_{i=1}^m \hat{S}_{i,a}^2 = o_{\mathcal{P}_0}(1), \quad \frac{1}{m} \sum_{i=1}^m \hat{S}_{i,b}^2 = o_{\mathcal{P}_0}(1) \quad \text{and} \quad \frac{1}{m} \sum_{i=1}^m \hat{S}_{i,c}^2 = o_{\mathcal{P}_0}(1).$$

This can be done as the way how (III), (III), and (IV) are handled earlier along with Markov's inequality. This completes the proof.

A.4 Proof of Corollary 2

For each $j \in \{1, \dots, K\}$, the proof of Theorem 4 shows that $\sigma_P/\hat{\sigma}_j = 1 + o_{\mathcal{P}_0}(1)$. Thus, by Lemma 4(b), it is enough to show the asymptotic normality of

$$\frac{1}{K} \sum_{j=1}^K \frac{\sqrt{n} \bar{S}_j}{\sigma}.$$

Without loss of generality, denote the sample indices of D_1, D_2, \dots, D_K as

$$I_1 = \left\{ 1, \dots, \frac{2n}{K} \right\}, \quad I_2 = \left\{ \frac{2n}{K} + 1, \dots, \frac{4n}{K} \right\}, \dots, I_K = \left\{ 2n - \frac{2n}{K} + 1, \dots, 2n \right\},$$

and let the first n/K elements of I_j as I'_j (e.g., $I'_1 = \{1, \dots, n/K\}$). Then with $m' = n/K$, the proof of Theorem 4 establishes that

$$\sum_{j=1}^K \bar{S}_j = \frac{K}{n} \sum_{j=1}^K \sum_{i \in I'_j} S_i + o_{\mathcal{P}_0}(n^{-1/2}),$$

where

$$S_i = \langle \psi(V_i^{(1)}) - r_X(X_i^{(2)})\psi(V_i^{(2)}), \psi(V_{i+m'}^{(1)}) - r_X(X_{i+m'}^{(2)})\psi(V_{i+m'}^{(2)}) \rangle.$$

Notably, $\sum_{i \in I'_1} S_i, \dots, \sum_{i \in I'_K} S_i$ are mutually independent. Hence

$$\text{Var}\left(\frac{K}{n} \sum_{j=1}^K \sum_{i \in I'_j} \sigma_P^{-1} S_i\right) = \frac{K^2}{n^2} \cdot K \cdot \frac{n}{K} = \frac{K^2}{n}.$$

By the central limit theorem (Lemma 2),

$$\frac{1}{K} \sum_{j=1}^K \sum_{i \in I'_j} \frac{\sqrt{n} S_i}{\sigma_P} = \frac{1}{K} \sum_{j=1}^K \frac{\sqrt{n} \bar{S}_j}{\sigma} + o_{\mathcal{P}_0}(1)$$

converges to $N(0, 1)$ as desired.

A.5 Proof of Example 1

In this section, we aim to present a detailed analysis of the asymptotic equivalence between the GCM statistic T and its counterpart \tilde{T} constructed using $\{(\tilde{X}_i, \tilde{Y}_i, \tilde{Z}_i)\}_{i=1}^n$. Throughout this section, we assume that Y has a finite second moment and that $Y \perp\!\!\!\perp Z | X$, i.e., the null hypothesis holds. Let

$$T = \frac{\frac{1}{\sqrt{n}} \sum_{i=1}^n R_i}{\left\{ \frac{1}{n} \sum_{i=1}^n R_i^2 - \left(\frac{1}{n} \sum_{r=1}^n R_r \right)^2 \right\}^{1/2}} := \frac{\nu_R}{\hat{\sigma}_R},$$

and $\tilde{T} := \nu_{\tilde{R}}/\hat{\sigma}_{\tilde{R}}$. Let $\sigma_R^2 > 0$ denote the variance of $\{\tilde{Y} - f(\tilde{X})\}\{\tilde{Z} - g(\tilde{X})\}$ where $(\tilde{X}, \tilde{Y}, \tilde{Z})$ is a random draw from the joint distribution P_{XYZ} . We begin with an upper bound for $|T - \tilde{T}|$:

$$\begin{aligned} \left| \frac{\nu_R}{\hat{\sigma}_R} - \frac{\nu_{\tilde{R}}}{\hat{\sigma}_{\tilde{R}}} \right| &\leq \left| \frac{\nu_R}{\hat{\sigma}_R} - \frac{\nu_{\tilde{R}}}{\hat{\sigma}_R} \right| + \left| \frac{\nu_{\tilde{R}}}{\hat{\sigma}_R} - \frac{\nu_{\tilde{R}}}{\hat{\sigma}_{\tilde{R}}} \right| \\ &\leq \frac{1}{\hat{\sigma}_R} |\nu_R - \nu_{\tilde{R}}| + \frac{|\nu_{\tilde{R}}|}{(\hat{\sigma}_R + \hat{\sigma}_{\tilde{R}})\hat{\sigma}_R \hat{\sigma}_{\tilde{R}}} |\hat{\sigma}_R^2 - \hat{\sigma}_{\tilde{R}}^2|, \end{aligned}$$

from which the proof boils down to showing the convergence of the following four terms to zero in probability:
(a) $\nu_R - \nu_{\tilde{R}}$, (b) $\hat{\sigma}_R^2 - \hat{\sigma}_{\tilde{R}}^2$, (c) $\hat{\sigma}_R^2 - \sigma_R^2$ and (d) $\hat{\sigma}_{\tilde{R}}^2 - \sigma_{\tilde{R}}^2$. Under these convergence results, the asymptotic

equivalence follows by the continuous mapping theorem along with the fact that $\nu_{\tilde{R}}$ is stochastically bounded by the central limit theorem. In what follows, we establish convergence of (a), (b), (c), and (d) to zero in probability in order.

