CONTROLLED MORRIS METHOD: A NEW DISTRIBUTION-FREE SEQUENTIAL TESTING PROCEDURE FOR FACTOR SCREENING

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ABSTRACT
Morris’s elementary effects method (MM) has been known as a model-free factor screening approach especially well-suited when the number of factors is large or when the computer model is computationally expensive to run. In this paper, we propose the controlled Morris method (CMM) that acts in a sequential manner to keep the computational effort down to a minimum. The sequential probability ratio test-based multiple testing procedure adopted by CMM enables to identify the factors with significant main and/or interaction effects while controlling Type I and Type II familywise error rates at desired levels. A numerical example is provided to demonstrate the efficacy and efficiency of CMM.

1 INTRODUCTION
Factor screening refers to the process of identifying, through design of experiments, statistical modeling and sampling, those factors that have a significant influence on the model output. Proposed for factor screening in the context of deterministic computer experiments by Morris (1991), Morris’s elementary effects method (MM) has been known as a model-free approach particularly well-suited when the number of factors is relatively large or when the computer model is computationally expensive to run. Recently, Campolongo and Braddock (1999) and Cropp and Braddock (2002) extend standard MM by providing estimates of two-factor interaction effects. Campolongo et al. (2007) propose to use normalized elementary effects as compared to those used in standard MM so that the performance of MM can be more robust. Boukouvalas et al. (2014) propose to implement MM in a sequential way so that factors having nonlinear effects can be identified more efficiently. Most recently, Fédou and Rendas (2015) present a fast mixed effects screening method that enables efficient estimation of the interaction graph of factors. Shi et al. (2016) propose an effective error control mechanism for controlling the overall false discovery rate achieved by MM, and reveal its connections with other screening methods such as sequential bifurcation (e.g., Bettonvil and Kleijnen 1997, Wan et al. 2010).

Despite the aforementioned improvements made to standard MM, little attention has been given to establishing an adaptive sampling procedure for MM with a rigorous statistical guarantee on its screening performance. In this paper, we propose the controlled Morris method (CMM) that adopts a novel distribution-free sequential probability ratio test (SPRT)-based multiple testing procedure for identifying factors that have significant main and/or interaction effects while ensuring the Type I and II familywise error rates controlled at desired levels.

While SPRT-based procedures have been proposed for multiple hypothesis testing, most of them rely on the assumption that the underlying distribution from which the observations are sampled is known (Wald 1992; De and Baron 2012b; De and Baron 2012a; Bartroff and Song 2014); among those distributions stipulated, Gaussian is particularly popular (Wan et al. 2010; Ankenman et al. 2014; Shi et al. 2014).
However, the distribution is typically unknown in practice, or even if known it cannot be specified by a simple distribution function. Despite the significant role played by nonparametric estimation methods in modern statistics, nonparametric SPRT-based hypothesis testing procedures have been rarely studied in the literature, to the best of our knowledge. Antoniak and Dillard (1968) and Yu and Su (2004) are among the few that have relaxed the distribution-known assumption, and they investigate Wilcoxon signed rank statistics but still have to assume that the underlying distribution is symmetric around the median. The SPRT-based sequential multiple testing procedure adopted by CMM, on the other hand, is fully distribution free, thanks to the use of online kernel density estimation.

The remainder of this paper is organized as follows. In Section 2, we give a brief review of the Morris’s elementary effects method. In Section 3, we provide details on the controlled Morris method (CMM). Section 4 provides a numerical evaluation of CMM. Section 5 concludes the paper.

2 A REVIEW ON MORRIS METHOD

The Morris’s elementary effects method (MM) is originally proposed for factor screening in the context of deterministic computer experiments (Morris 1991). Suppose that there are $k$ factors in total in the simulation model and each factor is scaled to take values from $[0, 1]$. For the purpose of factor screening, MM considers varying the value of each factor across $p$ pre-selected levels in $[0, 1]$; that is, the experimental region $\Omega$ for MM is a $k$-dimensional $p$-level grid in $[0, 1]^k$.

