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# A Kolmogorov-Smirnov type test for independence between marks and points of marked point processes* 

Tonglin Zhang<br>Department of Statistics, Purdue University 250 North University Street, West Lafayette, IN 47907-2066<br>e-mail: tlzhang@purdue.edu


#### Abstract

Marked point processes are commonly used stochastic models for representing a finite number of natural hazard events located in space and time, because these kinds of data often associate measurements (i.e. marks) with locations (i.e. points) of events. Methods of marked point processes when marks and points are interacting have been proposed, but it is still necessary to know whether the interaction must be considered. This article presents a Kolmogorov-Smirnov type method to test the independence between points and marks of marked point processes. The asymptotic distribution of the test statistic under a few weak regularity conditions is derived. According to the asymptotic result, a specific way to construct the test statistic is recommended as its null distribution can be approximated by the absolute maximum of the two-dimensional standard Brownian pillow. The simulation results and real data analyses demonstrated that the proposed method is powerful in detecting weak dependence between marks and points and performs well with a moderate sample size.


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## 1. Introduction

Marked point processes are important models for a wide variety of scientific disciplines. A typical way to describe a marked point process is to use an artificial order such that data can be represented as

$$
\sum_{i=1}^{n} \epsilon_{\left(S_{i}, M_{i}\right)}
$$

where $n$ is the total number of events, $S_{i}$ are the locations of points, $M_{i}$ are the associated marks, and $\epsilon$ is the Dirac measure. Both spatial and spatiotemporal processes can be represented by the above expression: if $S_{i}$ represents a spatial

[^0]location, then the marked point process is purely spatial; if $S_{i}$ represents a spatiotemporal location, then the marked point process is spatiotemporal.

Specific statistical methods include variogram analysis, various kinds of kriging, and geostatistical simulation techniques may be used to model a marked point process [7], but these methods rely on a fundamental assumption that point locations appear independently of marks because the correlation function used in a geostatistical method often ignores the point distributions [14]. However, the independence assumption is often violated in real world data. For instance, the relative positions of trees in a forest have repercussions on their size owing to their competition for light or nutrients [48], indicating that tree sizes and locations of trees may not be independent. Forest wildfire activities exhibit power-law relationships between frequency and burned area [30], indicating that the burned area and the locations of forest wildfires may not be independent. The purpose of this article is to explore a Kolmogorov-Smirnov type method to test the independence between points and marks of marked point processes.

It is especially convenient in modeling, estimation, and prediction in a marked point process if marks and points are independent. Many commonly used Hawkes models, such as the epidemic-type aftershock sequences (ETAS) model [43], may exhibit the independence between marks and points [50]. The spatstat [2] and PtProcess [18] packages in $\mathbf{R}$ have been developed for marked point processes. Methods of marked point processes when points and marks are interacted have been proposed $[20,32,38]$. However, it is still important to answer whether it is necessary to account for the interaction between points and marks. Recently, a few methods have been proposed to assess the independence between points and marks, including a test for stationarity and isotropy of a marked point process using variograms [1, 48], a nonparametric kernel-based test to assess the separability of the conditional intensity function [50], and a $\chi^{2}$-based test to assess the interaction between points and marks for marked point processes with a stationary spatial projection of the marked point process [16]. These methods, however, contain kernel estimation of a nonparametric smoothing function which may result in a loss of power.

In this article, we propose a Kolmogorov-Smirnov type statistic ([53], p. 265) to test the independence between marks and points. Intuitively, the test statistic is derived by computing the absolute maximum difference between the joint empirical process and the product of the marginal empirical processes of marks and points, where the maximum is computed over a collection of subsets of the study area. We show that the choice of the collection of subsets affects the asymptotic distribution of the test statistic. With a careful choice of the collection, the asymptotic null distribution may be distribution-free and approximated by the absolute maximum of the two-dimensional standard Brownian pillow.

The Brownian pillow is an extension of the Brownian bridge from $[0,1]$ to $[0,1]^{d}$ for $d \geq 1$. It is the Brownian bridge if $d=1$. Many test statistics for multivariate independence converge to the Brownian pillow under the null hypothesis $[5,11,12,39,40,49]$. Many articles have evaluated the probability distribution of the Brownian bridge. However, there are just a few for the Brownian pillow. For example, analytical bounds for the maximum of the Brownian pillow were
derived $[9,15,27,29,52]$, but there are few known results related to more important features of the Brownian pillow, such as the maximum absolute value [28]. Therefore, a simulation method is natural and convenient to be used.

Marked point processes are often used to model a number of natural hazard events located in space and time. Many successful applications of marked point processes can be found in the literature. These include marked point modeling and prediction of earthquakes, where each earthquake is represented by a magnitude and a space-time coordinate [21, 42, 54]. The three-dimensional space coordinate contains the longitude, latitude and depth of earthquake occurrences. Marked point processes for wildfires have been discussed by [45], where each wildfire is represented by its area burned and space-time coordinate. The two-dimensional space coordinate contains the longitude and latitude of wildfire occurrences. A few statistical methods of marked point processes with the independence assumption may be used [34, 46]. If the independence assumption is violated, one may consider intensity-dependent methods [20, 32, 38]. However, a test for the independence between marks and points is still important in the analysis.

To the best of our knowledge, the present article provides the first KolmogorovSmirnov type approach to assess the independence between marks and points of marked point processes. In principle, the Kolmogorov-Smirnov type approach used in this paper is nonparametric. It is easily implemented in all kinds of distributions of marks and points. The proposed method can be applied to both continuous and discrete marks. It is prudent, however, to start with the continuous marks for our asymptotic results because the core derivation does not contain complicated mathematical inferences.

The rest of the article is organized as follows. In Section 2, we briefly review the classical Kolmogorov-Smirnov test for independence. In Section 3, we propose the test statistic and derive its asymptotic distribution. In Section 4, we provide our simulation results, focusing on the marked Poisson process, the marked mixed Poisson process, and the marked Neyman-Scott cluster process. In Section 5, we apply our test statistic to real world data: the Ambrosia dumosa data [35] and the Alberta wildfire data. Our results showed that the marks and the points were not independent in both data, and a supplementary nonparametric method was proposed and applied to the Alberta wildfire data to account for dependence between points and marks. In the end, we provide a discussion.

## 2. Kolmogorov-Smirnov test for multivariate independence

The Kolmogorov-Smirnov test for multivariate independence has been discussed about fifty years ago (e.g. [5]). It usually considers a $p$-dimensional multivariate variable $\mathbf{x}=\left(X_{1}, \ldots, X_{p}\right)$ with the joint $\operatorname{CDF} F\left(x_{1}, \ldots, x_{p}\right)$. Let $F_{j}$ be the $j$-th marginal CDF of $F$. Then, the Kolmogorov-Smirnov test for multivariate independence is to assess the null hypothesis

$$
\begin{equation*}
H_{0}: F\left(x_{1}, \ldots, x_{p}\right)=\prod_{j=1}^{p} F_{j}\left(x_{j}\right) \tag{2.1}
\end{equation*}
$$

Suppose samples distributions of $F$ are observed. Then, the Kolmogorov-Smirnov statistic for multivariate independence is

$$
\begin{equation*}
T_{n}=\sup _{x_{1}, \ldots, x_{p} \in \mathbb{R}} \sqrt{n}\left|\hat{F}\left(x_{1}, \ldots, x_{p}\right)-\prod_{j=1}^{p} \hat{F}_{j}\left(x_{j}\right)\right| \tag{2.2}
\end{equation*}
$$

where $n$ is the sample size, $\hat{F}$ and $\hat{F}_{j}$ are the joint sample CDF and marginal sample CDF, respectively. It has been shown that under $H_{0}$ the limiting distribution of $T_{n}$ is a $p$-dimensional Brownian pillow, which is a Gaussian random field on $[0,1]^{p}$ with mean 0 and covariance function equal to $\prod_{j=1}^{p}\left(t_{j} \wedge t_{j}^{\prime}-t_{j} t_{j}^{\prime}\right)$ for $0<t_{j}, t_{j}^{\prime}<1$ with $j=1, \ldots, p$. To understand the Brownian pillow, we briefly review the case when $p=2$ below.

A Gaussian random field $W(u, v)$ on $[0,1]^{2}$ is called the two-dimensional standard Brownian sheet if $W[(0, v)]=W[(u, 0)]=0, E[W(u, v)]=0$, and $E\left[W(u, v) W\left(u^{\prime}, v^{\prime}\right)\right]=\left(u \wedge u^{\prime}\right)\left(v \wedge v^{\prime}\right)$ for all $(u, v) \in[0,1]^{2}$, see [56]. Important results on sample path properties of a Brownian sheet were obtained and it has been shown that the sample paths of a Brownian sheet are continuous with probability one and $W(u, v)$ has independent stationary increments on $[0,1]^{2}$, see $[8,9,44,55]$. Let

$$
\begin{equation*}
\tilde{W}(u, v)=W(u, v)-u W(1, v)-v W(u, 1)+u v W(1,1) \tag{2.3}
\end{equation*}
$$

Then, $\tilde{W}(u, v)$ is called the two-dimensional standard Brownian pillow (or the two-dimensional tied-down Brownian sheet) [12]. In general, the two-dimensional standard Brownian pillow $\tilde{W}(u, v)$ is a Gaussian random field on $[0,1]^{2}$ with $E[\tilde{W}(u, v)]=0, \tilde{W}(u, 0)=\tilde{W}(0, v)=\tilde{W}(u, 1)=\tilde{W}(1, v)=0$, and covariance function given by

$$
\begin{equation*}
E\left[\tilde{W}(u, v) \tilde{W}\left(u^{\prime}, v^{\prime}\right)\right]=\left(u \wedge u^{\prime}-u u^{\prime}\right)\left(v \wedge v^{\prime}-v v^{\prime}\right) \tag{2.4}
\end{equation*}
$$

for $0 \leq u, v, u^{\prime}, v^{\prime} \leq 1$. Neither the exact nor the approximate distribution of $\sup _{u, v \in[0,1]}|\tilde{W}(u, v)|$ has not been derived yet [26, 27, 28, 29].

We used a simulation method to evaluate the values of $\tilde{W}_{\alpha}$, where $\tilde{W}_{\alpha}$ is the upper $\alpha$ quantile of the distribution of $\sup _{u, v \in[0,1]}|\tilde{W}(u, v)|$. In the method, we independently generated $m^{2}$ random variables from $N\left(0,1 / m^{2}\right)$. Let them be denoted by $X_{i j}$ for $i, j=1, \ldots, m$. Let $Y_{k l}=\sum_{i=1}^{k} \sum_{j=1}^{l} X_{i j}$ and $Z_{k l}=Y_{k l}-$ $(l / m) Y_{k m}-(k / m) Y_{m l}+(k / m)(l / m) Y_{m m}$. Then, $E\left(Z_{k l}\right)=0$ and $E\left(Z_{k l} Z_{k^{\prime} l^{\prime}}\right)=$ $\left(k \wedge k^{\prime} / m-k k^{\prime} / m^{2}\right)\left(l \wedge l^{\prime} / m-l l^{\prime} / m^{2}\right)$. Thus, the distribution of $\sup _{1 \leq k, l \leq m}\left|Z_{k l}\right|$ was approximately equal to the distribution of $\sup _{0 \leq u, v \leq 1}|\tilde{W}(u, v)|$ if $m$ is large. In a simulation with $10^{5}$ replications for $m=1000$, we had $\tilde{W}_{0.1}=0.7298$, $\tilde{W}_{0.05}=0.7948$, and $\tilde{W}_{0.01}=0.9234$. Since we had $P\left(\sup _{u, v \in[0,1]}|\tilde{W}(u, v)|>\right.$ $1.5)<10^{-5}$, we may conclude that the $p$-value is almost 0 if $T_{n}>1.5$.

We attempt to modify the classical test statistic $T_{n}$ displayed by Equation (2.2) such that the modified test can be used to assess the independence between points and marks of marked point processes. Note that the CDF used in the
definition of $T_{n}$ cannot be naturally defined in a marked point process. We rephrase the definition of $T_{n}$ such that it can be easily modified.

Let $\mathcal{R}=\{(-\infty, t],-\infty<t<\infty\}$ and $\mathscr{R}$ be the collection of all Borel sets in $\mathbb{R}$. Then, $\mathscr{R}$ can be generated by $\mathcal{R}$. Let $\mu$ be the joint probability measure on $\mathbb{R}^{p}$ satisfying

$$
\mu\left\{\left(-\infty, x_{1}\right], \ldots,\left(-\infty, x_{p}\right]\right\}=F\left(x_{1}, \ldots, x_{p}\right)
$$

Then, $\mu$ is uniquely determined. Let $\mu_{j}$ be the $j$-th marginal probability measure of $\mu$. Then, $\mu_{j}$ satisfies

$$
\mu_{j}\left\{\left(-\infty, x_{j}\right]\right\}=F_{j}\left(x_{j}\right)
$$

The null hypothesis given by Equation (2.1) can be rephrased as

$$
H_{0}: \mu\left\{\left(-\infty, x_{1}\right], \ldots,\left(-\infty, x_{p}\right]\right\}=\prod_{j=1}^{p} \mu_{j}\left\{\left(-\infty, x_{j}\right]\right\}
$$

for all $x_{1}, \ldots, x_{p}$, which is also equivalent to

$$
H_{0}: \mu\left\{A_{1}, \ldots, A_{p}\right\}=\prod_{j=1}^{p} \mu_{j}\left\{A_{j}\right\}
$$

for all $A_{1}, \ldots, A_{p} \in \mathcal{R}$. The test statistic $T_{n}$ given by Equation (2.2) can be written as

$$
\begin{equation*}
T_{n}=\sup _{A_{1}, \ldots A_{p} \in \mathcal{R}} \sqrt{n}\left|\hat{\mu}\left(A_{1}, \ldots, A_{p}\right)-\prod_{j=1}^{p} \hat{\mu}_{j}\left(A_{j}\right)\right| \tag{2.5}
\end{equation*}
$$

where $\hat{\mu}$ and $\hat{\mu}_{j}$ are the sample measures of $\mu$ and $\mu_{j}$, respectively. In the following of this paper, we focus on a modification of $T_{n}$ given by (2.5) such that it can be used to test the independence between points and marks in a marked point process.