1. Term (a): $\nu_R - \nu_{\tilde{R}}$. Starting with the term (a), the difference between ν_R and $\nu_{\tilde{R}}$ can be written as

$$\frac{1}{\sqrt{n}} \sum_{i=1}^n R_i - \frac{1}{\sqrt{n}} \sum_{i=1}^n \tilde{R}_i = \underbrace{\frac{1}{\sqrt{n}} \sum_{i=n_1+1}^{\bar{n}_1} (R_i - \tilde{R}_i) \cdot \mathbb{1}(\bar{n}_1 > n_1)}_{:=\Delta_1} + \underbrace{\frac{1}{\sqrt{n}} \sum_{i=\bar{n}_1+1}^{n_1} (R_i - \tilde{R}_i) \cdot \mathbb{1}(\bar{n}_1 \leq n_1)}_{:=\Delta_2}.$$

Remark that Z_i is a fixed constant for a given index i , which allows us to show that $\mathbb{E}[R_i] = 0$ for any $i \in [n]$. For example, when $i = 1$, Z_1 equals 1 (since $X_1 = X_1^{(1)}$) and thus the law of total expectation yields

$$\mathbb{E}[R_1] = \mathbb{E}[\{1 - g(X_1)\}\{Y_1 - f(X_1)\}] = \mathbb{E}[\{1 - g(X_1)\}\mathbb{E}[\{Y_1 - f(X_1)\} | X_1]] = 0,$$

where we recall $f(X_1) = \mathbb{E}[Y_1 | X_1]$. It also follows that $\mathbb{E}[\tilde{R}_i] = 0$ for any $i \in [n]$ under the null hypothesis. This together with the law of total expectation shows that

$$\mathbb{E}[\Delta_1] = \mathbb{E}\left[\frac{1}{\sqrt{n}} \sum_{i=n_1+1}^{\bar{n}_1} \mathbb{1}(\bar{n}_1 > n_1) \mathbb{E}[R_i - \tilde{R}_i | \bar{n}_1]\right] = 0,$$

and similarly $\mathbb{E}[\Delta_2] = 0$. Thus, we have $\mathbb{E}[\nu_R - \nu_{\tilde{R}}] = 0$.

Now consider the variance of $\nu_R - \nu_{\tilde{R}}$. Since $\mathbb{E}[\nu_R - \nu_{\tilde{R}}] = 0$ and $\mathbb{E}[\Delta_1 \Delta_2] = 0$, we have

$$\text{Var}[\nu_R - \nu_{\tilde{R}}] = \text{Var}[\Delta_1] + \text{Var}[\Delta_2].$$

For $\text{Var}[\Delta_1]$, we have

$$\begin{aligned} \text{Var}[\Delta_1] &= \mathbb{E}[\text{Var}\{\Delta_1 | \bar{n}_1\}] + \text{Var}\left[\underbrace{\mathbb{E}\{\Delta_1 | \bar{n}_1\}}_{=0}\right] \\ &= \mathbb{E}\left[\text{Var}\left\{\frac{1}{\sqrt{n}} \sum_{i=n_1+1}^{\bar{n}_1} (R_i - \tilde{R}_i) \cdot \mathbb{1}(\bar{n}_1 > n_1) \mid \bar{n}_1\right\}\right] \\ &= \mathbb{E}\left[\frac{1}{n} \sum_{i=n_1+1}^{\bar{n}_1} \mathbb{1}(\bar{n}_1 > n_1) \mathbb{E}\{(R_i - \tilde{R}_i)^2 \mid \bar{n}_1\}\right] \\ &\stackrel{(i)}{\leq} 2\mathbb{E}\left[\frac{\bar{n}_1 - n_1}{n} \cdot \frac{1}{\bar{n}_1 - n_1} \sum_{i=n_1+1}^{\bar{n}_1} \mathbb{E}(R_i^2 + \tilde{R}_i^2 \mid \bar{n}_1)\right] \stackrel{(ii)}{\leq} \frac{4}{n} \mathbb{E}[|\bar{n}_1 - n_1|] \text{Var}(Y_1) \stackrel{(iii)}{\leq} \frac{2}{\sqrt{n}} \text{Var}(Y_1), \end{aligned}$$

where (i) follows from the inequality $(x - y)^2 \leq 2x^2 + 2y^2$ and (ii) uses the law of total variance along with the fact that $R_i^2 \leq \{Y_i - f(X_i)\}^2$ and $\tilde{R}_i^2 \leq \{\tilde{Y}_i - f(\tilde{X}_i)\}^2$ since $Z_i, \tilde{Z}_i \in \{1, 2\}$. For the last inequality (iii), we use $\mathbb{E}[|\bar{n}_1 - n_1|] \leq \sqrt{n}/2$. The same bound holds for $\text{Var}[\Delta_2]$ and thus

$$\text{Var}[\nu_R - \nu_{\tilde{R}}] \leq \frac{4}{\sqrt{n}} \text{Var}(Y_1).$$

Combining the results with Chebyshev's inequality now shows that $\nu_R - \nu_{\tilde{R}}$ converges to zero in probability.

2. Term (b): $\hat{\sigma}_R^2 - \hat{\sigma}_{\tilde{R}}^2$. We next aim to show that

$$\hat{\sigma}_R^2 - \hat{\sigma}_{\tilde{R}}^2 = \left\{ \frac{1}{n} \sum_{i=1}^n R_i^2 - \left(\frac{1}{n} \sum_{r=1}^n R_r \right)^2 \right\} - \left\{ \frac{1}{n} \sum_{i=1}^n \tilde{R}_i^2 - \left(\frac{1}{n} \sum_{r=1}^n \tilde{R}_r \right)^2 \right\}$$

converges to zero in probability. We decompose this into two terms

$$(\mathbb{I}) := \frac{1}{n} \sum_{i=1}^n R_i^2 - \frac{1}{n} \sum_{i=1}^n \tilde{R}_i^2 \quad \text{and} \quad (\mathbb{III}) := \left(\frac{1}{n} \sum_{r=1}^n R_r \right)^2 - \left(\frac{1}{n} \sum_{r=1}^n \tilde{R}_r \right)^2,$$