Let $Y(x)$ be the output obtained by running a deterministic computer experiment at factor combination $x = (x_1, \ldots, x_k) \in \Omega$. The elementary effect of the $j$th factor at $x$ is defined as

$$d_j(x) = \frac{Y(x + e_j \Delta) - Y(x)}{\Delta}, \quad j = 1, \ldots, k,$$

(1)

where $e_j$ denotes the unit vector in the direction of the $j$th axis; $\Delta$ is a predefined integer multiple of $1/(p-1)$ such that $x + e_j \Delta \in \Omega$. Hence, the $j$th factor $x_j$ assumes values in $\{0, 1/(p-1), 2/(p-1), \ldots, 1\}$. Intuitively speaking, $d_j(x)$ can be thought of as the partial derivative of $Y(x)$ with respect to $x_j$ when $\Delta$ is small (Woods and Lewis 2016).

The elementary effects corresponding to the $j$th factor, $d_j(x)$, follow a finite distribution, $F_j$, which can be obtained by randomly sampling the factor combination $x$ from $\Omega$. The number of elements of $F_j$ is then $p^{k-1}[p - \Delta(p-1)]$, where $p^{k-1}$ is the number of factor combinations formed by the remaining $k-1$ factors, and $p - \Delta(p-1)$ is the number of possible levels that factor $j$ can take to obtain elementary effects; for example, when $\Delta = 1/(p-1)$, factor $j$ can only have $p-1$ levels (i.e., the level corresponding to $x_j = 1$ is excluded). A recommended choice of $p$ is even and $\Delta = p(2(p-1))^{-1}$ (Morris 1991).

A highly centralized distribution $F_j$ suggests a consistent importance of factor $j$ across the experimental region $\Omega$, and a highly decentralized distribution indicates a strong dependence of factor $j$ on the other factors (i.e., nonlinear or interaction effects may be present). MM determines the importance of the $j$th factor in terms of two measures, $\mu_j$ and $\sigma_j$, respectively, the mean and standard deviation of $F_j$. To estimate these two measures, MM samples $N$ elementary effects from $F_j$ via a sampling design that generates $N$ trajectories in $\Omega$ (Morris 1991), and the $i$th trajectory provides $k$ elementary effects estimates $d_{ji}$ for $j = 1, 2, \ldots, k$. The following unbiased estimators of the mean and variance of $F_j$ are then used to assess the importance of factor $j$,

$$\hat{\mu}_j = \frac{1}{N} \sum_{i=1}^{N} d_{ji}, \quad \hat{\sigma}_j^2 = \frac{1}{N-1} \sum_{i=1}^{N} (d_{ji} - \hat{\mu}_j)^2,$$

(2)

(3)

where recall that $d_{ji}$ denotes the $i$th elementary effect randomly generated for factor $j$ from $F_j$. 

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In practice, Morris (1991) recommends to use a graph plotting $\hat{\mu}_j$ vs. $\hat{\sigma}_j$ with two lines corresponding to $\hat{\mu}_j = \pm 2\hat{\sigma}_j/\sqrt{N}$ for assessing the importance of factor $j$: If the point $(\hat{\mu}_j, \hat{\sigma}_j)$ locates outside the wedge formed by the two lines, then factor $j$ is deemed important. Such a practice, however, is more of a commonsense rule than a rigorously justified screening method.

To generate elementary effects from $F_j$, Morris introduces using sampling matrix and design matrix. For example,

$$B = \begin{bmatrix}
0 & 0 & \ldots & 0 \\
1 & 0 & \ldots & 0 \\
1 & 1 & \ldots & 0 \\
\vdots & \vdots & \ddots & 0 \\
1 & 1 & \ldots & 1 \\
\end{bmatrix}$$

is a $(k + 1) \times k$ sampling matrix, which consists of a $1 \times k$ vector of zeros in the first row and a $k \times k$ lower triangular matrix whose entries below the main diagonal are all ones. The design matrix corresponding to $B$ is $\Delta B := \Delta \times B$. Though a sampling matrix is easier to understand, its corresponding design matrix is the one actually used for running simulation experiments.

Notice that $B$ (or equivalently, $\Delta B$) can only generate one elementary effect for each factor, which holds true for any random form (also called random orientation or trajectory) of $\Delta B$, generically denoted by $B^\perp$. For example, when $k = 3$ and $p = 4$ (so that $\Delta = p (2 (p - 1))^{-1} = 2/3$), one random form of $\Delta B$ can be

$$B^\perp = \begin{bmatrix}
2/3 & 0 & 1 \\
2/3 & 2/3 & 1 \\
0 & 2/3 & 1 \\
0 & 2/3 & 1/3 \\
\end{bmatrix}.$$  

Morris (1991) provides a special algorithm to convert a sampling matrix $B$ to a random design matrix $B^\perp$ given specified $k$ and $p$. To obtain $N \geq 2$ independent elementary effects for each factor, Morris suggests to use $N$ random forms of $\Delta B$ and form the ultimate design matrix for running simulation experiments as $(B_1^\perp, \ldots, B_N^\perp)^\top$.