## 3. Method

The definition of marked point processes is well-established and can be found in $[10,25]$. In general, a marked point process is a point process defined on the product space of points and marks, but the concept has its own life in applications. To well express our method, we introduce the following notations. Let $\|\cdot\|$ be the $L_{2}-$ norm and $\|\cdot\|_{\gamma}$ be the $L_{\gamma}-$ norm for $\gamma>1$ over an Euclidean space. Let $\phi$ be the empty set, $\bar{C}$ be the complementary set of $C$. Denote $x \wedge x^{\prime}=$ $\min \left(x, x^{\prime}\right)$ if $x, x^{\prime} \in \mathbb{R}, \mathbf{x} \wedge \mathbf{x}^{\prime}=\left(x_{1} \wedge x_{1}^{\prime}, \ldots, x_{k} \wedge x_{k}^{\prime}\right)$ if $x, x^{\prime} \in \mathbb{R}^{k}$, and write $\mathbf{x} \preceq \mathbf{x}^{\prime}$ if $x_{i} \leq x_{i}^{\prime}$ for all $i=1, \ldots, k$, where $x_{i}$ and $x_{i}^{\prime}$ are the $i$-th component of $\mathbf{x}$ and $\mathbf{x}^{\prime}$, respectively.

### 3.1. Statistical formulation

A marked point process $N=\left(N_{s}, N_{m}\right)$ with points in a complete separable metric space $\mathcal{S}$ and marks in a complete separable metric space $\mathcal{M}$ is a point process on $\mathcal{S} \times \mathcal{M}$ with the additional property that the spatial projection of
marked point process $N_{s}$ is itself a point process and for any bounded $A \in \mathcal{S}$ there is $N_{s}(A)=N(A \times \mathcal{M})<\infty$, where $N(A \times B)$ is number of points in $A \times B$ for Borel sets $A \subseteq \mathcal{S}$ and $B \subseteq \mathcal{M}$. In modeling the occurrence of ecological or geographical events when depth is not involved, we may have $\mathcal{S}=\mathbb{R}^{d}$ with $d=2$ if time is not considered or $d=3$ if time is considered.

The distribution of a marked point process can be mathematically defined using the idea for a purely point process [37]. Let $n=N(\mathcal{S} \times \mathcal{M})$ be the total number of observations and assume $n$ is finite. Then, $n$ is a discrete random variable. Assume $n \geq 1$ and observations are given by an artificial order such that the data can be expressed into $\left\{\left(S_{i}, M_{i}\right): i=1, \ldots, n\right\}$. Then, a probability distribution $\pi_{n}$ can be defined on the Borel sets of $(\mathcal{S} \times \mathcal{M})^{n}$ as ([4], p. 232)
$\pi_{n}\left(A_{1} \times B_{1}, \ldots, A_{n} \times B_{n}\right)=P\left(\left(S_{1}, M_{1}\right) \in A_{1} \times B_{1}, \ldots,\left(S_{n}, M_{n}\right) \in A_{n} \times B_{n}\right)$.
To be consistent with treating the marked point process as a theory of unordered data, $\pi_{n}$ should be permutation invariant as

$$
\pi_{n}\left(A_{1} \times B_{1}, \ldots, A_{n} \times B_{n}\right)=\pi_{n}\left(A_{i_{1}} \times B_{i_{1}}, \ldots, A_{i_{n}} \times B_{i_{n}}\right)
$$

where $\left(i_{1}, \ldots, i_{n}\right)$ is a permutation of $(1, \ldots, n)$. Let

$$
\pi_{s, n}\left(A_{1}, \ldots, A_{n}\right)=\pi_{n}\left(A_{1} \times \mathcal{M}, \ldots, A_{n} \times \mathcal{M}\right)
$$

and

$$
\pi_{m, n}\left(B_{1}, \ldots, B_{n}\right)=\pi_{n}\left(\mathcal{S} \times B_{1}, \ldots, \mathcal{S} \times B_{n}\right)
$$

Then, $\pi_{s, n}$ is the marginal distribution of points and $\pi_{m, n}$ is the marginal distribution of marks.

Definition 1. The marked point process $N$ is independent if for all $n \geq 1$ there is

$$
\begin{equation*}
\pi_{n}\left(A_{1} \times B_{1}, \ldots, A_{n} \times B_{n}\right)=\pi_{s, n}\left(A_{1}, \ldots, A_{n}\right) \pi_{m, n}\left(B_{1}, \ldots, B_{n}\right) \tag{3.1}
\end{equation*}
$$

for any Borel sets $A_{1}, \ldots, A_{n} \in \mathscr{S}$ and $B_{1}, \ldots, B_{n} \in \mathscr{M}$, where $\mathscr{S}$ and $\mathscr{M}$ are the collections of all Borel sets of $\mathcal{S}$ and $\mathcal{M}$, respectively.

In applications, it is often assumed that events appear independently such that data can be modeled by a marked Poisson process. For Poisson point processes with simple ground processes (i.e., with no two points at exactly the same location), the intensity (if it exists) uniquely characterizes all of the finitedimensional distribution of the process [10, 24]. For a marked Poisson process, one may consider a similar intensity function defined as the expected rate of occurrences at a certain location of points with a certain value of marks [45]. This idea has been used to define a marked Poisson process by [50], which has been summarized in the following.
Definition 2. The marked point process $N$ is called a marked Poisson process if there exists a Borel measure $\mu$ on $\mathcal{S} \times \mathcal{M}$ such that for any disjoint $C_{1}, \ldots, C_{k}$, $N\left(C_{1}\right), \ldots, N\left(C_{k}\right)$ are independent Poisson distributed with expected values $\mu\left(C_{i}\right)$ for $i=1, \ldots, k$, respectively.

Proposition 1. The necessary condition for $N$ to be independent is

$$
\begin{equation*}
\pi(A \times B)=\pi_{s}(A) \pi_{m}(B) \tag{3.2}
\end{equation*}
$$

for any $A \in \mathscr{S}$ and $B \in \mathscr{M}$, where $\pi(A \times B)=\pi_{1}(A \times B)$, $\pi_{s}(A)=\pi_{s, 1}(A)$, and $\pi_{m}(B)=\pi_{m, 1}(B)$. If $N$ is a marked Poisson process, then condition (3.2) is also sufficient.

It can be seen that all marked point processes with intensity-dependent marks violate Equation (3.2). However, marked point processes with intensityindependent marks may also exhibit dependence between points and marks. Therefore, Equation (3.2) is only a special case of Equation (3.1). In applications, it is often assumed that events appear independently so that data can be modeled by a marked Poisson process. For Poisson point processes with simple ground processes, the first-order intensity uniquely characterizes all of the finite-dimensional distribution of the process.

### 3.2. Derivation of the test statistic

In order to avoid the complexity of Equation (3.1), our interest is to investigate whether Equation (3.2) holds. We consider the null hypothesis

$$
\begin{equation*}
H_{0}: \varphi(A, B)=0 \tag{3.3}
\end{equation*}
$$

for any $A \in \mathscr{S}$ and $B \in \mathscr{M}$ against the alternative hypothesis

$$
\begin{equation*}
H_{1}: \varphi(A, B) \neq 0 \tag{3.4}
\end{equation*}
$$

for some $A \in \mathscr{S}$ and $B \in \mathscr{M}$, where $\varphi(A, B)=\pi(A \times B)-\pi_{s}(A) \pi_{m}(B)$.
Lemma 1. A necessary condition for Equation (3.2) to be held is that $\varphi(A, B)=$ 0 for any $A \in \mathcal{A} \subseteq \mathscr{S}$ and $B \in \mathcal{B} \subseteq \mathscr{M}$. If $\mathscr{S}$ and $\mathscr{M}$ can be generated by $\mathcal{A}$ and $\mathcal{B}$, respectively, then the condition is also sufficient.

Note that $\pi(A \times B)$ is the joint probability distribution for points in $A$ and marks in $B, \pi_{s}(A)$ is the marginal probability for points in $A$, and $\pi_{m}(B)$ is the marginal probability for marks in $B$. If $n \geq 1$, then they can be estimated by $N(A \times B) / n, N_{s}(A) / n$, and $N_{m}(B) / n$, respectively. Let

$$
\begin{equation*}
\hat{\varphi}_{n}(A, B)=\hat{\pi}_{n}(A \times B)-\hat{\pi}_{s, n}(A) \hat{\pi}_{m, n}(B) \tag{3.5}
\end{equation*}
$$

where $\hat{\pi}_{n}(A \times B)=N(A \times B) / n, \hat{\pi}_{s, n}(A)=N_{s}(A) / n$, and $\hat{\pi}_{m, n}(B)=N_{m}(B) / n$ if $n>0$ and $\hat{\pi}_{n}(A \times B)=\hat{\pi}_{s, n}(A)=\hat{\pi}_{m, n}(B)=0$ if $n=0$. Then, $\hat{\varphi}_{n}(A, B)$ is an estimator of $\varphi(A, B)$. Our test statistic therefore is

$$
\begin{equation*}
T_{n}=\sup _{A \in \mathcal{A}, B \in \mathcal{B}} \sqrt{n}\left|\hat{\pi}_{n}(A \times B)-\hat{\pi}_{s, n}(A) \hat{\pi}_{m, n}(B)\right| \tag{3.6}
\end{equation*}
$$

where $\mathcal{A} \subseteq \mathscr{S}$ and $\mathcal{B} \subseteq \mathscr{M}$. Then, $T_{n}$ is a Kolmogorov-Smirnov type test statistic with $T_{n}=0$ if $n=0$. It can be seen that the test statistic $T_{n}$ given by Equation (3.6) is special case of Equation (2.5). It is recommended to reject $H_{0}$ if $T_{n}$ is large.

Even though there are many ways to choose $\mathcal{A}$ and $\mathcal{B}$ in Equation (3.6), it is not necessary to consider all of them. In this article, we only recommend to consider two of them. In the first, we choose $\mathcal{A}$ and $\mathcal{B}$ such that $\mathscr{S}$ and $\mathscr{M}$ can be generated by $\mathcal{A}$ and $\mathcal{B}$, respectively. In the second method, we choose $\mathcal{A}$ and $\mathcal{B}$ such that the $p$-value can be easily derived using the limiting distribution. Because the second method is more convenient in applications, we focus on the second method in the rest of the paper. The detail derivations of the two methods will be discussed in the next subsection.

### 3.3. Asymptotics

We derive the asymptotic properties of $T_{n}$ in this subsection. The asymptotic properties include the asymptotic null distribution, consistency, and asymptotic power functions. These properties are particularly considered in the marked Poisson process as it has independent increments ([25], p. 6), which is an important assumption in our main conclusion for asymptotic power functions given by Theorem 2. If the marked point process is not Poisson, then we can only provide the asymptotic null distribution since the basic theorem that we cite requires $N_{s}$ be strong stationary.

### 3.3.1. Asymptotics in marked poisson processes

Let $\varphi_{\mathcal{A}, \mathcal{B}}$ be the vector composed by $\varphi(A, B)$ and $\hat{\varphi}_{\mathcal{A}, \mathcal{B}, n}$ be the vector composed by $\hat{\varphi}_{n}(A, B)$ for all $A \in \mathcal{A}$ and $B \in \mathcal{B}$, respectively. Then, $\varphi_{\mathcal{A}, \mathcal{B}}$ and $\hat{\varphi}_{\mathcal{A}, \mathcal{B}, n}$ are finite if both $\mathcal{A}$ and $\mathcal{B}$ are finite; otherwise they are infinite. For example, if $\mathcal{A}=\left\{A_{1}, \ldots, A_{I}\right\}$ and $\mathcal{B}=\left\{B_{1}, \ldots, B_{J}\right\}$, then

$$
\begin{equation*}
\varphi_{\mathcal{A}, \mathcal{B}}=\left(\varphi\left(A_{1}, B_{1}\right), \varphi\left(A_{1}, B_{2}\right), \ldots, \varphi\left(A_{I}, B_{J}\right)\right) \tag{3.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\hat{\varphi}_{\mathcal{A}, \mathcal{B}, n}=\left(\hat{\varphi}_{n}\left(A_{1}, B_{1}\right), \hat{\varphi}_{n}\left(A_{1}, B_{2}\right), \ldots, \hat{\varphi}_{n}\left(A_{I}, B_{J}\right)\right) \tag{3.8}
\end{equation*}
$$

are both $I J$-dimensional vectors.
In order to derive the asymptotic distribution of $T_{n}$, we impose the following regularity conditions:
(C1) $N$ is a marked Poisson process.
(C2) The intensity function of $N_{s}$ has the form of $\lambda_{s}(\mathbf{s})=\kappa \lambda_{0}(\mathbf{s})$, where $\kappa$ is a finite positive number and $\int_{\mathcal{S}} \lambda_{0}(\mathbf{s}) d \mathbf{s}=1$.
(C3) Conditioning on points, the marks are continuous random variables with conditional density functions given by $\lambda_{1}(\mathbf{m} \mid \mathbf{s})=\lambda(\mathbf{s}, \mathbf{m}) / \lambda_{s}(\mathbf{s})$ for $\mathbf{m} \in$ $\mathcal{M}$ and $\mathbf{s} \in \mathcal{S}$.