and show each of them converges to zero in probability. For the first term (\mathbb{I}) , we have

$$\frac{1}{n} \sum_{i=1}^n R_i^2 - \frac{1}{n} \sum_{i=1}^n \tilde{R}_i^2 = \underbrace{\frac{1}{n} \sum_{i=n_1+1}^{\bar{n}_1} (R_i^2 - \tilde{R}_i^2) \cdot \mathbf{1}(\bar{n}_1 > n_1)}_{:= \tilde{\Delta}_1} + \underbrace{\frac{1}{n} \sum_{i=\bar{n}_1+1}^{n_1} (R_i^2 - \tilde{R}_i^2) \cdot \mathbf{1}(\bar{n}_1 \leq n_1)}_{:= \tilde{\Delta}_2}.$$

Using the law of total expectation, we obtain

$$\begin{aligned} \mathbb{E}[|\tilde{\Delta}_1|] &= \mathbb{E}\left[\left|\frac{1}{n} \sum_{i=n_1+1}^{\bar{n}_1} (R_i^2 - \tilde{R}_i^2) \cdot \mathbf{1}(\bar{n}_1 > n_1)\right|\right] \\ &\leq \mathbb{E}\left[\frac{\bar{n}_1 - n_1}{n} \cdot \frac{1}{\bar{n}_1 - n_1} \sum_{i=n_1+1}^{\bar{n}_1} \mathbf{1}(\bar{n}_1 > n_1) \cdot \mathbb{E}\{|R_i^2 - \tilde{R}_i^2| \mid \bar{n}_1\}\right] \\ &\leq \frac{2}{n} \mathbb{E}[|\bar{n}_1 - n_1|] \text{Var}(Y_1) \leq \frac{\text{Var}(Y_1)}{\sqrt{n}}, \end{aligned}$$

where the first inequality is derived from Jensen's inequality, and the remaining steps follow from the previous results. A similar argument applies to $\tilde{\Delta}_2$. By Markov's inequality, (\mathbb{I}) converges to zero. For the second term (\mathbb{III}) , we have

$$\left| \left(\frac{1}{n} \sum_{r=1}^n R_r \right)^2 - \left(\frac{1}{n} \sum_{r=1}^n \tilde{R}_r \right)^2 \right| \leq \left| \frac{1}{n} \sum_{r=1}^n R_r + \frac{1}{n} \sum_{r=1}^n \tilde{R}_r \right| \cdot \left| \frac{1}{n} \sum_{r=1}^n R_r - \frac{1}{n} \sum_{r=1}^n \tilde{R}_r \right|,$$

which can be shown to converge to zero in probability using the previous results.

3. **Terms (c) and (d):** $\hat{\sigma}_R^2 - \sigma_{\tilde{R}}^2$ and $\hat{\sigma}_{\tilde{R}}^2 - \sigma_R^2$. We can see that the term (d) converges to zero in probability by the conventional law of large numbers. The term (c) also converges to zero as well since $\hat{\sigma}_R^2 - \sigma_{\tilde{R}}^2 = (c)' + (c)''$, where $(c)' = \hat{\sigma}_R^2 - \hat{\sigma}_{\tilde{R}}^2$ and $(c)'' = \hat{\sigma}_{\tilde{R}}^2 - \sigma_{\tilde{R}}^2$. In fact, $(c)' = (b)$ and $(c)'' = (d)$, and both are known to converge to zero in probability based on the previous results.

This completes the proof which shows that T and \tilde{T} are asymptotically equivalent for the stable case.

B Supporting Lemmas

In this section, we collect several lemmas from the existing literature for completeness. The proof of the following lemma can be found, for example, in [Mulzer \(2018\)](#).

Lemma 1. *Let Z_1, \dots, Z_n be i.i.d. Bernoulli random variables with success probability $p \in [0, 1]$ and $S_n = \sum_{i=1}^n Z_i$. For any $\delta \in [0, 1]$, it holds that*

$$\mathbb{P}\{S_n \geq (1 + \delta)np\} \leq e^{-\frac{np\delta^2}{3}} \quad \text{and} \quad \mathbb{P}\{S_n \leq (1 - \delta)np\} \leq e^{-\frac{np\delta^2}{3}}.$$

The following is the uniform central limit theorem result in [Shah and Peters \(2020, Lemma 18\)](#).

Lemma 2. ([Shah and Peters, 2020, Lemma 18](#)) *Let \mathcal{P} be a family of distributions for a random variable $\zeta \in \mathbb{R}$ and suppose that ζ_1, ζ_2, \dots are i.i.d. copies of ζ . For each $n \in \mathbb{N}$, let $S_n := n^{-1/2} \sum_{i=1}^n \zeta_i$. Suppose that for all $P \in \mathcal{P}$, we have $\mathbb{E}_P(\zeta) = 0$, $\mathbb{E}_P(\zeta^2) = 1$ and $\mathbb{E}_P(|\zeta|^{2+\eta}) < c$ for some $\eta, c > 0$. We have that*

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(S_n \leq t) - \Phi(t)| = 0.$$

The next lemma corresponds to Lundborg et al. (2022, Lemma S8) on conditional uniform central limit theorem.