3 THE CONTROLLED MORRIS METHOD

In this section we describe the controlled Morris method (CMM) that acts in a sequential manner to keep the number of simulation runs down to a minimum, while identifying the factors that have significant main and/or interaction effects with Type I and Type II familywise error rates controlled at desired levels.

3.1 The Set-Up of CMM

Suppose that there are $k$ factors in a deterministic simulation model. Let $d_j = \{d_{j,1}, d_{j,2}, \ldots\}$ denote the set of elementary effects generated sequentially for factor $j$, where $d_{j,1}, d_{j,2}, \ldots$ are independent and identically distributed (i.i.d.) following a common distribution given by $F_j$ whose corresponding probability density function (PDF) is given by $f_j$, for $j = 1, 2, \ldots, k$. The goal is to determine whether each factor has a significant main or interaction effect or both, via simultaneously performing the following $2k$ individual hypothesis tests.

For a given factor $j$, we want to determine if it has a significant main effect using the following one-sided hypothesis test:

$$H_l : |\mu_j| \leq \Delta_{EE}^{(0)} \quad \text{vs.} \quad G_l : |\mu_j| \geq \Delta_{EE}^{(1)}, \quad \text{for} \quad l = 2j - 1; j = 1, 2, \ldots, k,$$  

$$H_l : |\mu_j| \leq \Delta_{EE}^{(0)} \quad \text{vs.} \quad G_l : |\mu_j| \geq \Delta_{EE}^{(1)},$$

$$\text{for} \quad l = 2j - 1; j = 1, 2, \ldots, k,$$  

3 THE CONTROLLED MORRIS METHOD
where the null and alternative are respectively denoted by $H_l$ and $G_l$; $\Delta_{EE}^{(0)}$ and $\Delta_{EE}^{(1)}$ are respectively the user specified parameters that define the thresholds of nonsignificant and significant main effects and they satisfy $0 \leq \Delta_{EE}^{(0)} \leq \Delta_{EE}^{(1)}$.

We note that (5) is equivalent to the following two alternative hypotheses tests depending on the sign of $\mu_j$:

$$
\begin{align*}
H_l: 0 \leq \mu_j \leq \Delta_{EE}^{(0)} & \quad \text{vs.} \quad G_l: \mu_j \geq \Delta_{EE}^{(1)}, \quad \text{if } \mu_j \geq 0; \\
H_l: -\Delta_{EE}^{(0)} \leq \mu_j \leq 0 & \quad \text{vs.} \quad G_l: \mu_j \leq -\Delta_{EE}^{(1)}, \quad \text{if } \mu_j < 0.
\end{align*}
$$

That is, when $\mu_j \geq 0$, we consider the main effect of factor $j$ to fall into one of the following two categories: (i) nonsignificant, if $\mu_j \leq \Delta_{EE}^{(0)}$; (ii) significant, if $\mu_j \geq \Delta_{EE}^{(1)}$; moreover, if $\mu_j \in (\Delta_{EE}^{(0)}, \Delta_{EE}^{(1)})$, then we want to provide a sound power to identify it as significant. Analogous description holds for the alternative case where $\mu_j < 0$. Notice that the sign of $\mu_j$ can be observed through that of $\bar{\mu}_j$ calculated with a sample of elementary effects generated for factor $j$ (Campolongo et al. 2007), and it may change as new elementary effects continue to arrive.

Similarly, we use the following one-sided hypothesis test to determine if factor $j$ has a significant interaction effect with other factors:

$$
H_l: \sigma_j \leq \Delta_{IE}^{(0)} \quad \text{vs.} \quad G_l: \sigma_j \geq \Delta_{IE}^{(1)}, \quad \text{for } l = 2j; j = 1, 2, \ldots, k,
$$

where $\Delta_{IE}^{(0)}$ and $\Delta_{IE}^{(1)}$ respectively denote the thresholds of nonsignificant and significant interaction effects that satisfy $0 \leq \Delta_{IE}^{(0)} \leq \Delta_{IE}^{(1)}$. The interaction effect of factor $j$ will be classified into the following two categories: (i) nonsignificant, if $\sigma_j \leq \Delta_{IE}^{(0)}$; and (ii) significant, if $\sigma_j \geq \Delta_{IE}^{(1)}$; moreover, if $\sigma_j \in (\Delta_{IE}^{(0)}, \Delta_{IE}^{(1)})$, then we want to provide a sound power to identify it as significant.