Under Conditions (C1)-(C3), we have $N_{\kappa}(A \times B) \sim \operatorname{Poisson}(\kappa \pi(A \times B))$, $N_{\kappa}(A \times B) \mid n \sim \operatorname{Bin}(n, \pi(A \times B))$ if $n \neq 0$, and $\pi(A \times B)=\int_{A \times B} \lambda(\mathbf{s}, \mathbf{m}) d \mathbf{s} d \mathbf{m} / \kappa$ if $A \subseteq \mathcal{S}$. We denote $\xrightarrow{P}$ as convergence in probability and $\xrightarrow{D}$ as convergence in distribution as $\kappa \rightarrow \infty$.

Lemma 2. If Conditions (C1)-(C3) hold, then $\hat{\varphi}_{n}(A, B) \xrightarrow{P} \varphi(A, B)$ for any $A \in \mathscr{S}$ and $B \in \mathscr{M}$.

Lemma 3. Assume Conditions (C1)-(C3) hold. Let $\varphi_{\mathcal{A}, \mathcal{B}}$ and $\hat{\varphi}_{\mathcal{A}, \mathcal{B}, n}$ be given by (3.7) and (3.8) for $\mathcal{A}=\left\{A_{1}, \ldots, A_{I}\right\} \subseteq \mathscr{S}$ and $\mathcal{B}=\left\{B_{1}, \ldots, B_{J}\right\} \subseteq \mathscr{M}$, respectively. Let $\nu_{1 i}=\pi\left(C_{i} \times D_{i}\right), \nu_{2 j}=\pi\left(C_{j} \times D_{j}\right), u_{i j}=\pi\left[\left(C_{i} \cap C_{j}\right) \times\left(D_{i} \cap D_{j}\right)\right]$ for $i, j=1,2,3,4$, where $C_{i}=A_{i_{1}}$ or $C_{j}=A_{i_{2}}$ if $i=1,2$ or $j=1,2, C_{i}=\bar{A}_{i_{1}}$ or $C_{j}=\bar{A}_{i_{2}}$ if $i=3,4$ or $j=3,4, D_{i}=B_{j_{1}}$ or $D_{j}=B_{j_{2}}$ if $i=1,3$ or $j=1,3$, and $D_{i}=\bar{B}_{j_{1}}$ or $D_{j}=\bar{B}_{j_{2}}$ if $i=2,4$ or $j=2,4$ for $i_{1}, i_{2}=1, \ldots, I$ and $j_{1}, j_{2}=1, \ldots, J$. Then

$$
\begin{equation*}
\sqrt{n}\left(\hat{\varphi}_{\mathcal{A}, \mathcal{B}, n}-\varphi_{\mathcal{A}, \mathcal{B}}\right) \xrightarrow{D} \mathcal{N}(0, \Sigma), \tag{3.9}
\end{equation*}
$$

where

$$
\Sigma=\left(\begin{array}{cccc}
\sigma_{A_{1}, B_{1}, A_{1}, B_{1}} & \sigma_{A_{1}, B_{1}, A_{1}, B_{2}} & \cdots & \sigma_{A_{1}, B_{1}, A_{I}, B_{J}} \\
\sigma_{A_{1}, B_{2}, A_{1}, B_{1}} & \sigma_{A_{1}, B_{2}, A_{1}, B_{2}} & \cdots & \sigma_{A_{1}, B_{2}, A_{I}, B_{J}} \\
\vdots & \vdots & \ddots & \vdots \\
\sigma_{A_{I}, B_{J}, A_{1}, B_{1}} & \sigma_{A_{I}, B_{J}, A_{1}, B_{2}} & \cdots & \sigma_{A_{I}, B_{J}, A_{I}, B_{J}}
\end{array}\right)
$$

with

$$
\begin{align*}
& \sigma_{A_{i_{1}}, B_{j_{1}}, A_{i_{2}}, B_{j_{2}}} \\
= & 4 \varphi\left(A_{i_{1}}, B_{j_{1}}\right) \varphi\left(A_{i_{2}}, B_{j_{2}}\right)+4 \varphi\left(A_{i_{1}}, B_{j_{1}}\right)\left(\nu_{21} \nu_{24}-\nu_{22} \nu_{23}\right) \\
& +4 \varphi\left(A_{i_{2}}, B_{j_{2}}\right)\left(\nu_{11} \nu_{14}-\nu_{12} \nu_{13}\right) \\
& +\nu_{14}\left(u_{11} \nu_{24}-u_{12} \nu_{23}-u_{13} \nu_{22}+u_{14} \nu_{21}\right)  \tag{3.10}\\
& -\nu_{13}\left(u_{21} \nu_{24}-u_{22} \nu_{23}-u_{23} \nu_{22}+u_{24} \nu_{21}\right) \\
& -\nu_{12}\left(u_{31} \nu_{24}-u_{32} \nu_{23}-u_{33} \nu_{22}+u_{34} \nu_{21}\right) \\
& +\nu_{11}\left(u_{41} \nu_{24}-u_{42} \nu_{23}-u_{43} \nu_{22}+u_{44} \nu_{21}\right)
\end{align*}
$$

If Equation (3.2) holds, then

$$
\begin{align*}
\sigma_{A_{i_{1}}, B_{j_{1}}, A_{i_{2}}, B_{j_{2}}}= & {\left[\pi_{s}\left(A_{i_{1}} \cap A_{i_{2}}\right)-\pi_{s}\left(A_{i_{1}}\right) \pi_{s}\left(A_{i_{2}}\right)\right] }  \tag{3.11}\\
& \times\left[\pi_{m}\left(B_{j_{1}} \cap B_{j_{2}}\right)-\pi_{m}\left(B_{j_{1}}\right) \pi_{m}\left(B_{j_{2}}\right)\right] .
\end{align*}
$$

Corollary 1. If Conditions (C1)-(C3) hold, then for any $A \in \mathscr{S}$ and $B \in \mathscr{M}$ there is $\sqrt{n}\left[\hat{\varphi}_{n}(A, B)-\varphi(A, B)\right] \xrightarrow{D} N\left(0, \sigma_{A, B}^{2}\right)$, where

$$
\begin{align*}
\sigma_{A, B}^{2}= & 4 \varphi^{2}(A, B)+8 \varphi(A, B)[\pi(A \times \bar{B}) \pi(\bar{A} \times B)-\pi(A \times B) \pi(\bar{A} \times \bar{B})] \\
& +\left[\pi^{2}(A \times B) \pi(\bar{A} \times \bar{B})+\pi(A \times B) \pi^{2}(\bar{A} \times \bar{B})\right.  \tag{3.12}\\
& \left.+\pi^{2}(A \times \bar{B}) \pi(\bar{A} \times B)+\pi(A \times \bar{B}) \pi^{2}(\bar{A} \times B)\right]
\end{align*}
$$

Theorem 1. If the null hypothesis of Equation (3.3) is violated, then there is $\lim _{\kappa \rightarrow \infty} P\left(T_{n}>x\right)=1$ for any $x>0$.
Corollary 2 (Consistency). Suppose Conditions (C1)-(C3) hold. Assume $\mathscr{S}$ and $\mathscr{M}$ can be generated by $\mathcal{A}$ and $\mathcal{B}$, respectively. If the null hypothesis given by Equation (3.3) is violated, then $\lim _{\kappa \rightarrow \infty} P\left(T_{n}>x\right)=1$ for any $x>0$.

Theorem 2 (Asymptotic performance of the power function). Suppose Conditions (C1)-(C3) hold. Denote $A_{\mathbf{a}}$ and $B_{\mathbf{b}}$ as sets indexed by $\mathbf{a}=\left(a_{1}, \ldots, a_{\alpha}\right) \in$ $\mathbb{R}^{\alpha}$ and $\mathbf{b}=\left(b_{1}, \ldots, b_{\beta}\right) \in \mathbb{R}^{\beta}$, respectively, where $\alpha$ and $\beta$ are positive integers. Let $\mathcal{A}=\left\{A_{\mathbf{a}}: \mathbf{a} \in \mathbb{R}^{\alpha}\right\}$ and $\mathcal{B}=\left\{B_{\mathbf{b}}: \mathbf{b} \in \mathbb{R}^{\beta}\right\}$. Assume
(A1) $A_{\mathbf{a}^{\prime}} \subseteq A_{\mathbf{a}}$ if $\mathbf{a}^{\prime} \preceq \mathbf{a}$ and $B_{\mathbf{b}^{\prime}} \subseteq B_{\mathbf{b}}$ if $\mathbf{b}^{\prime} \preceq \mathbf{b}$;
(A2) $\lim _{\left\|\mathbf{a}^{\prime}-\mathbf{a}\right\| \rightarrow 0} A_{\mathbf{a}^{\prime}}=A_{\mathbf{a}}$ and $\lim _{\left\|\mathbf{b}^{\prime}-\mathbf{b}\right\| \rightarrow 0} B_{\mathbf{b}^{\prime}}=B_{\mathbf{b}}$;
(A3) $\lim _{a_{i} \rightarrow-\infty} A_{\mathbf{a}}=\lim _{b_{j} \rightarrow-\infty} B_{\mathbf{b}}=\phi$ for any $i$ and $j$; and
(A4) $\lim _{a_{1} \rightarrow \infty, \ldots, a_{\alpha} \rightarrow \infty} A_{\mathbf{a}}=\mathcal{S}$ and $\lim _{b_{1} \rightarrow \infty, \ldots, b_{\beta} \rightarrow \infty} B_{\mathbf{b}}=\mathcal{M}$.
Then,

$$
\begin{equation*}
\sup _{A \in \mathcal{A}, B \in \mathcal{B}} \sqrt{n}\left|\hat{\varphi}_{\mathcal{A}, \mathcal{B}, n}-\varphi_{\mathcal{A}, \mathcal{B}}\right| \xrightarrow{D} \sup _{A \in A, B \in \mathcal{B}}\left|Z_{A, B}\right| \tag{3.13}
\end{equation*}
$$

where $Z_{A, B}$ is a Gaussian random field with mean 0 and covariance function $\sigma_{A, B, A^{\prime}, B^{\prime}}$ for any $A, A^{\prime} \in \mathcal{A}$ and $B, B^{\prime} \in \mathcal{B}$.

Theorem 2 can be used to derive the asymptotic distribution of $T_{n}$ under both the null and alternative hypotheses. However, the derivation of $T_{n}$ under the null hypothesis is more important since it is related to the $p$-value of the test. If the null hypothesis holds, then the conditional distribution of $M$ does not depend on its location, which implies that $\lambda_{1}(\mathbf{m} \mid \mathbf{s})$ is independent of $\mathbf{s}$. Using Condition (C1), an iid sample of the marks is derived.

Theorem 3 (Asymptotic null distribution). Assume all conditions in Theorem 2 hold. If the null hypothesis given by Equation (3.3) also holds, then

$$
\lim _{\kappa \rightarrow \infty} P\left(T_{n} \geq x\right)=P\left(\sup _{A \in \mathcal{A}, B \in \mathcal{B}}\left|Z_{A, B}\right| \geq x\right)
$$

for all $x>0$, where $Z_{A, B}$ is a Gaussian random field with mean 0 and covariance function

$$
\tau_{A, B, A^{\prime} B^{\prime}}=\left[\pi_{s}\left(A \cap A^{\prime}\right)-\pi_{s}(A) \pi_{s}\left(A^{\prime}\right)\right]\left[\pi_{m}\left(B \cap B^{\prime}\right)-\pi_{m}(B) \pi_{m}\left(B^{\prime}\right)\right]
$$

for any $A, A^{\prime} \in \mathcal{A}$ and $B, B^{\prime} \in \mathcal{B}$.
It is generally not easy to directly apply Theorem 3 to a real-world application because the asymptotic null distribution of $T_{n}$ depends on $\pi_{s}$ and $\pi_{m}$, which are unknown. If we choose specific $\mathcal{A}$ and $\mathcal{B}$, then the asymptotic null distribution of $T_{n}$ does not depend on $\pi_{s}$ and $\pi_{m}$ any more. Then, the $p$-value of $T_{n}$ can be easily derived using Theorem 3. The method is given below.
Corollary 3. Assume all conditions in Theorem 3 hold. If $\alpha=\beta=1$ in Condition (A1) of Theorem 2, then

$$
T_{n} \xrightarrow{D} \sup _{u, v \in[0,1]}|\tilde{W}(u, v)|
$$

where $\tilde{W}$ is the standard Brownian pillow on $[0,1]^{2}$.

### 3.3.2. Asymptotic null distribution for stationary marked point processes

Even though we have derived the asymptotic null distribution and asymptotic power functions for marked Poisson process, it is generally hard to derive the asymptotic power functions when the marked point process is not Poisson. Therefore, we decide to only provide the limiting distribution of $T_{n}$ under the null hypothesis. To derive the limiting distribution of $T_{n}$ for a marked Point process which is not Poisson, we need to assume that $N_{s}$ is stationary and satisfies the strong mixing condition.