Lemma 3. (*Lundborg et al., 2022, Lemma S8*) Let $(X_{n,i})_{n \in \mathbb{N}, i \in [n]}$ be a triangular array of real-valued random variables and let $(\mathcal{F}_n)_{n \in \mathbb{N}}$ be a filtration on \mathcal{F} . Assume that

1. $X_{n,1}, \dots, X_{n,n}$ are conditionally independent given \mathcal{F}_n , for each $n \in \mathbb{N}$;
2. $\mathbb{E}_P(X_{n,i} | \mathcal{F}_n) = 0$ for all $n \in \mathbb{N}, i \in [n]$;
3. $|n^{-1} \sum_{i=1}^n \mathbb{E}_P(X_{n,i}^2 | \mathcal{F}_n) - 1| = o_P(1)$;
4. there exists $\delta > 0$ such that

$$\frac{1}{n} \sum_{i=1}^n \mathbb{E}_P(|X_{n,i}|^{2+\delta} | \mathcal{F}_n) = o_P(n^{\delta/2}).$$

Then $S_n = n^{-1/2} \sum_{m=1}^n X_{n,m}$ converges uniformly in distribution to $N(0, 1)$, i.e.

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}} \sup_{x \in \mathbb{R}} |\mathbb{P}_P(S_n \leq x) - \Phi(x)| = 0.$$

The next lemma corresponds to Lundborg et al. (2022, Lemma 20) on uniform Slutsky's theorem.

Lemma 4. (*Shah and Peters, 2020, Lemma 20*) Let \mathcal{P} be a family of distributions that determines the law of a sequences $(V_n)_{n \in \mathbb{N}}$ and $(W_n)_{n \in \mathbb{N}}$ of random variables. Suppose

$$\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(V_n \leq t) - \Phi(t)| = 0.$$

Then we have the following.

- (a) If $W_n = o_P(1)$, we have $\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(V_n + W_n \leq t) - \Phi(t)| = 0$.
- (b) If $W_n = 1 + o_P(1)$, we have $\lim_{n \rightarrow \infty} \sup_{P \in \mathcal{P}} \sup_{t \in \mathbb{R}} |\mathbb{P}_P(V_n/W_n \leq t) - \Phi(t)| = 0$.

C Additional Numerical Experiments

In this section, we present the experimental details and additional numerical experiments not included in the main body of the paper. Specifically, we describe the experimental details in Appendix C.1 and provide additional simulation results, including an empirical analysis of the impact of density ratio estimation errors on DRT methods in Appendix C.2, a real data analysis for CIT methods in Appendix C.3, a sensitivity analysis of the application of Algorithm 1 to the CIT approach in Appendix C.4, and a sensitivity analysis of Algorithm 1 to adjustment parameter ε in Appendix C.5.

C.1 Experimental Details

We begin with the implementation details of our numerical experiments, including density ratio estimation techniques, linear-time MMD test, classifier-based test, and conditional independence testing approaches.

Density Ratio Estimation. We estimate the density ratio $r_X(x)$ defined in (4) using a probabilistic classification-based approach described in Sugiyama et al. (2010, Section 3). Specifically, we focus on two classifiers: linear logistic regression (LL) and kernel logistic regression (KLR).

Given samples $\{(X_i^{(1)}, Y_i^{(1)})\}_{i=1}^{n_1} \stackrel{\text{i.i.d.}}{\sim} P_{XY}^{(1)}$ and $\{(X_j^{(2)}, Y_j^{(2)})\}_{j=1}^{n_2} \stackrel{\text{i.i.d.}}{\sim} P_{XY}^{(2)}$, consider $\{(X_i, \ell_i)\}_{i=1}^n$, where $(X_1, \dots, X_n) = (X_1^{(1)}, \dots, X_{n_1}^{(1)}, X_1^{(2)}, \dots, X_{n_2}^{(2)})$ and $(\ell_1, \dots, \ell_{n_1}, \ell_{n_1+1}, \dots, \ell_n) = (0, \dots, 0, 1, \dots, 1)$ with $\ell_i = \mathbf{1}(i \geq n_1 + 1)$. Further denote $X_i = (X_i(1), \dots, X_i(p))^\top$ where p is the dimension of X_i and let $\beta := (\beta_0, \beta_1, \dots, \beta_p)^\top$.

- For LL method, we model the posterior probability as

$$\eta(X_i; \beta) = \mathbb{P}(\ell = 1 | X_i) = \frac{1}{1 + \exp(-\beta_0 + \sum_{j=1}^p \beta_j X_i(j))}.$$

The estimated coefficients $\hat{\beta}$ are obtained by minimizing the negative log-likelihood.

- For KLR method (Zhu and Hastie, 2005), we use $\eta(X_i; \beta) = 1 / (1 + \exp(-\theta(X_i; \beta)))$, where $\theta(X_i; \beta) = \beta_0 + \sum_{j=1}^p \beta_j k(X_i(j), x)$ and $k(x, y) = \exp(-\|x - y\|^2/\sigma^2)$. The estimated coefficients $\hat{\beta}$ are obtained by minimizing the following penalized negative log-likelihood:

$$-\sum_{i=1}^n [\ell_i \theta(X_i; \beta) - \log(1 + \exp(\theta(X_i; \beta)))] + \frac{\lambda}{2} \|\theta\|_{\mathcal{H}_k}^2,$$

where \mathcal{H}_k is the reproducing kernel Hilbert space generated by k and λ is a regularization parameter.

The density ratio estimate is then:

$$\hat{r}_X(X_i) = \frac{n_2}{n_1} \cdot \frac{\eta(X_i; \hat{\beta})}{1 - \eta(X_i; \hat{\beta})}. \quad (12)$$

For the joint density ratio, we use (X_i, Y_i) instead of X_i alone. We set $\sigma^2 = 200$, following Hu and Lei (2024), and fix $\lambda = 0.0005$ throughout our simulations.

Linear-Time MMD Test. For the linear-time MMD tests, we use a Gaussian kernel with the bandwidth parameter fixed at 1 across all experiments. In the cross-validated version (${}^{\dagger}\text{MMD-}\ell$), we use 2-fold cross validation (i.e., $K = 2$) with an equal splitting ratio.