For notational convenience, let $\theta \triangleq (\theta_1, \theta_2, \ldots, \theta_{2k})^T$ be the vector $(\mu_1, \sigma_1, \mu_2, \sigma_2, \ldots, \mu_k, \sigma_k)^T$ whose dimensions are $2k \times 1$. That is, for $l \in \{1, 2, \ldots, 2k\}$, $\theta_l$ denotes the main effect (respectively, interaction effect) for factor $[l/2]$ if $l$ is odd (resp., even). Denote $\mathcal{H}(\theta) = \{l \in \{1, 2, \ldots, 2k\} : \theta_l \in H_l\}$ the set of indices whose corresponding null hypotheses are true; that is, factor $[l/2]$ has a truly nonsignificant main or interaction effect. Let $\mathcal{I}(\theta) = \{l \in \{1, 2, \ldots, 2k\} : \theta_l \in G_l\}$ be the set of indices whose corresponding null hypotheses are false; that is, factor $[l/2]$ has a truly significant main or interaction effect.

The aforementioned factor screening problem naturally falls into the multiple hypothesis test setting, where it is vital to control the error rates achieved especially when simultaneously testing a considerably large number of hypotheses (i.e., when $k$ is large). Upon a stopping rule $T$ is given which will be specified later, with decisions regarding acceptance or rejection of each of the $2k$ null hypotheses, CMM aims to control the \textit{Type I and Type II familywise error rates} (De and Baron 2012b, Bartroff and Song 2014) by guaranteeing that

$$
\begin{align*}
\text{FWE}_I(\theta) &= P\{H_l \text{ is rejected for some } l \in \mathcal{H}(\theta)\} \leq \alpha \\
\text{FWE}_II(\theta) &= P\{G_l \text{ is rejected for some } l \in \mathcal{I}(\theta)\} \leq \beta,
\end{align*}
$$

where $\alpha, \beta \in (0, 1)$ are two user-specified parameters, in addition to $\Delta_{EE}^{(0)}, \Delta_{EE}^{(1)}, \Delta_{IE}^{(0)},$ and $\Delta_{IE}^{(1)}$. Notice that the quantity $1 - \text{FWE}_II(\theta)$ is also known as “\textit{familywise power}”; equivalently, CMM aims to provide $\gamma \triangleq 1 - \beta$ familywise power for the entire factor screening procedure.

\subsection{3.2 Description of the CMM's Sequential Multiple Testing Procedure}

To effectively control the Type I and Type II familywise error rates simultaneously, CMM adopts a novel distribution-free sequential probability ratio test (SPRT)-based multiple testing procedure. Essentially, this procedure can be thought of as an ensemble of $2k$ individual SPRTs for identifying the significance of the main and interaction effects associated with each of the $k$ factors.

We now describe the procedure in terms of stages of sampling, between which accept/reject decisions are made to each hypothesis test. Without loss of generality, let $n$ denote the cumulative sample size of...
elementary effects collected for any active test (i.e., the $H_l$ for which no decision has been reached yet) that have up to and including the current stage. For the $l$th ($l = 1, 2, \ldots, 2k$) pair of hypotheses given by either (6) or (7), the testing of $H_l$ vs. $G_l$ is based on the following test statistic (De and Baron 2012a, De and Baron 2012b, Bartroff and Song 2014, Wang and Wan 2014),

$$A_l(n) = \sum_{j=1}^{n} \left( \log f_{i}(d_{j,i} | \theta_l \in G_l) - \log f_{i}(d_{j,i} | \theta_l \in H_l) \right), \quad \text{with } j = \lfloor \frac{l}{2} \rfloor,$$  \hfill (9)

where $f_{i}(\cdot | \theta_l \in H_l)$ and $f_{i}(\cdot | \theta_l \in G_l)$ are, respectively, the PDFs of elementary effects for factor $\lfloor \frac{l}{2} \rfloor$ given that $H_l$ and $G_l$ are true, for $\ell = 1, 2, \ldots, 2k$. Notice that $A_l(n)$ is known as Kullback-Leibler information numbers (De and Baron 2012b), and it is used for measuring the “distance” between two probability measures defined on a common measurable space (Dykstra 2005).