Let $N_{s}$ be a spatial point process on $\mathcal{S}=\mathbb{R}^{d}$. We say $N_{s}$ is strong stationary if the joint distribution of $\left(N_{s}\left(A_{1}\right), \ldots, N_{s}\left(A_{k}\right)\right)$ is equal to the joint distribution of $\left(N_{s}\left(A_{1}+\mathbf{s}\right), \ldots, N_{s}\left(A_{k}+\mathbf{s}\right)\right)$ for any $A_{1}, \ldots, A_{k} \in \mathscr{S}$, where $A_{i}+\mathbf{s}=\left\{\mathbf{s}^{\prime}+\mathbf{s}: \mathbf{s}^{\prime} \in\right.$ $\left.A_{i}\right\}$. Assume the $k$ th-order intensity function of $N_{s}$ exists. Then, the logarithm of the probability generating function of $N_{s}$ is defined as

$$
\psi(\zeta)=\log \left[E e^{\int_{\mathbb{R}^{d}} \log \zeta(\mathbf{s}) d N_{s}(\mathbf{s})}\right]
$$

where $\zeta, 0 \leq \zeta \leq 1$, is a function with compact support on $\mathbb{R}^{d}$. The $k$ th-order factorial cumulant measure of $N_{s}$ is defined as

$$
\begin{equation*}
C_{(k)}\left(A_{1} \times \cdots \times A_{k}\right)=\lim _{\eta \uparrow 1}\left[\frac{\partial^{k}}{\partial a_{1} \cdots \partial a_{k}} \psi\left(\eta+\sum_{i=1}^{k} a_{i} I_{A_{i}}\right)\right]_{a_{1}=\cdots=a_{k}=0} \tag{3.14}
\end{equation*}
$$

where $I_{A}$ is the indicator function of $A \in \mathscr{B}\left(\mathbb{R}^{d}\right)$. If $C_{(k)}$ is absolutely continuous, then its $k$ th-order density function, denoted by $Q_{k}\left(\mathbf{s}_{1}, \ldots, \mathbf{s}_{k}\right)$ for distinct $\mathbf{s}_{1}, \ldots, \mathbf{s}_{k} \in \mathbb{R}^{d}$, is called the $k$ th-order factorial cumulant density of $N_{s}$. Let $P$ be the distribution of $N_{s}$. We say $N_{s}$ satisfies the strong mixing condition if for any positive $r$ and $d$ there is

$$
\begin{equation*}
\lim _{a \rightarrow \infty} \xi(a r, a d)=0 \tag{3.15}
\end{equation*}
$$

where

$$
\xi(r, d)=\sup _{\substack{d\left(E_{1}, E_{2}\right) \geq r \\ d\left(E_{1}\right) \leq d, d\left(E_{2}\right) \leq d}} \sup _{\substack{U_{1} \in \mathscr{F}\left(E_{1}\right) \\ U_{2} \in \mathscr{F}\left(E_{2}\right)}}\left|P\left(U_{1} \cap U_{2}\right)-P\left(U_{1}\right) P\left(U_{2}\right)\right|,
$$

$d(E)=\sup _{\mathbf{s}, \mathbf{s}^{\prime} \in E} \rho\left(\mathbf{s}, \mathbf{s}^{\prime}\right)$ is the diameter of $E, d\left(E_{1}, E_{2}\right)=\sup _{\mathbf{s} \in E, \mathbf{s}^{\prime} \in E^{\prime}} \rho\left(\mathbf{s}, \mathbf{s}^{\prime}\right)$ is the maximum distance between disjoint sets of points and $\rho$ is the Euclidean distance function. The detailed interpretation of those definitions can be found in $[23,37]$.

The asymptotic null distribution of $T_{n}$ is also investigated by assuming $E(n) \rightarrow$ $\infty$. To explain the condition $E(n) \rightarrow \infty$, we write $\mathcal{S}_{\kappa}=[-\kappa, \kappa]^{d}$ and assume points are restricted on $\mathcal{S}_{\kappa}$. Denote $N_{\kappa}(A \times B)=N\left[\left(A \cap \mathcal{S}_{\kappa}\right) \times B\right]$ and $N_{s, \kappa}(A)=N_{\kappa}(A \times \mathcal{M})$. The asymptotic properties of $T_{n}$ is evaluated by considering

$$
T_{n_{\kappa}}=\sup _{A \in \mathcal{A}_{\kappa}, B \in \mathcal{B}} \sqrt{n_{\kappa}}\left|\hat{\varphi}_{n_{\kappa}}(A, B)\right|,
$$

as $\kappa \rightarrow \infty$, where $n_{\kappa}=N_{\kappa}\left(\mathcal{S}_{\kappa} \times \mathcal{M}\right), \mathcal{A}_{\kappa}$ is a collection of Borel sets of $\mathcal{S}_{\kappa}, \mathcal{B}$ is a collection of Borel sets of $\mathcal{M}, \hat{\varphi}_{n_{\kappa}}(A, B)=\hat{\pi}_{n_{\kappa}}(A \times B)-\hat{\pi}_{s, n_{\kappa}}(A) \hat{\pi}_{m, n_{\kappa}}(B)$ and $\hat{\pi}_{n_{\kappa}}, \hat{\pi}_{s, n_{\kappa}}, \hat{\pi}_{m, n_{\kappa}}$ are the restrictions of $\hat{\pi}, \hat{\pi}_{s}$, and $\hat{\pi}_{m}$ on $\mathcal{S}_{\kappa}$, respectively. To derive the asymptotic distirbution of $T_{n_{\kappa}}$, we need the following regularity conditions for $N_{s}$.
(D1) $N_{s}$ is strong stationary on $\mathcal{S}$.
(D2) $N_{s}$ satisfies the strong mixing condition.
(D3) The $k$ th-order factorial cumulant density functions $Q_{k}$ of $N_{s}$ satisfies

$$
\int_{\mathbf{s}_{1}, \ldots, \mathbf{s}_{k} \in \mathcal{S}}\left|Q_{k}\left(\mathbf{s}_{1}, \ldots, \mathbf{s}_{k}\right)\right| d \mathbf{s}_{1} \cdots d \mathbf{s}_{k}<C_{1}, k=2,3,4
$$

and

$$
\int_{\mathbf{s}_{1}, \ldots, \mathbf{s}_{k} \in \mathcal{S}}\left|Q_{k}\left(\mathbf{s}_{1}, \ldots, \mathbf{s}_{k}\right)\right| d \mathbf{s}_{1} \cdots d \mathbf{s}_{k-1}<C_{2}, k=2,3,4
$$

for some constants $C_{1}, C_{2} \in \mathbb{R}$.
Note that $Q_{k}$ is permutation invariant. The first integral in (D3) is over all of the $k$-th arguments but the second one is just over $k-1$ of them.

The conditional density of marks in a marked Poisson process given by Condition (C3) can also be used to describe the conditional density of marks in a stationary marked point process. Let the conditional density still be denoted by $\lambda_{1}(m \mid \mathbf{s})$. If the null hypothesis holds, then $\lambda_{1}(m \mid \mathbf{s})$ does not depend on $\mathbf{s}$, which implies that the conditional distributions of marks given their point locations are identical. If marks are also independent, then they are iid.

Theorem 4. Assume Conditions (D1)-(D3) hold and the distribution of marks is independent of their point locations. If marks are identically independently distributed, then $T_{n_{\kappa}}$ weakly converges to the absolute maximum of a Gaussian random field with mean 0 and covariance function given by Equation (3.11).

Corollary 4. Assume all conditions in Theorems 3 and 4 hold. If $\alpha=\beta=1$ in Condition (A1) of Theorem 2, then

$$
T_{n} \xrightarrow{D} \sup _{u, v \in[0,1]}|\tilde{W}(u, v)| .
$$

To choose $\mathcal{A}$ and $\mathcal{B}$ in Equation (3.6), we need to consider the asymptotic null distribution of $T_{n}$ given by Theorems 3 and 4 as well as their corollaries. There are two methods recommended. In the first, we choose $\mathcal{A}$ and $\mathcal{B}$ as the minimum collections of Borel sets in $\mathcal{S}$ and $\mathcal{M}$ such that $\mathscr{S}$ and $\mathscr{M}$ can be generated by them, respectively. A significant $T_{n}$ is enough to conclude $H_{1}$ given by Equation (3.4) and an insignificant $T_{n}$ is enough to conclude $H_{0}$ given by Equation (3.3). In the second, we choose $\mathcal{A}$ and $\mathcal{B}$ specifically such that both of them can be generated by a univariate functions from $\mathcal{S}$ to $\mathbb{R}$ or $\mathcal{M}$ to $\mathbb{R}$, respectively. A significant $T_{n}$ is enough to conclude $H_{1}$ given by Equation (3.4) but an insignificant $T_{n}$ is not enough to conclude $H_{a}$ given by Equation (3.3).

Because the second method is more convenient than the first method in applications, we focus on the second method in this paper. In the next subsection, we present a few examples about the choice of $\mathcal{S}$ and $\mathcal{M}$ in $T_{n}$.

### 3.4. Examples

We present a few examples to illustrate our ideas. We use $\alpha$ to represent the significance level and $N$ to represent a marked Poisson process (or a stationary marked point process which satisfies the strong mixing condition). In all of these examples, we focus of the choices of $\mathcal{A}$ and $\mathcal{B}$ in $T_{n}$.
Example 1. Assume $\mathcal{S}=\mathbb{R}^{2}$ and $\mathcal{M}=\mathbb{R}$ such that the points are twodimensional random vectors and the marks are one-dimensional random variables. Let $A_{\mathbf{a}}=A_{a_{1}, a_{2}}=\left(-\infty, a_{1}\right] \times\left(-\infty, a_{2}\right], B_{b}=(-\infty, b], \mathcal{A}=\left\{A_{\mathbf{a}}: \mathbf{a} \in\right.$ $\left.\mathbb{R}^{2}\right\}$, and $\mathcal{B}=\left\{B_{b}: b \in \mathbb{R}\right\}$. If $n=1$ with the observation being denoted by $(S, M)$, then $F(\mathbf{a}, b)=P\left(S \in A_{\mathbf{a}}, M \in B_{b}\right)$ is the joint CDF of $S$ and $M$, $F_{S}(\mathbf{a})=P\left(S \in A_{\mathbf{a}}\right)$ is the marginal CDF of $S$, and $F_{m}(b)=P\left(M \in B_{b}\right)$ is the marginal CDF of $M$. If $n \geq 1$, then

$$
\begin{aligned}
& T_{n}=\sup _{a \in \mathbb{R}^{2}, b \in \mathbb{R}} \sqrt{n} \left\lvert\, \frac{1}{n} N\left(\left(-\infty, a_{1}\right] \times\left(-\infty, a_{2}\right] \times(-\infty, b]\right)\right. \\
& \left.-\frac{N_{s}\left(\left(-\infty, a_{1}\right] \times\left(-\infty, a_{2}\right]\right) N_{m}((-\infty, b])}{n^{2}} \right\rvert\,
\end{aligned}
$$

If $F$ is absolutely continuous with respect to the Lebesgue measure on $\mathbb{R}^{3}$, then under $H_{0}$ there is

$$
\lim _{\kappa \rightarrow \infty} P\left(T_{n} \geq x\right)=P\left(\sup _{\mathbf{u} \in \mathbb{R}^{2}, 0 \leq v \leq 1}|Z(\mathbf{u}, v)| \geq x\right), \mathbf{u} \in[0,1]^{2}, v \in[0,1]
$$

where $Z(\mathbf{u}, v)$ is a mean zero Gaussian random field on $\mathbb{R}^{3}$ with given by $E\left[Z(\mathbf{u}, v) Z\left(\mathbf{u}^{\prime}, v^{\prime}\right)\right]=\left[F_{s}\left(\mathbf{u} \wedge \mathbf{u}^{\prime}\right)-F_{s}(\mathbf{u}) F_{s}\left(\mathbf{u}^{\prime}\right)\right]\left(v \wedge v^{\prime}-v v^{\prime}\right)$, which depends on $F_{s}$. If we choose $a_{1}=a_{2}=a$ such that $A_{\mathbf{a}}=(-\infty, a]^{2}$ with $\mathbf{a}=(a, a)$, then under $H_{0}$ there is

$$
\lim _{\kappa \rightarrow \infty} P\left(T_{n} \geq x\right)=P\left(\sup _{0 \leq u, v \leq 1}|\tilde{W}(u, v)| \geq x\right)
$$

which does not depend on $F_{s}$. Therefore, Equation (3.2) is rejected if $T_{n}>\tilde{W}_{\alpha}$.
Example 2. Let $\mathcal{S}=\mathbb{R}^{d}$ and $\mathcal{M}=\mathbb{R}^{k}$. Let $A_{\mathbf{a}}=\left\{\mathbf{x}: \mathbf{x} \preceq \mathbf{a}, \mathbf{a} \in \mathbb{R}^{d}\right\}$, $B_{\mathbf{b}}=\left\{\mathbf{x}: \mathbf{x} \preceq \mathbf{b}, \mathbf{b} \in \mathbb{R}^{k}\right\}, \mathcal{A}=\left\{A_{\mathbf{a}}: \mathbf{a} \in \mathbb{R}^{d}\right\}$, and $\mathcal{B}=\left\{B_{\mathbf{b}}: \mathbf{b} \in \mathbb{R}^{k}\right\}$. If $n>0$, then

$$
T_{n}=\sup _{\mathbf{a} \in \mathbb{R}^{d}, \mathbf{b} \in \mathbb{R}^{k}} \sqrt{n}\left|\frac{1}{n} \sum_{i=1}^{n} I_{S_{i} \preceq \mathbf{a}, M_{i} \preceq \mathbf{b}}-\left(\frac{1}{n} \sum_{i=1}^{n} I_{S_{i} \preceq \mathbf{a}}\right)\left(\frac{1}{n} \sum_{i=1}^{n} I_{M_{i} \preceq \mathbf{b}}\right)\right| .
$$