Classifier-based Test. As mentioned in Section 4.1, under the balanced-sample setting, the Bayes optimal classifier is defined as:

$$h^*(y, x) := \mathbf{1} \left(\frac{f_{YX}^{(1)}(y, x)}{f_{YX}^{(1)}(y, x) + f_{YX}^{(2)}(y, x)} > \frac{1}{2} \right).$$

This classifier can be equivalently expressed using density ratios:

$$h^*(y, x) = \mathbf{1} \left(\frac{r_X(x)}{r_X(x) + r_{YX}(y, x)} > \frac{1}{2} \right),$$

where $r_X(x) = f_X^{(1)}(x)/f_X^{(2)}(x)$ as in (4) and $r_{YX}(y, x) := f_{YX}^{(1)}(y, x)/f_{YX}^{(2)}(y, x)$. The empirical classifier is then defined as a plug-in estimator of h^* :

$$\hat{h}(y, x) := \mathbf{1} \left(\frac{\hat{r}_X(x)}{\hat{r}_X(x) + \hat{r}_{YX}(y, x)} > \frac{1}{2} \right), \quad (13)$$

where \hat{r}_X and $\hat{r}_{Y|X}$ are the estimated marginal and joint density ratios, respectively, obtained using the classification-based approach described above.

The classifier \hat{h} is constructed based on the training set, whereas the testing set is further split into two subsets with a ratio of 8:2. A larger subset is used for density ratio estimation, and the other subset is used for calculating the test statistic. In the cross-validated version (${}^{\dagger}\text{CLF}$), we use 2-fold cross-validation as in ${}^{\dagger}\text{MMD}-\ell$, maintaining 8:2 splitting ratio within each fold.

Randomized Conditional Independence Test. The RCIT method is implemented using the default hyperparameter settings specified in Strobl et al. (2019). Specifically, we use the default approximation method (Lindsay–Pilla–Basak method) for the null distribution. The number of random Fourier features is set to 100 for the conditioning set, and 5 for the non-conditioning sets.

Regression Methods for CIT. In our implementation of CIT approaches, we utilize several standard regression techniques. Table 1 below provides an overview of the key methods and their corresponding hyperparameter settings.

Table 1: Description of regression methods used in CIT approach.

Regression Method	R Implementation	Tuning Parameters	Description
Random Forests	<code>ranger</code>	$mtry = \sqrt{p}$	# of variables to split at each node
XGBoost	<code>xgboost</code>	max depth = 6 $\eta = 0.3$	maximum tree depth learning rate

The code for reproducing all of our simulation results (including those in Appendices C.2 to C.5) and for more detailed settings is available on GitHub: <https://github.com/suman-cha/CDTST>.

C.2 Impact of Density Ratio Estimation Errors on DRT Methods

To complement the real data analysis presented in Section 5.3, we conduct experiments to examine the relationship between density ratio estimation errors and type I errors of DRT methods. Figure 5 illustrates the log-scaled mean squared error (MSE) of the marginal density ratio r_X and the conditional density ratio $r_{Y|X}$ estimates for both the LL and KLR methods across various sample sizes. Our experimental setup involves 500 simulations for each combination of sample size, dataset, and estimation method. We report the median MSE to provide a robust measure of estimation accuracy. For better visualization, the log-scaled MSE values are clipped: marginal density ratio errors above 10 are capped at 10 and conditional density ratio errors are limited to a maximum of 1. Notably, the true error values for the LL method in high-dimensional settings, significantly exceed these clipped limits.

In the low-dimensional diamonds dataset, both LL and KLR methods show relatively low MSE values for both marginal and conditional density ratio estimation. The performance gap between LL and KLR methods diminishes as the sample size increases. This observation aligns with the findings in Section 5.3, where simpler methods like LL suffice to control the type I error in low-dimensional settings. The low estimation errors explain their similar performance in such scenarios.

In contrast, the high-dimensional superconductivity dataset shows significant differences between the methods. KLR consistently outperforms LL in both marginal and conditional density ratio estimation, maintaining low and stable MSE values across all sample sizes. On the other hand, LL shows extremely high MSE values, particularly for small sample sizes. Although LL shows some improvement as the sample size increases, it remains inferior to KLR in terms of estimation accuracy.

The high estimation errors for LL in high-dimensional settings, even beyond the clipping applied, account for poor type I error control observed in Figure 4. These results highlight the need for more advanced density

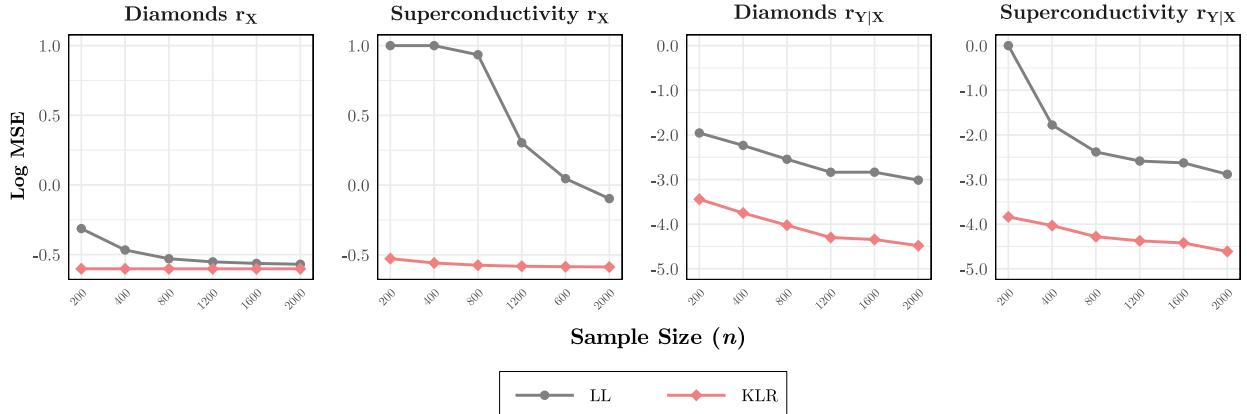


Figure 5: Log-scaled mean squared errors of marginal density ratio $r_X(x)$ (*left*) and conditional density ratio $r_{Y|X}(y|x)$ (*right*) estimates for LL and KLR methods across various sample sizes. Results are shown for diamonds and superconductivity datasets, based on median values from 500 simulations under the null hypothesis.

ratio estimation techniques to ensure the validity of tests in complex and high-dimensional scenarios.

C.3 Real Data Analysis for CIT Methods

We present the results for CIT methods applied to the diamonds and superconductivity datasets, complementing the analysis discussed in Section 5.3. Figure 6 shows the rejection rates for these methods under both the null and alternative hypotheses across various sample sizes.

For the low-dimensional diamonds dataset, the CIT methods generally exhibit good type I error control, with rejection rates close to the significance level $\alpha = 0.05$ under the null hypothesis. Under the alternative hypothesis, we observe increasing power for all methods except for WGSC as the sample size grows. Notably, RCIT and GCM show superior performance in terms of power.

In the high-dimensional superconductivity dataset, the performance of CIT methods is similar to that observed in the diamonds dataset, with no significant differences compared to DRT methods as shown in Figure 4. In terms of type I error control, GCM exhibits increasing rejection rates under the null hypothesis as the sample size grows. RCIT shows more inflated type I error, especially at small sample size. Regarding power, the CIT methods demonstrate relatively consistent performance across both datasets.

C.4 Sensitivity Analysis of CIT Methods to Algorithm 1

We examine the impact of Algorithm 1 on CIT methods across scenarios with unbounded marginal density ratios, as outlined in Section 5.2. Our analysis encompasses RCIT and the regression-based methods; PCM, GCM, and WGSC. The latter three are implemented using various regression models, such as linear models (`lm`), Random Forests (`rf`), and XGBoost (`xgb`). Tables 2 to 4 show the results for Scenarios 1(U), 2(U), and 3(U), respectively. In these tables, a checkmark (\checkmark) indicates that Algorithm 1 is applied, whereas a cross (\times) indicates it is not.

RCIT exhibits significant sensitivity to the application of Algorithm 1, particularly in Scenario 1(U). In this case, without the algorithm, the type I error rates of RCIT increase with sample size, whereas with the algorithm, these rates decrease as the sample size grows. This behavior highlights the potential stabilizing effect of Algorithm 1 on the performance of RCIT. On the other hand, GCM shows consistent performance across different regression methods, suggesting the robustness to the choice of underlying regression models.

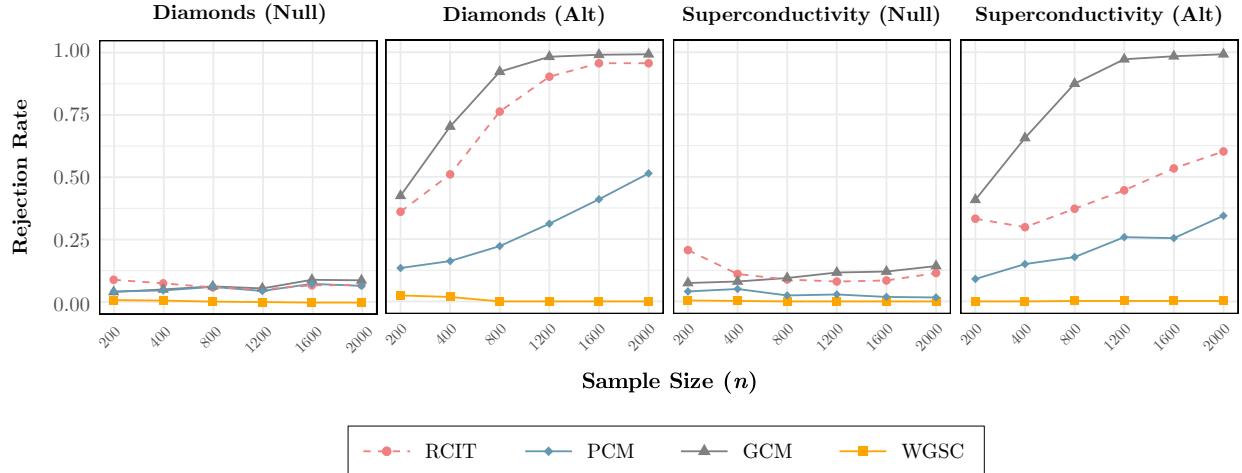


Figure 6: Rejection rates of CIT methods on the diamonds and superconductivity datasets under null and alternative hypotheses. Results are averaged over 500 repetitions with significance level $\alpha = 0.05$.

In contrast, the performance of PCM varies significantly with the choice of regression method. WGSC shows inconsistent patterns across scenarios, indicating potential sensitivity to specific data properties or model assumptions. This variability underscores the need for careful consideration when applying WGSC for conditional two-sample testing.

C.5 Sensitivity Analysis of Algorithm 1 to ε

We conduct a sensitivity analysis of Algorithm 1 with respect to the adjustment parameter ε . This parameter determines the size of the constructed testing dataset $\mathcal{D}_{\tilde{n}}$, through the equation $\tilde{n} = kn$, where k is a function of ε , as defined in Algorithm 1. The goal of this section is to examine the impact of different ε values on the performance of conditional two-sample testing via CIT methods. We consider 3 candidates for ε : $\{1/n, 1/\log(n), 1/\sqrt{\log(n)}\}$. The analysis covers scenarios with unbounded marginal density ratios, as described in Section 5.2, examining both null and alternative hypotheses across different sample sizes. The settings of these experiments correspond to the settings detailed in Section 5 with $n \in \{200, 500, 1000\}$.

Tables 5 to 7 present the results for each scenario, comparing the performance of RCIT, PCM, GCM, and WGSC across different ε values and sample sizes. Although the performance of the CIT methods remains relatively stable across different choices of ε , some variations in rejection rates are observed, particularly for small sample sizes.