If $F$ is absolutely continuous with respect to the Lebesgue measure on $\mathbb{R}^{d+k}$, then under $H_{0}$ there is

$$
\lim _{\kappa \rightarrow \infty} P\left(T_{n} \geq x\right)=P\left(\sup _{\mathbf{u} \in \mathbb{R}^{d}, \mathbf{v} \in \mathbb{R}^{k}}|Z(\mathbf{u}, \mathbf{v})| \geq x\right)
$$

where $Z(\mathbf{u}, \mathbf{v})$ is a mean zero Gaussian random field with the covariance function given by
$E\left[Z(\mathbf{u}, \mathbf{v}) Z\left(\mathbf{u}^{\prime}, \mathbf{v}^{\prime}\right)\right]=\left[F_{s}\left(\mathbf{u} \wedge \mathbf{u}^{\prime}\right)-F_{s}(\mathbf{u}) F_{s}\left(\mathbf{u}^{\prime}\right)\right]\left[F_{m}\left(\mathbf{v} \wedge \mathbf{v}^{\prime}\right)-F_{m}(\mathbf{v}) F_{m}\left(\mathbf{v}^{\prime}\right)\right]$,
$\mathbf{u} \in \mathbb{R}^{d}$ and $\mathbf{v} \in \mathbb{R}^{k}$, which depends on both $F_{s}$ and $F_{m}$. If we choose $A_{\mathbf{a}}=$ $(-\infty, a]^{d}, B_{\mathbf{b}}=(-\infty, b]^{k}, \mathcal{A}=\left\{A_{a}: a \in \mathbb{R}\right\}$, and $\mathcal{B}=\left\{B_{b}: b \in \mathbb{R}\right\}$, then under $H_{0}$ there is

$$
\lim _{\kappa \rightarrow \infty} P\left(T_{n} \geq x\right)=P\left(\sup _{0 \leq u, v \leq 1}|\tilde{W}(u, v)| \geq x\right)
$$

Therefore, Equation (3.2) is rejected if $T_{n}>\tilde{W}_{\alpha}$.
Example 3. In Example 1, let $A_{a}=\left\{\mathbf{s}:\left\|\mathbf{s}-\mathbf{s}_{0}\right\| \leq a\right\}$ and $\mathcal{A}=\left\{A_{a}: a \in \mathbb{R}^{+}\right\}$ for some $\mathbf{s}_{0}=\left(s_{01}, s_{02}\right) \in \mathbb{R}^{2}$. Then under $H_{0}$ there is

$$
\lim _{\kappa \rightarrow \infty} P\left(T_{n} \geq x\right)=P\left(\sup _{0 \leq u, v \leq 1}|\tilde{W}(u, v)| \geq x\right)
$$

Therefore, we reject Equation (3.2) if $T_{n}>\tilde{W}_{\alpha}$.
Example 4. In Example 2, let $A_{a, \gamma}=\left\{\mathbf{s}:\left\|\mathbf{s}-\mathbf{s}_{0}\right\|_{\gamma} \leq a\right\}, B_{b, \omega}=\{\mathbf{m}$ : $\left.\left\|\mathbf{m}-\mathbf{m}_{0}\right\|_{\omega} \leq b\right\}, \mathcal{A}_{\gamma}=\left\{A_{a, \gamma}: a \in \mathbb{R}^{+}\right\}, \mathcal{B}_{\omega}=\left\{B_{b, \omega}: b \in \mathbb{R}^{+}\right\}$for some $\mathbf{s}_{0} \in \mathbb{R}^{d}, \mathbf{m}_{0} \in \mathbb{R}^{k}$, and $\gamma, \omega \geq 1$, where $\|\cdot\|_{\gamma}$ is the $L_{\gamma}$-norm. Then under $H_{0}$ there is

$$
\lim _{\kappa \rightarrow \infty} P\left(T_{n} \geq x\right)=P\left(\sup _{0 \leq u, v \leq 1}|\tilde{W}(u, v)| \geq x\right)
$$

Therefore, we reject Equation (3.2) if $T_{n}>\tilde{W}_{\alpha}$.

## 4. Simulation

We simulated realizations from marked point processes on $[0,1]^{2} \times \mathbb{R}$ for points and marks. We considered marked Poisson processes, marked mixed Poisson processes, and marked Neyman-Scott cluster processes, where only the marked Poisson processes had independent increments. We chose these processes here due to their popularity in modeling ecological data. In order to generate a marked Poisson process specifically, we first generated the purely point process $N_{s}$ and then generated the associated marks according to a selected conditional distribution function. It is known that $N_{s}$ in a marked Neyman-Scott cluster process may satisfy the stong mixing condition. However, $N_{s}$ in a marked mixed Poisson process does not satisfy the strong mixing condition because the dependence between number of points in subregions does not approach to zero if the conditional mean measure is not degenerated. Therefore, we used the marked mixed Poisson process for the robustness of the strong mixing condition in our asymptotics. In our setting, both the Neyman-Scott process and the mixed Poisson process can be thought as special cases of the Cox process (e.g. Chapter 5 in [36]).

We considered both stationary $N_{s}$ and nonstationary $N_{s}$ in our simulations. To generate data from a stationary $N_{s}$, we used the uniform intensity function
of $N_{s}$ in $[0,1]^{2}$ such that the true intensity function of $N_{s}$ was

$$
\begin{equation*}
\lambda_{0}(u, v)=Y \kappa \tag{4.1}
\end{equation*}
$$

for $0 \leq u, v \leq 1$, where $\kappa$ was a constant and $Y$ was a positive random variable. The random variable $Y$ in (4.1) was only used in marked mixed Poisson processes, where we chose $E(Y)=1$ and $V(Y)=1 / \gamma$, which was determined by the $\Gamma(\gamma, \gamma)$ distribution. For othere cases, we fixed $Y=1$. Clearly, a marked mixed Poisson process reduced to a marked Poisson process if $Y$ was a constant.

Let $S_{i}=\left(S_{i 1}, S_{i 2}\right), i=1, \ldots, n$, be the points that were generated from a simulation. The values of marks, denoted by $M_{1}, \ldots, M_{n}$, respectively, were generated independently from the $N\left(\nu_{i}, 1\right)$ distribution, where

$$
\begin{equation*}
\nu_{i}=\nu\left(S_{i}\right)=\eta \sqrt{S_{i 1}^{2}+S_{i 2}^{2}} \tag{4.2}
\end{equation*}
$$

To generate data from a nonstationary $N_{s}$, we used the intensity function of $N_{s}$ as

$$
\begin{equation*}
\lambda_{0}(u, v)=Y \kappa \beta(u) \beta(v), u, v \in[0,1] \tag{4.3}
\end{equation*}
$$

where $\kappa$ is a constant and $\beta(u)=30 u^{2}(1-u)^{2}$ is the density of the $\operatorname{Beta}(3,3)$ distribution. The values of marks were also generated independently from the $N\left(\nu_{i}, 1\right)$ distribution, where

$$
\begin{equation*}
\nu_{i}=\nu\left(S_{i}\right)=\eta \beta\left(S_{i 1}\right) \beta\left(S_{i 2}\right), \eta \geq 0 \tag{4.4}
\end{equation*}
$$

It can be seen that $\nu_{i}$ in the stationary $N_{s}$ did not depend on the intensity of $N_{s}$ but $\nu_{i}$ in the nonstationary $N_{s}$ depended on the intensity of $N_{s}$.

Before we were able to use Equation (4.2) or Equation (4.4) for marks, we needed to generate $N_{s}$ according to the spatial Poisson process, the NeymanScott cluster process, and the mixed Poisson process, respectively. To generate the marked Neyman-Scott cluster process, we first generated parent points from a Poisson point process with intensity function $\lambda_{0}(u, v) / k$ and then generated a Poisson number with mean $k$ of offsprings for each parent point, where the position of offspring points relative to its parent point was generated from an independent bivariate normal distribution with standard deviation $\gamma$. Therefore, the expected value of the number of parent points was $\kappa / k$. In all of our simulation studies, we have $E(n)=\kappa, E\left(M_{i} \mid S_{i}\right)=\nu_{i}$ and $V\left(M_{i} \mid S_{i}\right)=1$. If $\eta=0$, then $\nu_{i}=0$ for all $i$. Overall, in all the cases that we considered, we had that the marked point process was independent if and only if $\eta=0$.

We had two methods to define $T_{n}$ : one was given by Theorems 3 and 4 and the other was given by Corollaries 3 and 4 . As there were many intensity functions considered in the simulation studies, it was more convenient to focus on the latter one. To define the test statistic $T_{n}$ recommended by Corollaries 3 and 4, we used $\mathcal{A}=\left\{A_{a}: 0 \leq a \leq 1\right\}$ with $A_{a}=\left\{(u, v): \sqrt{u^{2}+v^{2}} \leq a\right\}$ for points and $\mathcal{B}=\{(-\infty, b]:-\infty<b<\infty\}$ for marks. We tested hypotheses

$$
H_{0}: \eta=0 \text { versus } H_{1}: \eta \neq 0
$$

TAble 1
Type I error probabilities and power functions of $T_{n}$ for stationary $N_{s}$ at the significance level 0.05, where $\gamma$ represented the standard deviation for clusters in the Neyman-Scott cluster process or the value of $V^{-1}(Y)$ in the mixed Poisson process

| Process | $\gamma$ | $\eta$ |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $\kappa=1000$ |  |  |  | $\kappa=5000$ |  |  |  |
|  |  | 0.0 | 0.1 | 0.2 | 0.3 | 0.0 | 0.1 | 0.2 | 0.3 |
| Poisson Mixed |  | 0.0490 | 0.1290 | 0.3913 | 0.7248 | 0.0555 | 0.4949 | 0.9831 | 1.0000 |
|  | 1 | 0.0501 | 0.1262 | 0.3389 | 0.5526 | 0.0532 | 0.4235 | 0.7495 | 0.8631 |
|  | 2 | 0.0473 | 0.1247 | 0.3643 | 0.6133 | 0.0512 | 0.4730 | 0.8594 | 0.9530 |
|  | 4 | 0.0470 | 0.1325 | 0.3838 | 0.6614 | 0.0517 | 0.4661 | 0.9126 | 0.9857 |
|  | 8 | 0.0544 | 0.1285 | 0.3859 | 0.6982 | 0.0558 | 0.4781 | 0.9537 | 0.9985 |
| Cluster$k=10$ | 16 | 0.0493 | 0.1247 | 0.3887 | 0.7072 | 0.0554 | 0.4888 | 0.9689 | 0.9997 |
|  | 0.01 | 0.0510 | 0.1377 | 0.4249 | 0.7339 | 0.0575 | 0.5077 | 0.9800 | 1.0000 |
|  | 0.05 | 0.0523 | 0.1381 | 0.4184 | 0.7399 | 0.0565 | 0.5161 | 0.9839 | 1.0000 |
|  | 0.1 | 0.0480 | 0.1395 | 0.4378 | 0.7701 | 0.0544 | 0.5413 | 0.9893 | 1.0000 |
| $\begin{aligned} & \text { Cluster } \\ & \kappa / k=10 \end{aligned}$ | 0.2 | 0.0501 | 0.1582 | 0.4923 | 0.8354 | 0.0526 | 0.6229 | 0.9968 | 1.0000 |
|  | 0.01 | 0.0511 | 0.2112 | 0.5752 | 0.7941 | 0.0560 | 0.6593 | 0.9233 | 0.9729 |
|  | 0.05 | 0.0503 | 0.1967 | 0.5204 | 0.7505 | 0.0525 | 0.6108 | 0.9031 | 0.9630 |
|  | 0.1 | 0.0508 | 0.1877 | 0.5019 | 0.7446 | 0.0558 | 0.5966 | 0.9166 | 0.9758 |
|  | 0.2 | 0.0515 | 0.1804 | 0.5185 | 0.7939 | 0.0539 | 0.6201 | 0.9626 | 0.9943 |

Table 2
Type I error probabilities and power functions of $T_{n}$ for nonstationary $N_{s}$ at the significance level 0.05, where $\gamma$ was the same as that in Table 1

| Process | $\gamma$ | $\eta$ |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $\kappa=1000$ |  |  |  | $\kappa=5000$ |  |  |  |
|  |  | 0.0 | 0.1 | 0.2 | 0.3 | 0.0 | 0.1 | 0.2 | 0.3 |
| Poisson Mixed |  | 0.0514 | 0.0961 | 0.3065 | 0.6295 | 0.0567 | 0.3874 | 0.9709 | 1.0000 |
|  | 1 | 0.0458 | 0.1050 | 0.2860 | 0.4830 | 0.0496 | 0.3589 | 0.7105 | 0.8414 |
|  | 2 | 0.0498 | 0.1032 | 0.2918 | 0.5536 | 0.0546 | 0.3843 | 0.8063 | 0.9388 |
|  | 4 | 0.0523 | 0.0996 | 0.3026 | 0.5819 | 0.0513 | 0.3772 | 0.8830 | 0.9805 |
|  | 8 | 0.0513 | 0.1029 | 0.3066 | 0.6015 | 0.0534 | 0.3886 | 0.9195 | 0.9957 |
|  | 16 | 0.0491 | 0.0955 | 0.2910 | 0.6184 | 0.0588 | 0.3914 | 0.9451 | 0.9994 |
| $\begin{gathered} \text { Cluster } \\ \kappa=10 \end{gathered}$ | 0.01 | 0.0514 | 0.1127 | 0.3405 | 0.6665 | 0.0529 | 0.4041 | 0.9620 | 0.9999 |
|  | 0.05 | 0.0509 | 0.1194 | 0.3653 | 0.6910 | 0.0584 | 0.4430 | 0.9798 | 1.0000 |
|  | 0.1 | 0.0534 | 0.1310 | 0.4348 | 0.7932 | 0.0539 | 0.5437 | 0.9963 | 1.0000 |
|  | 0.2 | 0.0516 | 0.1723 | 0.5767 | 0.9100 | 0.0577 | 0.7159 | 0.9996 | 1.0000 |
| Cluster$\kappa / k=10$ | 0.01 | 0.0520 | 0.1968 | 0.5475 | 0.7965 | 0.0583 | 0.6367 | 0.9342 | 0.9820 |
|  | 0.05 | 0.0517 | 0.1896 | 0.5028 | 0.7557 | 0.0550 | 0.5939 | 0.9184 | 0.9739 |
|  | 0.1 | 0.0517 | 0.1850 | 0.5245 | 0.7929 | 0.0555 | 0.6254 | 0.9475 | 0.9889 |
|  | 0.2 | 0.0487 | 0.1928 | 0.5885 | 0.8672 | 0.0558 | 0.6996 | 0.9860 | 0.9988 |

by $T_{n}$. We used 0.05 significance level test. By Corollary 3, we rejected $H_{0}$ if $T_{n} \geq \tilde{W}_{0.05}=0.7948$.