Table 2: Performance comparison of CIT methods for Scenario 1(U) under null and alternative hypotheses. Rejection rates are provided for RCIT and regression-based methods (PCM, GCM, WGSC), implemented using `lm`, `rf`, and `xgb`. Results are shown for various sample sizes, with and without Algorithm 1.

n	Hypothesis	Algorithm 1	RCIT	PCM			GCM			WGSC		
				1m	rf	xgb	1m	rf	xgb	1m	rf	xgb
200	Null	✓	0.166	0.034	0.060	0.104	0.048	0.028	0.082	0.000	0.278	0.076
		✗	0.202	0.044	0.064	0.098	0.032	0.036	0.074	0.000	0.062	0.052
	Alternative	✓	0.240	0.038	0.068	0.102	0.164	0.076	0.144	0.000	0.284	0.082
		✗	0.362	0.046	0.078	0.104	0.300	0.178	0.268	0.000	0.082	0.074
500	Null	✓	0.056	0.066	0.066	0.074	0.036	0.024	0.060	0.000	0.248	0.050
		✗	0.210	0.040	0.060	0.078	0.046	0.050	0.064	0.000	0.070	0.076
	Alternative	✓	0.228	0.070	0.092	0.100	0.356	0.170	0.296	0.000	0.252	0.048
		✗	0.536	0.038	0.138	0.126	0.552	0.444	0.602	0.000	0.086	0.078
1000	Null	✓	0.064	0.050	0.050	0.060	0.048	0.024	0.060	0.000	0.270	0.056
		✗	0.296	0.056	0.038	0.068	0.044	0.062	0.058	0.000	0.068	0.026
	Alternative	✓	0.394	0.050	0.098	0.124	0.586	0.262	0.516	0.000	0.268	0.054
		✗	0.762	0.054	0.178	0.182	0.764	0.650	0.858	0.000	0.054	0.044
2000	Null	✓	0.048	0.054	0.030	0.056	0.046	0.028	0.062	0.000	0.282	0.016
		✗	0.500	0.036	0.042	0.056	0.048	0.054	0.048	0.000	0.038	0.022
	Alternative	✓	0.610	0.050	0.112	0.110	0.796	0.502	0.820	0.000	0.282	0.024
		✗	0.914	0.036	0.292	0.272	0.890	0.812	0.998	0.000	0.086	0.028

Table 3: Performance comparison of CIT methods for Scenario 2(U) under null and alternative hypotheses. Rejection rates are provided for RCIT and regression-based methods (PCM, GCM, WGSC), implemented using `lm`, `rf`, and `xgb`. Results are shown for various sample sizes, with and without Algorithm 1.

n	Hypothesis	Algorithm 1	RCIT	PCM			GCM			WGSC		
				1m	rf	xgb	1m	rf	xgb	1m	rf	xgb
200	Null	✓	0.162	0.036	0.074	0.098	0.072	0.044	0.084	0.000	0.256	0.060
		✗	0.116	0.046	0.048	0.090	0.048	0.058	0.110	0.000	0.056	0.040
	Alternative	✓	0.730	0.038	0.466	0.188	0.064	0.044	0.104	0.000	0.284	0.192
		✗	0.808	0.050	0.504	0.254	0.054	0.050	0.110	0.000	0.200	0.174
500	Null	✓	0.074	0.056	0.060	0.104	0.044	0.046	0.094	0.000	0.270	0.042
		✗	0.094	0.042	0.042	0.076	0.060	0.056	0.072	0.000	0.058	0.056
	Alternative	✓	0.906	0.056	0.864	0.564	0.054	0.050	0.092	0.000	0.320	0.330
		✗	0.880	0.040	0.924	0.702	0.044	0.048	0.074	0.000	0.364	0.494
1000	Null	✓	0.086	0.046	0.032	0.070	0.048	0.020	0.056	0.000	0.272	0.030
		✗	0.070	0.048	0.040	0.094	0.048	0.044	0.094	0.000	0.046	0.040
	Alternative	✓	0.968	0.046	0.990	0.974	0.052	0.026	0.066	0.000	0.398	0.604
		✗	0.954	0.048	1.000	0.998	0.040	0.048	0.080	0.000	0.724	0.832
2000	Null	✓	0.048	0.036	0.030	0.064	0.032	0.018	0.048	0.000	0.272	0.030
		✗	0.076	0.040	0.036	0.070	0.042	0.046	0.072	0.000	0.046	0.020
	Alternative	✓	0.980	0.034	1.000	1.000	0.030	0.020	0.044	0.000	0.486	0.900
		✗	0.980	0.042	1.000	1.000	0.040	0.026	0.066	0.000	0.960	0.994

Table 4: Performance comparison of CIT methods for Scenario 3(U) under null and alternative hypotheses. Rejection rates are provided for RCIT and regression-based methods (PCM, GCM, WGSC), implemented using `lm`, `rf`, and `xgb`. Results are shown for various sample sizes, with and without Algorithm 1.

n	Hypothesis	Algorithm 1	RCIT	PCM			GCM			WGSC		
				1m	rf	xgb	1m	rf	xgb	1m	rf	xgb
200	Null	✓	0.166	0.048	0.062	0.100	0.052	0.028	0.096	0.000	0.254	0.064
		✗	0.142	0.042	0.048	0.086	0.058	0.056	0.098	0.000	0.066	0.056
	Alternative	✓	0.688	0.046	0.726	0.578	0.214	0.204	0.266	0.000	0.406	0.636
		✗	0.718	0.040	0.758	0.646	0.236	0.248	0.282	0.000	0.658	0.670
500	Null	✓	0.076	0.058	0.064	0.072	0.046	0.030	0.088	0.000	0.268	0.038
		✗	0.066	0.042	0.034	0.070	0.046	0.050	0.080	0.000	0.056	0.060
	Alternative	✓	0.734	0.062	0.804	0.762	0.268	0.250	0.294	0.000	0.426	0.708
		✗	0.764	0.046	0.808	0.794	0.264	0.270	0.292	0.000	0.742	0.754
1000	Null	✓	0.064	0.052	0.034	0.072	0.050	0.034	0.066	0.000	0.272	0.020
		✗	0.086	0.048	0.046	0.072	0.036	0.036	0.080	0.000	0.052	0.036
	Alternative	✓	0.774	0.056	0.826	0.832	0.244	0.214	0.246	0.000	0.470	0.762
		✗	0.818	0.048	0.830	0.838	0.236	0.236	0.268	0.000	0.800	0.806
2000	Null	✓	0.070	0.034	0.024	0.062	0.078	0.050	0.084	0.000	0.274	0.026
		✗	0.082	0.046	0.036	0.084	0.076	0.068	0.098	0.000	0.044	0.024
	Alternative	✓	0.790	0.038	0.806	0.816	0.224	0.220	0.250	0.000	0.472	0.790
		✗	0.810	0.044	0.814	0.824	0.232	0.232	0.268	0.000	0.816	0.802