We simulated 10000 realizations for each selected case in the three types of marked point processes. For each realization, we used $T_{n}$ to test whether points and marks were independent. We computed the type I error probabilities (when $\eta=0$ ) and power functions (when $\eta>0$ ) for stationary $N_{s}$ (Table 1) and nonstationary $N_{s}$ (Table 2), respectively. The results showed that the type I error probabilities were all close to 0.05 . The behaviors for marked Poisson processes and marked point processes were expected because of Corollaries 3 and 4.

To evaluate the power functions, we increased $\eta$ from 0 to a certain positive value for all the cases that we have considered. We found that the power functions increased as either $\eta$ or $\kappa$ increased but they were slightly different within these processes as $\gamma$ varied. The power functions between the marked Poisson process and the marked mixed Poisson process behaved similarly, but they were slightly different from the power functions of the marked Neyman-Scott cluster processes. This was expected because the dependence between points and marks became stronger as $\eta$ increased and the expected value of the normalized constant $\sqrt{n}$ in $T_{n}$ increased as $\kappa$ increased. However, the value of $\gamma$ only affected the variability of the normalized constant $\sqrt{n}$ in $T_{n}$. The power functions of the marked Neyman-Scott cluster processes behaved differently because the presence of cluster effects might slightly affect the power functions.

## 5. Applications

We applied the proposed approach to the Ambrosia Dumosa data and the $A l$ berta Forest Wildfire data. The Ambrosia Dumosa data were previously studied by Miriti, Howe, and Wright [35] and Guan and Afshartous [16]. The Alberta Forest Wildfire data consisted of forest wildfire activities in Alberta, Canada, from 1931 to 2012. The test statistic $T_{n}$ was applied to both data, where the significance level 0.05 was used. For Alberta Forest Wildfire data, we initiated a nonparametric method to assess the dependence between points and marks when the test was significant.

### 5.1. Ambrosia dumosa data

Ecologists are interested in plant interactions which can cause low productivity in ecosystems. It has been found that both positive and negative plant interactions are common in desert plant communities [6, 22]. To understand these interactions, Ambrosia dumosa data were collected within a hectare ( $100 \times 100 \mathrm{~m}^{2}$ ) area in the Colorado Desert in 1984 [35]. Ambrosia dumosa is a drought deciduous shrub with $20-60 \mathrm{~cm}$ in height which is abundant on well drained soils below one thousand meter elevation. To formulate hypothesis about plant interaction, we investigated the spatial patterns between plant locations and plant measurements, in which we treated the locations of Ambrosia dumosa plants as points and plant measurements as marks.

The Ambrosia dumosa data consisted of locations and several important measurements of 4358 Ambrosia dumosa (Figure 1(a) and 1(b)), including the height of the plant canopy $\left(M^{1}\right)$, the length of the major axis of the plant canopy $\left(M^{2}\right)$, the length of the minor axis of the plant canopy $\left(M^{3}\right)$, and the volume of the plant canopy $\left(M^{4}\right)$. It is of interest to know whether any of them are independent of their point locations.

It has been pointed out by Guan and Afshartous [16] that the approach of Schlather et. al [48] is not appropriate because the marginal distribution of marks is not normal in this example. The nonparametric approach proposed by


Fig 1. Tree locations and density in Ambrosia dumosa plant data, and fire locations and density in Alberta Forest Wildfire data.
[50] requires nonparametric estimates of intensity functions, which relies on the choice of a kernel function and a bandwidth. A kernel-based method proposed by Guan and Afshartous [16] has been used to assess the independence between points and $M^{4}$, indicating a significance $p$-value of 0.0031 . Here we used $T_{n}$ to test the independence between points and every of $M^{1}, M^{2}, M^{3}$ and $M^{4}$, separately.

To define the test statistic $T_{n}$, we chose $\mathcal{A}=\left\{A_{a}: A_{a}=\{\|\mathbf{s}\| \leq a\}\right.$, $a \geq 0\}$ and $\mathcal{B}=\{(-\infty, b]: b \in \mathbb{R}\}$. From Lemma 1 , it is enough to conclude that points and marks were not independent if $T_{n} \geq 0.7948$. The values of $T_{n}$ for $M^{1}$, $M^{2}, M^{3}$, and $M^{4}$ were $1.5635,1.9370,1.9423$, and 2.0861 , respectively. All of the $p$-values were almost equal to 0 (all were $<10^{-5}$ ). Therefore, we concluded that the locations of trees were not independent of the tree canopy, the length of the major axis, the length of the minor axis, or the volume of the canopy.

### 5.2. Alberta wildfire data

The Canadian Alberta Forest Service initiated the modern era of wildfire record keeping in 1931. Over the years, this fire information has been recorded, stored and made available in different formats. Beginning in 1996, paper copies of wild-


Fig 2. Values of $\sqrt{n} \hat{\varphi}_{n}\left(A_{a}, B_{b}\right)$ as functions of $a$ and $\log (1+b)$ in the Alberta Forest Wildfire data (1996-2010), where the maximum was attained at around $a=5438 \mathrm{~km}$ and $b=42\left(\mathrm{~km}^{2}\right)$.
fire history information were no longer retained. The wildfire history data were entered at the field level on the Fire Information Resource Evaluation System (FIRES), which are available at http://www.srd.alberta.ca/Wildfire.

We collected the historical forest wildfire data from 2006 to 2010 from the website. The dataset contained forest wildfire activities with area burned greater than or equal to 0.01 hectares (Figure 1(c) and 1(d)). The largest wildfire occurred in 2007 and had an area burned of $630 \mathrm{~km}^{2}$. To formulate hypothesis about the interaction between wildfire activities, we investigated the independence between the fire occurrences and area burned, in which we treated the spatial locations of wildfires as points and area burned as marks.

We used $T_{n}$ to assess the independence for each particular year and all years combined. To define the test statistic $T_{n}$, we chose $\mathcal{A}=\left\{A_{a}: A_{a}=\left\{\rho\left(\mathbf{s}, \mathbf{s}_{0}\right) \leq\right.\right.$ $a\}, a \geq 0\}$ and $\mathcal{B}=\left\{B_{b}: B_{b}=(-\infty, b]\right\}$, where $\rho\left(\mathbf{s}, \mathbf{s}_{0}\right)$ was the spherical distance between $\mathbf{s}$ and $\mathbf{s}_{0}$. We chose $\mathbf{s}_{0}$ at $60^{\circ}$ latitude north and $120^{\circ}$ longitude west because this point was at the upper-left corner of the whole study area. We derived the values $\sqrt{n} \hat{\varphi}_{n}\left(A_{a}, B_{b}\right)$ for all possible values of $a$ and $b$. The combined case is displayed in Figure 2. We computed the values and $p$-values of $T_{n}$ and displayed them in Table 3. It showed that the locations and area burned were not independent in all the cases.

To account for the interaction between points and marks, we used a model with intensity-dependent marks. The idea of intensity-dependent marks can be found in $[20,38]$. To apply the idea, we needed to derive an estimate of the intensity function of $N_{s}$. Therefore, we used a two-step method. In the

Table 3
Values $T_{n}$ for Alberta forest wildfire data, where all of the p-values of $T_{n}$ were almost 0

| Year | Total Area <br> burned (hectares) | Total No. <br> of Fires | $T_{n}$ |
| :---: | ---: | ---: | :---: |
| 2006 | 118762 | 1954 | 2.4168 |
| 2007 | 103669 | 1349 | 3.0343 |
| 2008 | 20787 | 1712 | 2.8652 |
| 2009 | 66947 | 1710 | 2.6096 |
| 2010 | 83656 | 1840 | 3.8439 |
| Combined | 393821 | 8565 | 6.4260 |

first step, we used a nonparametric kernel approach to estimate the intensity function $\lambda_{0}(s)$ of fire occurrences $[17,51]$. Let $\mathbf{s}_{i}$ be the location and $m_{i}$ be the area burned of the $i$-th fire given in the dataset. Based on the nonparametric kernel approach, the well-known Berman-Diggle estimator [3,13] was given by

$$
\begin{equation*}
\hat{\lambda}_{0}(\mathbf{s})=\frac{1}{\omega^{2} C_{\omega}(\mathbf{s})} \sum_{i=1}^{n} f\left(\rho\left(\mathbf{s}, \mathbf{s}_{i}\right) / \omega\right) \tag{5.1}
\end{equation*}
$$

where $\rho$ was the spherical distance function, $f$ was a two-dimensional kernel function, $\omega$ was the bandwidth, and $C_{\omega}(\mathbf{s})$ was the edge correction. We chose $f$ as the bivariate standard normal density function. Based on the $L_{2}$-norm cross validation criterion, we had the best $\omega$ was about $\hat{\omega}=77.42 \mathrm{~km}$. In the second step, we used a linear regression method to model area burned. According to the power-law in literature of forest fire research [31], we proposed to consider the following linear regression model

$$
\log m_{i}=\beta_{0}+\beta_{1} \log \hat{\lambda}_{0}\left(\mathbf{s}_{i}\right)+N\left(0, \sigma^{2}\right)
$$

We had $\hat{\beta}_{0}=-8.2845$ with standard error $0.1705, \hat{\beta}_{1}=-1.3825$ with standard error 0.0420 , and $\hat{\sigma}^{2}=2.307^{2}$. As $\beta_{1}$ was significantly different from 0 , we concluded there was a strong linear effect between fire occurrence and area burned. Because $\hat{\beta}_{1}<0$, our result implied that large forest fires were dominant if the relative frequency was small, which was consistent with the literature (e.g. [19, 33, 41, 47]).

## 6. Discussion

This research has proposed a Kolmogorov-Smirnov type method to test the independence between points and marks of marked point processes. The method primarily relies on the test statistic $T_{n}$ given by Equation (3.6). Two methods for $\mathcal{A}$ and $\mathcal{B}$ are recommended according to Theorems 3 and 4 as well as their corollaries. In the first method, the asymptotic distribution of $T_{n}$ depends on the unknown intensity functions of the marked point process but in the second method it does not. Therefore, we focus on the second method in this paper. The difference between the two methods is that: it is enough to conclude dependence if the null hypothesis is rejected by either methods but only the first method can accept the null hypothesis if the test is insignificant.

There are several possible extensions of our method. According to the random field model for marked point processes considered by [48], it would be interesting to study the mark-point intensity function and use it to model the local variation of the process. High dimensional nonparametric estimation methods may be used to estimate the mark-point intensity function, which may describe the interaction effect between points and marks. In the future, one could consider partitioning the mark-point intensity function into several important components so that they can represent pure point effects, pure mark effects, and point-mark interaction effects, respectively.

## Appendix A: Proofs

Proof of Proposition 1. The necessity is implied by Definition 1. According to the theorem of $\pi-\lambda$ system ([4], p. 42), it is enough to consider the proof for disjoint $A_{1}, \ldots, A_{n} \in \mathscr{S}$ and $B_{1}, \ldots, B_{n} \in \mathscr{M}$. The proof is obvious for a marked Poisson process as we have $\pi_{n}\left(A_{1} \times B_{1}, \ldots, A_{n} \times B_{n}\right)=\prod_{i=1}^{n} \pi\left(A_{i} \times B_{i}\right)$, $\pi_{s, n}\left(A_{1}, \ldots, A_{n}\right)=\prod_{i=1}^{n} \pi_{s}\left(A_{i}\right)$, and $\pi_{m, n}\left(B_{1}, \ldots, B_{n}\right)=\prod_{i=1}^{n} \pi_{m}\left(B_{i}\right)$ when $n \geq 1$.

Proof of Lemma 1. It is enough to consider the sufficiency because the necessity can be directly implied by Equation (3.2). If $\mathscr{S}$ and $\mathscr{M}$ can be generated by $\mathcal{A}$ and $\mathcal{B}$, respectively, then for any $\epsilon>0, C \in \mathscr{S}$, and $D \in \mathscr{M}$ there are $A \in \mathcal{A}$ and $B \in \mathcal{B}$ such that $\left|\pi_{s}(A)-\pi_{s}(C)\right| \leq \epsilon,\left|\pi_{m}(B)-\pi_{m}(D)\right| \leq \epsilon$, and $|\pi(A \times B)-\pi(C \times D)| \leq \epsilon([4]$, p. 169 $)$. Then,

$$
\begin{aligned}
\varphi(C, D)= & \pi(C \times D)-\pi_{s}(C) \pi_{m}(D) \\
= & {[\pi(C \times D)-\pi(A \times B)]+\left[\pi_{s}(A)-\pi_{s}(C)\right] \pi_{m}(B) } \\
& +\pi_{s}(C)\left[\pi_{m}(B)-\pi_{m}(D)\right]
\end{aligned}
$$

which implies $|\varphi(C, D)| \leq 3 \epsilon$. This immediately leads the conclusion by letting $\epsilon \rightarrow 0$.

Proof of Lemma 2. Since $n \sim \operatorname{Poisson}(\kappa), E\left[\hat{\pi}_{n}(A \times B) \mid n\right]=\pi(A \times B) I_{n>0}$ and $V\left[\hat{\pi}_{n}(A \times B) \mid n\right]=n^{-1} \pi(A \times B)[1-\pi(A \times B)] I_{n>0}$. Then, $E\left[\hat{\pi}_{n}(A \times B)\right]=$ $E\left[\pi(A \times B) I_{n>0}\right]=\pi(A \times B)\left(1-e^{-\kappa}\right) \rightarrow \pi(A \times B)$ and

$$
\begin{aligned}
& V\left[\hat{\pi}_{n}(A \times B)\right] \\
= & V\left[\pi(A \times B) I_{n>0}\right]+E\left\{n^{-1} \pi(A \times B)[1-\pi(A \times B)] I_{n>0}\right\} \\
= & \pi^{2}(A \times B) e^{-\kappa}\left(1-e^{-\kappa}\right)+2 \pi(A \times B)[1-\pi(A \times B)] \sum_{n=1}^{\infty} \frac{1}{2 n} \frac{\kappa^{n}}{n!} e^{-\kappa} \\
\leq & \pi^{2}(A \times B) e^{-\kappa}\left(1-e^{-\kappa}\right)+2 \pi(A \times B)[1-\pi(A \times B)] \sum_{n=1}^{\infty} \frac{\kappa^{n}}{(n+1)!} e^{-\kappa} \\
\leq & \pi^{2}(A \times B) e^{-\kappa}\left(1-e^{-\kappa}\right)+2 \pi(A \times B)[1-\pi(A \times B)] / \kappa \\
\rightarrow & 0
\end{aligned}
$$

as $\kappa \rightarrow \infty$. Using the Chebyshev inequality, we have $\hat{\pi}_{n}(A, B) \xrightarrow{P} \pi(A \times B)$. Similarly, we have $\hat{\pi}_{s, n}(A) \xrightarrow{P} \pi_{s}(A)$ and $\hat{\pi}_{m, n}(B) \xrightarrow{P} \pi_{m}(B)$. The final conclusion is drawn from the Continuous Mapping Theorem.

Proof of Lemma 3. For any $A_{i_{1}}, A_{i_{2}} \in \mathcal{A}$ and $B_{j_{1}}, B_{j_{2}} \in \mathcal{B}$, let $Y_{11}=N\left(A_{i_{1}} \times\right.$ $\left.B_{j_{1}}\right) / \kappa, Y_{12}=N\left(A_{i_{1}} \times \bar{B}_{j_{1}}\right) / \kappa, Y_{13}=N\left(\bar{A}_{i_{1}} \times B_{j_{1}}\right) / \kappa, Y_{14}=N\left(\bar{A}_{i_{1}} \times \bar{B}_{j_{1}}\right) / \kappa$, $Y_{21}=N\left(A_{i_{2}} \times B_{j_{2}}\right) / \kappa, Y_{22}=N\left(A_{i_{2}} \times \bar{B}_{j_{2}}\right) / \kappa, Y_{23}=N\left(\bar{A}_{i_{2}} \times B_{j_{2}}\right) / \kappa, Y_{24}=$ $N\left(\bar{A}_{i_{2}} \times \bar{B}_{j_{2}}\right) / \kappa$. Let $\mathbf{y}_{i}=\left(Y_{i 1}, Y_{i 2}, Y_{i 3}, Y_{i 4}\right)^{t}, \nu_{i}=\left(\nu_{i 1}, \nu_{i 2}, \nu_{i 3}, \nu_{i 4}\right)^{t}$ for $i=1,2$, $\mathbf{y}=\left(\mathbf{y}_{1}^{t}, \mathbf{y}_{2}^{t}\right)^{t}$, and $\nu=\left(\nu_{1}^{t}, \nu_{2}^{t}\right)^{t}$. Then,

$$
\sqrt{\kappa}(\mathbf{y}-\nu) \xrightarrow{D} N\left(0,\left(\begin{array}{cc}
\mathbf{V}_{1} & \mathbf{U}  \tag{A.1}\\
\mathbf{U}^{t} & \mathbf{V}_{2}
\end{array}\right)\right),
$$

where $\mathbf{V}_{1}=\operatorname{diag}\left(\nu_{1}\right), \mathbf{V}_{2}=\operatorname{diag}\left(\nu_{2}\right)$, and $\mathbf{U}=\left(u_{i j}\right)_{i, j=1,2,3,4}$. Denote $\mathbf{x}=$ $\left(\mathbf{x}_{1}^{t}, \mathbf{x}_{2}^{t}\right)^{t}$ with $\mathbf{x}_{1}=\left(x_{11}, x_{12}, x_{13}, x_{14}\right)^{t}$ and $\mathbf{x}_{2}=\left(x_{21}, x_{22}, x_{23}, x_{24}\right)^{t}$. Let $H(\mathbf{x})=$ $\left(h\left(\mathbf{x}_{1}\right), h\left(\mathbf{x}_{2}\right)\right)^{t}$, where

$$
\begin{equation*}
h(\mathbf{z})=\frac{z_{1}}{z_{1}+z_{2}+z_{3}+z_{4}}-\frac{\left(z_{1}+z_{2}\right)\left(z_{1}+z_{3}\right)}{\left(z_{1}+z_{2}+z_{3}+z_{4}\right)^{2}}, \mathbf{z}=\left(z_{1}, \ldots, z_{8}\right) \in \mathbb{R}^{8} \tag{A.2}
\end{equation*}
$$

If $n>0$, then $H(\mathbf{y})=\left(\hat{\varphi}_{n}\left(A_{i_{1}} \times B_{j_{1}}\right), \hat{\varphi}_{n}\left(A_{i_{2}} \times B_{j_{2}}\right)\right)^{t}$ and $H(\nu)=\left(\varphi\left(A_{i_{1}} \times\right.\right.$ $\left.\left.B_{j_{1}}\right), \varphi\left(A_{i_{2}} \times B_{j_{2}}\right)\right)^{t}$. Note that $\lim _{\kappa \rightarrow \infty} P(n>0)=1$ and $n / \kappa \xrightarrow{P} 1$, we can show the final conclusion using the $\Delta$-Method and Continuous Mapping Theorem.
Proof of Corollary 1. The conclusion can be directly implied from Lemma 3.

Proof of Theorem 1. Use the same notations as in Lemma 3. If the null hypothesis of Equation (3.3) is violated, then there are $\tilde{A} \in \mathcal{A}$ and $\tilde{B} \in \mathcal{B}$ such that $\varphi(\tilde{A}, \tilde{B})=\pi(\tilde{A}, \tilde{B})-\pi_{s}(\tilde{A}) \pi_{m}(\tilde{B}) \neq 0$. Without loss of generality, assume $\varphi(\tilde{A}, \tilde{B})$ is positive. Let $Z \sim \mathcal{N}\left(0, \sigma_{\tilde{A}, \tilde{B}}^{2}\right)$. Then for any $\epsilon>0$, there are $z_{0}>0$, $n_{0}>0$, and $\kappa_{0}>0$ such that $P\left(Z>-z_{0}\right) \geq 1-\epsilon / 3$, and for any $\kappa>\kappa_{0}$ there are $P\left(n \leq n_{0}\right) \leq \epsilon / 3$ and

$$
\left|P\left\{\sqrt{n}\left[\hat{\varphi}_{n}(\tilde{A}, \tilde{B})-\varphi(\tilde{A}, \tilde{B})\right]>-z_{0}\right\}-P\left(Z>-z_{0}\right)\right| \leq \epsilon / 3
$$

Then, we have

$$
P\left\{\sqrt{n}\left[\hat{\varphi}_{n}(\tilde{A}, \tilde{B})-\varphi(\tilde{A}, \tilde{B})\right]>-z_{0}\right\} \geq P\left(Z>-z_{0}\right)-\epsilon / 3 \geq 1-2 \epsilon / 3
$$

Note that the way to choose $n_{0}$ depends on $\kappa_{0}$ and $n_{0}$ can be arbitrary large if $\kappa_{0}$ becomes large. Therefore, we can increase $\kappa_{0}$ such that $\sqrt{n_{0}} \varphi(\tilde{A}, \tilde{B})-z_{0}>x$. Then

$$
\begin{aligned}
P\left\{T_{n}>x\right\} & =P\left\{\max _{A \in \mathcal{A}, B \in \mathcal{B}} \sqrt{n}\left|\hat{\varphi}_{n}(A, B)\right| \geq x\right\} \\
& \geq P\left\{\sqrt{n} \hat{\varphi}_{n}(\tilde{A}, \tilde{B})>x\right\} \\
& \geq P\left\{\sqrt{n}\left[\hat{\varphi}_{n}(\tilde{A}, \tilde{B})-\varphi(\tilde{A}, \tilde{B})\right]>-z_{0}, n>n_{0}\right\}
\end{aligned}
$$

$$
\begin{aligned}
& \quad+P\left\{\sqrt{n}\left[\hat{\varphi}_{n}(\tilde{A}, \tilde{B})-\varphi(\tilde{A}, \tilde{B})\right]>-z_{0}, n \leq n_{0}\right\}-P\left(n \leq n_{0}\right) \\
& \geq P\left\{\sqrt{n}\left[\hat{\varphi}_{n}(\tilde{A}, \tilde{B})-\varphi(\tilde{A}, \tilde{B})\right]>-z_{0}\right\}-\epsilon / 3 \\
& \geq \\
& \geq 1-\epsilon
\end{aligned}
$$

Therefore, $\lim _{\kappa \rightarrow \infty} P\left(T_{n}>x\right)=1$ for all $x>0$.
Proof of Corollary 2. The conclusion can be directly drawn using Lemma 1 and Theorem 1.

Proof of Theorem 2. We use the empirical process approach to show the conclusion of the theorem. Let $(S, M)$ be a single observation in the data. Then, $N$ can be understood as being independently obtained from the distribution of $(S, M)$ with total $n$ observations, where $n \sim \operatorname{Poisson}(\kappa)$. Let $F(x)=\pi\left(A_{\mathbf{a}} \times B_{\mathbf{b}}\right)=$ $P\left[(S, M) \in A_{\mathbf{a}} \times B_{\mathbf{b}}\right]$ for $\mathbf{x}=(\mathbf{a}, \mathbf{b})$ with $\mathbf{a} \in \mathbb{R}^{\alpha}$ and $\mathbf{b} \in \mathbb{R}^{\beta}$. Then, $F(\mathbf{x})$ is continuous in a and $\mathbf{b}$ because its density is $\lambda_{1}(m \mid \mathbf{s}) \lambda_{0}(\mathbf{s})$ according to Conditions (C2) and (C3) as well as Assumptions (A1)-(A3). Let $I_{\mathrm{a}}=\left\{\mathbf{s}: \mathbf{s} \in A_{\mathrm{a}}\right\}$ be the indicator function for the point in $A_{\mathbf{a}}$ and $I_{\mathbf{b}}=\left\{\mathbf{m}: \mathbf{m} \in B_{\mathbf{b}}\right\}$ be the indicator function for the mark in $B_{\mathbf{b}}$. Then, $\hat{F}(\mathbf{x})=\hat{\pi}_{n}\left(A_{\mathbf{a}} \times B_{\mathbf{b}}\right)$ is the empirical distribution of $F(x)$. Let $\mathcal{G}_{1}=\left\{I_{\mathbf{a}} I_{\mathbf{b}}: \mathbf{a} \in \mathbb{R}^{\alpha}, \mathbf{b} \in \mathbb{R}^{\beta}\right\}$, $\mathcal{G}_{2}=\left\{I_{\mathbf{a}}\left(1-I_{\mathbf{b}}\right): \mathbf{a} \in \mathbb{R}^{\alpha}, \mathbf{b} \in \mathbb{R}^{\beta}\right\}, \mathcal{G}_{3}=\left\{\left(1-I_{\mathbf{a}}\right) I_{\mathbf{b}}: \mathbf{a} \in \mathbb{R}^{\alpha}, \mathbf{b} \in \mathbb{R}^{\beta}\right\}$, and $\mathcal{G}_{4}=\left\{\left(1-I_{\mathbf{a}}\right)\left(1-I_{\mathbf{b}}\right): \mathbf{a} \in \mathbb{R}^{\alpha}, \mathbf{b} \in \mathbb{R}^{\beta}\right\}$. To apply the empirical process approach, it is enough to show $\mathcal{G}_{1}, \mathcal{G}_{2}, \mathcal{G}_{3}$, and $\mathcal{G}_{4}$ are $F$-Donsker, where the definition can be found on page 270 in [53]. In the following, we only display the proof for $\mathcal{G}_{1}$ because the proofs for $\mathcal{G}_{2}, \mathcal{G}_{3}$, and $\mathcal{G}_{4}$ are similar.