Table 5: Sensitivity analysis of Algorithm 1 for Scenario 1(U) under null and alternative hypotheses. The table shows rejection rates of four CIT methods (RCIT, PCM, GCM, WGSC) for ε values and sample sizes.

n	Hypothesis	ε	RCIT	PCM	GCM	WGSC
200	Null	$1/n$	0.164	0.064	0.022	0.072
		$1/\log(n)$	0.166	0.064	0.036	0.076
		$1/\sqrt{\log(n)}$	0.168	0.072	0.018	0.062
	Alternative	$1/n$	0.220	0.070	0.080	0.068
		$1/\log(n)$	0.240	0.068	0.076	0.082
		$1/\sqrt{\log(n)}$	0.224	0.074	0.072	0.066
500	Null	$1/n$	0.090	0.038	0.026	0.052
		$1/\log(n)$	0.056	0.066	0.024	0.050
		$1/\sqrt{\log(n)}$	0.080	0.036	0.036	0.044
	Alternative	$1/n$	0.250	0.050	0.126	0.040
		$1/\log(n)$	0.228	0.092	0.170	0.048
		$1/\sqrt{\log(n)}$	0.244	0.062	0.166	0.046
1000	Null	$1/n$	0.046	0.052	0.030	0.038
		$1/\log(n)$	0.064	0.050	0.024	0.056
		$1/\sqrt{\log(n)}$	0.062	0.040	0.030	0.040
	Alternative	$1/n$	0.376	0.078	0.258	0.048
		$1/\log(n)$	0.394	0.098	0.262	0.054
		$1/\sqrt{\log(n)}$	0.394	0.078	0.288	0.038

Table 6: Sensitivity analysis of Algorithm 1 for Scenario 2(U) under null and alternative hypotheses. The table shows rejection rates of four CIT methods (RCIT, PCM, GCM, WGSC) for ε values and sample sizes.

n	Hypothesis	ε	RCIT	PCM	GCM	WGSC
200	Null	$1/n$	0.168	0.070	0.030	0.064
		$1/\log(n)$	0.162	0.074	0.044	0.060
		$1/\sqrt{\log(n)}$	0.164	0.068	0.028	0.064
	Alternative	$1/n$	0.694	0.418	0.026	0.200
		$1/\log(n)$	0.730	0.466	0.044	0.192
		$1/\sqrt{\log(n)}$	0.724	0.500	0.038	0.178
500	Null	$1/n$	0.080	0.054	0.020	0.046
		$1/\log(n)$	0.074	0.060	0.046	0.042
		$1/\sqrt{\log(n)}$	0.078	0.040	0.024	0.038
	Alternative	$1/n$	0.874	0.824	0.028	0.276
		$1/\log(n)$	0.906	0.864	0.050	0.330
		$1/\sqrt{\log(n)}$	0.902	0.812	0.030	0.324
1000	Null	$1/n$	0.094	0.044	0.028	0.038
		$1/\log(n)$	0.086	0.032	0.020	0.030
		$1/\sqrt{\log(n)}$	0.064	0.042	0.032	0.040
	Alternative	$1/n$	0.958	0.968	0.022	0.576
		$1/\log(n)$	0.968	0.990	0.026	0.604
		$1/\sqrt{\log(n)}$	0.964	0.980	0.022	0.626

Table 7: Sensitivity analysis of Algorithm 1 for Scenario 3(U) under null and alternative hypotheses. The table shows rejection rates of four CIT methods (RCIT, PCM, GCM, WGSC) for ε values and sample sizes.

n	Hypothesis	ε	RCIT	PCM	GCM	WGSC
200	Null	$1/n$	0.160	0.070	0.026	0.070
		$1/\log(n)$	0.166	0.062	0.028	0.064
		$1/\sqrt{\log(n)}$	0.112	0.060	0.026	0.076
	Alternative	$1/n$	0.688	0.704	0.216	0.626
		$1/\log(n)$	0.688	0.726	0.204	0.636
		$1/\sqrt{\log(n)}$	0.698	0.722	0.210	0.644
500	Null	$1/n$	0.074	0.038	0.006	0.048
		$1/\log(n)$	0.076	0.064	0.030	0.038
		$1/\sqrt{\log(n)}$	0.084	0.054	0.030	0.038
	Alternative	$1/n$	0.722	0.792	0.250	0.692
		$1/\log(n)$	0.734	0.804	0.250	0.708
		$1/\sqrt{\log(n)}$	0.730	0.782	0.252	0.694
1000	Null	$1/n$	0.076	0.042	0.038	0.048
		$1/\log(n)$	0.064	0.034	0.034	0.020
		$1/\sqrt{\log(n)}$	0.060	0.038	0.032	0.036
	Alternative	$1/n$	0.766	0.822	0.216	0.738
		$1/\log(n)$	0.774	0.826	0.214	0.762
		$1/\sqrt{\log(n)}$	0.768	0.822	0.208	0.760