Let $F_{i}$ be the $i$-th marginal CDF of $F$. For any $\epsilon \in(0,1)$, we can find $J+2$ points denoted by $x_{0 i 0}, x_{0 i 1}, \ldots, x_{0 i J}, x_{0 i(J+1)}, i=1, \ldots, \alpha+\beta$, with $x_{0 i 0}=-\infty$ and $x_{0 i(J+1)}=\infty$ such that

$$
\frac{\epsilon^{2}}{\epsilon+\alpha+\beta} \leq F_{i}\left(x_{0 i(j+1)}\right)-F_{i}\left(x_{0 i j}\right) \leq \frac{\epsilon^{2}}{\alpha+\beta}, j=0, \ldots, J,
$$

with $(\alpha+\beta) / \epsilon^{2} \leq J \leq(\alpha+\beta) / \epsilon^{2}+1$. Let

$$
\mathcal{X}_{\epsilon}=\left\{\mathbf{x}=\left(x_{1}, \ldots, x_{\alpha+\beta}\right): x_{i}=x_{0 i j} \text { for some } j=0,1, \ldots, J+1\right\}
$$

Then $\# \mathcal{X}_{\epsilon}=(J+2)^{\alpha+\beta} \leq(\alpha+\beta+3)^{\alpha+\beta} / \epsilon^{2(\alpha+\beta)}$. For any $g_{y} \in \mathcal{G}_{1}$, we can find $\mathrm{x}^{\prime}, \mathrm{x}^{\prime \prime} \in \mathcal{X}_{\epsilon}$ such that $\mathrm{x}^{\prime} \preceq \mathbf{y} \preceq \mathrm{x}^{\prime \prime}$ but there is no $\mathrm{x}^{*} \in \mathcal{X}_{\epsilon}$ satisfying $x_{i}^{\prime}<x_{i}^{*}<x_{i}^{\prime \prime}$ for some $i=1, \ldots, \alpha+\beta$, where $x_{i}^{\prime}, x_{i}^{*}$, and $x_{i}^{\prime \prime}$ are the $i$-th component of $\mathbf{x}^{\prime}, \mathbf{x}^{*}$, and $\mathbf{x}^{\prime \prime}$, respectively. Then, $g_{\mathbf{x}^{\prime}} \leq g_{\mathbf{y}} \leq g_{\mathbf{x}^{\prime \prime}}$ and

$$
\left\|g_{\mathbf{x}^{\prime}}-g_{\mathbf{x}^{\prime \prime}}\right\|^{2}=\int_{\mathbb{R}^{\alpha+\beta}}\left|g_{\mathbf{x}^{\prime \prime}}(\mathbf{x})-g_{\mathbf{x}^{\prime}}(\mathbf{x})\right|^{2} F(d \mathbf{x}) \leq \sum_{i=1}^{\alpha+\beta}\left[F_{i}\left(x_{i}^{\prime \prime}\right)-F_{i}\left(x_{i}^{\prime}\right)\right] \leq \epsilon^{2}
$$

Because

$$
\int_{0}^{1} \sqrt{\log \left(\# \mathcal{X}_{\epsilon}\right)} d \epsilon \leq \int_{0}^{1} \sqrt{(\alpha+\beta) \log (\alpha+\beta+3)-2(\alpha+\beta) \log \epsilon} d \epsilon<\infty
$$

$\mathcal{G}_{1}$ is $F$-Donsker. Similarly, we can show $\mathcal{G}_{2}, \mathcal{G}_{3}$, and $\mathcal{G}_{4}$ are also $F$-Donsker. The conclusion is drawn using the method in proofs of Corollary 1 and Lemma 3
with the functional $\Delta$-method ([53], p. 291) because the gradient of $h$ defined in (A.2) is uniformly bounded.

Proof of Theorem 3. The conclusion can be directly implied by Lemma 3 and Theorem 2.
Proof of Corollary 3. Let $u(a)=\pi_{s}\left(A_{a}\right)$ and $v(b)=\pi_{m}\left(A_{b}\right)$ for $-\infty<a, b<$ $\infty$. Then, $\lim _{a \rightarrow-\infty} u(a)=\lim _{b \rightarrow-\infty} v(b)=0$ and $\lim _{a \rightarrow \infty} u(a)=\lim _{b \rightarrow \infty} v(b)=$ 1. Denote $u=u(a), u^{\prime}=u\left(a^{\prime}\right), v=v(b)$, and $v^{\prime}=v\left(b^{\prime}\right)$. Then, $\pi_{s}\left(A_{a} \cap A_{a^{\prime}}\right)-$ $\pi_{s}\left(A_{a}\right) \pi_{s}\left(A_{a^{\prime}}\right)=u \wedge u^{\prime}-u u^{\prime}$ and $\pi_{m}\left(B_{b} \cap B_{b^{\prime}}\right)-\pi_{m}\left(B_{a}\right) \pi_{s}\left(B_{a^{\prime}}\right)=v \wedge v^{\prime}-v v^{\prime}$. Thus, $\sigma_{A_{a}, B_{b}, A_{a^{\prime}}, B_{b^{\prime}}}=\left(u \wedge u^{\prime}-u u^{\prime}\right)\left(v \wedge v^{\prime}-v v^{\prime}\right)$, which is the covariance function of the two-dimensional standard Brownian pillow given by (2.4).

Proof of Theorem 4. Let $B_{1}, \ldots, B_{J}$ be a partition of $\mathcal{M}$. Conditioning on $N_{s, \kappa}(A)$ for a given Borel $A \in \mathcal{S}_{\kappa}, N_{\kappa}\left(A \times B_{1}\right), \ldots, N_{\kappa}\left(A \times B_{J}\right)$ are multinomially distributed with the total count equal to $N_{s, \kappa}(A)$ and the proportion vector equal to $\left(\pi_{m}\left(B_{1}\right), \ldots, \pi_{m}\left(B_{J}\right)\right)^{t}$. For any measurable $A, A^{\prime} \in \mathcal{S}_{1}=[-1,1]^{d}$, let $\mathbf{p}_{B, B^{\prime}}=\left(\pi_{m}\left(B \cap B^{\prime}\right), \pi_{m}\left(B \cap \bar{B}^{\prime}\right), \pi_{m}\left(\bar{B} \cap B^{\prime}\right), \pi_{m}\left(\bar{B} \cap \bar{B}^{\prime}\right)\right)^{t}$ and

$$
\begin{aligned}
\tilde{\mathbf{y}}_{\kappa, A, A^{\prime}, B, B^{\prime}}= & \left(N_{\kappa}\left[\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right) \times\left(B \cap B^{\prime}\right)\right], N_{\kappa}\left[\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right) \times\left(B \cap \bar{B}^{\prime}\right)\right],\right. \\
& \left.N_{\kappa}\left[\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right) \times\left(\bar{B} \cap B^{\prime}\right)\right], N_{\kappa}\left[\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right) \times\left(\bar{B} \cap \bar{B}^{\prime}\right)\right]\right)^{t},
\end{aligned}
$$

where $A_{\kappa}=\kappa A$ and $A_{\kappa}^{\prime}=\kappa A^{\prime}$. Then,

$$
\sqrt{N_{s, \kappa}\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right)}\left[\frac{\tilde{\mathbf{y}}_{\kappa, A, A^{\prime} B, B^{\prime}}}{N_{s, \kappa}\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right)}-\mathbf{p}_{B, B^{\prime}}\right] \xrightarrow{D} N\left[0, \operatorname{diag}\left(\mathbf{p}_{B, B^{\prime}}\right)-\mathbf{p}_{B, B^{\prime}} \mathbf{p}_{B, B^{\prime}}^{t}\right],
$$

as $\kappa \rightarrow \infty$, which is equivalent to

$$
\begin{aligned}
& \frac{n_{\kappa}}{\sqrt{N_{s, \kappa}\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right)}}\left[\frac{\tilde{\mathbf{y}}_{\kappa, A, A^{\prime}, B, B^{\prime}}}{n_{\kappa}}-\mathbf{p}_{B, B^{\prime}} \frac{N_{s, \kappa}\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right)}{n_{\kappa}}\right] \\
& \xrightarrow{D} N\left[0, \operatorname{diag}\left(\mathbf{p}_{B, B^{\prime}}\right)-\mathbf{p}_{B, B^{\prime}} \mathbf{p}_{B, B^{\prime}}^{t}\right] .
\end{aligned}
$$

Note that $N_{s, \kappa}\left(A_{\kappa} \cap A_{\kappa}^{\prime}\right) / n_{\kappa} \xrightarrow{P} \pi_{s}\left(A \cap A^{\prime}\right)$. We have

$$
\sqrt{n_{\kappa}}\left(\frac{\tilde{\mathbf{y}}_{\kappa, A, A^{\prime}, B, B^{\prime}}}{n_{\kappa}}-\tilde{\mu}_{A, A^{\prime} B, B^{\prime}}\right) \xrightarrow{D} N\left(0, \tilde{\Sigma}_{A, A^{\prime} B, B^{\prime}}\right),
$$

where $\tilde{\mu}_{A, A^{\prime}, B, B^{\prime}}=\mathbf{p}_{B, B^{\prime}} \pi_{s}\left(A \cap A^{\prime}\right)$ and $\tilde{\Sigma}_{A, A^{\prime}, B, B^{\prime}}=\pi_{s}\left(A \cap A^{\prime}\right)\left[\operatorname{diag}\left(\mathbf{p}_{B, B^{\prime}}\right)-\right.$ $\left.\mathbf{p}_{B, B^{\prime}} \mathbf{p}_{B, B^{\prime}}^{t}\right]$. Let

$$
\begin{gathered}
\mathbf{y}_{\kappa, A, A^{\prime}, B, B^{\prime}}=\left(\tilde{\mathbf{y}}_{\kappa, A, A^{\prime}, B, B^{\prime}}, \tilde{\mathbf{y}}_{\kappa, A, \bar{A}^{\prime}, B, B^{\prime}}, \tilde{\mathbf{y}}_{\left.\kappa, \bar{A}, A^{\prime}, B, B^{\prime}, \tilde{\mathbf{y}}_{\kappa, \bar{A}, \bar{A}^{\prime}, B, B^{\prime}}\right)^{t}}^{\mu_{A, A^{\prime}, B, B^{\prime}}=\left(\tilde{\mu}_{A, A^{\prime}, B, B^{\prime}}, \tilde{\mu}_{A, \bar{A}^{\prime}, B, B^{\prime}}, \tilde{\mu}_{\bar{A}, A^{\prime}, B, B^{\prime}}, \tilde{\mu}_{\bar{A}, \bar{A}^{\prime}, B, B^{\prime}}\right)^{t}} \mathrm{t}\right.
\end{gathered}
$$

and

$$
\Sigma_{A, A^{\prime} B, B^{\prime}}=\operatorname{diag}\left(\tilde{\Sigma}_{A, A^{\prime}, B, B^{\prime}}, \tilde{\Sigma}_{A, \bar{A}^{\prime}, B, B^{\prime}}, \tilde{\Sigma}_{\bar{A}, A^{\prime}, B, B^{\prime}}, \tilde{\Sigma}_{\bar{A}, \bar{A}^{\prime}, B, B^{\prime}}\right)
$$

From Theorem 4.2 in [23], we have

$$
\sqrt{n_{\kappa}}\left(\mathbf{y}_{\kappa, A, A^{\prime}, B, B^{\prime}}-\mu_{A, A^{\prime}, B, B^{\prime}}\right) \xrightarrow{D} N\left(0, \Sigma_{A, A^{\prime}, B, B^{\prime}}\right) .
$$

Let $H\left(z_{1}, \ldots, z_{16}\right)=\left(h_{1}\left(z_{1}, \ldots, z_{16}\right), h_{2}\left(z_{1}, \ldots, z_{16}\right)\right)^{t}$ with $h_{1}\left(z_{1}, \ldots, z_{16}\right)=$ $\left(z_{1}+z_{2}+z_{5}+z_{6}\right)-\left(z_{1}+z_{2}+z_{5}+z_{6}+z_{9}+z_{10}+z_{13}+z_{14}\right)\left(z_{1}+z_{2}+z_{3}+\right.$ $\left.z_{4}+z_{5}+z_{6}+z_{7}+z_{8}\right), h_{2}\left(z_{1}, \ldots, z_{16}\right)=\left(z_{1}+z_{3}+z_{5}+z_{7}\right)-\left(z_{1}+z_{3}+z_{5}+\right.$ $\left.z_{7}+z_{9}+z_{11}+z_{13}+z_{15}\right)\left(z_{1}+z_{2}+z_{3}+z_{4}+z_{9}+z_{10}+z_{11}+z_{12}\right)$. Then, $H\left(\mathbf{y}_{\kappa, A, A^{\prime}, B, B^{\prime}}\right)=\left(\hat{\varphi}(A, B), \hat{\varphi}\left(A^{\prime}, B^{\prime}\right)\right)^{t}, H\left(\mu_{A, A^{\prime} B, B^{\prime}}\right)=(0,0)^{t}$. Using the $\Delta$ theorem, we have

$$
\sqrt{n}\binom{\hat{\varphi}(A, B)}{\hat{\varphi}\left(A^{\prime}, B^{\prime}\right)} \xrightarrow{D} N\left[\binom{0}{0},\left(\begin{array}{cc}
\sigma_{A, B, A, B} & \sigma_{A, B, A^{\prime}, B^{\prime}} \\
\sigma_{A, B, A^{\prime}, B^{\prime}} & \sigma_{A^{\prime}, B^{\prime}, A^{\prime}, B^{\prime}}
\end{array}\right)\right],
$$

where $\sigma_{A, B, A^{\prime}, B^{\prime}}$ is given in Equation (3.11). We finally have the conclusion of the theorem using Corollary 7.2 in [23].

Proof of Corollary 4. The conclusion can be proven using exactly the same way displayed in proof of Corollary 3.

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