Let $\alpha_l$ and $\beta_l$ be the prescribed levels of the Type I and Type II error rates to achieve for the $l$th test according to Wald’s SPRT, and let $a_l$ and $b_l$ be the upper and lower stopping boundaries, respectively. Wald’s SPRT for the single $l$th hypothesis $H_l$ rejects it (i.e., chooses $G_l$) upon obtaining $n$ elementary effects if $A_l(n) \geq a_l$, accepts it (i.e., chooses $H_l$) if $A_l(n) \leq b_l$, and continues sampling elementary effects for factor $\lfloor \frac{l}{2} \rfloor$ if $A_l(n) \in (a_l, b_l)$.

The sequential multiple testing procedure of CMM starts with an initial sample size of $n_0$ elementary effects, and examines whether the resulting test statistic $A_l(n_0)$ crosses one of the two stopping boundaries (i.e., $a_l$ or $b_l$), which leads to the decision of declaring effect $l$ significant (respectively, the upper boundary is crossed) or nonsignificant (resp., the lower boundary is crossed). If neither one is crossed, then the procedure continues sampling one elementary effect for factor $\lfloor \frac{l}{2} \rfloor$ per stage until $A_l(n)$ escapes from the continue-sampling region by crossing one of the two boundaries (i.e., $A_l(n) \notin (b_l, a_l)$). Let $T_l$ be the stopping time of the $l$th hypothesis test, namely,

$$T_l = \inf \left\{ n \geq n_0 : A_l(n) \notin (b_l, a_l) \right\}. \hfill (10)$$

It is clear that the procedure continues sampling until all $2k$ tests reach decisions, and the stopping time of the entire procedure follows as

$$T = \inf \left\{ n \geq n_0 : \bigcap_{l=1}^{2k} \{ A_l(n) \notin (b_l, a_l) \} \right\}. \hfill (11)$$

Lemma 2 of De and Baron (2012b) suggests that the $l$th test can control the Type I and Type II error probabilities by using the upper and lower stopping boundaries given by $a_l$ and $b_l$ because

$$P\{H_l \text{ is rejected for some } l \in \mathcal{H}(\theta)\} \leq P\{A_l(n) \geq a_l \mid l \in \mathcal{H}(\theta)\} \leq e^{-\alpha_l},$$

$$P\{H_l \text{ is accepted for some } l \in \mathcal{G}(\theta)\} \leq P\{A_l(n) \leq b_l \mid l \in \mathcal{G}(\theta)\} \leq e^{\beta_l}. \hfill (12)$$

Therefore, by setting the two boundaries respectively as $a_l = -\log \alpha_l$ and $b_l = \log \beta_l$, we can control the Type I and Type II error probabilities conservatively at levels $\alpha_l = e^{-\alpha_l}$ and $\beta_l = e^{\beta_l}$ for $l = 1, 2, \ldots, 2k$. Furthermore, it follows from Theorem 1 of De and Baron (2012b) that the sequential multiple testing procedure with the stopping time given by (11) can control FWE$k(\theta)$ and $FWE_q(\theta)$ respectively at levels $\alpha$ and $\beta$, if we set $\alpha_l = \alpha / 2k$ and $\beta_l = \beta / 2k$ thanks to the Bonferroni’s inequality.

### 3.3 Online Kernel Density Estimation

In this section we propose to approximate $f_l(\cdot)$ used in (9) via online kernel density estimation. We then use the estimated density function $\hat{f}_l(\cdot)$ to derive a distribution-free SPRT test statistic according to (9), denoted by $\hat{A}_l(n)$, by replacing $f_l(\cdot)$ with the estimate $\hat{f}_l(\cdot)$.
As an indispensable nonparametric estimation tool, kernel density estimation (KDE) specializes in capturing the behaviors of the distribution of interest without imposing any parametric assumptions (Härdle 1990; Zhang and Wang 2014; Parpas et al. 2015); more importantly, it is more flexible to accommodate complex settings such as the underlying density is multimodal (Han et al. 2008). Though different KDE methods can be employed under the CMM framework, in this paper we restrict our attention to the online (sequential) kernel density estimation (oKDE) method proposed by Kristan et al. (2010) and Kristan et al. (2011). oKDE is a type of machine learning approach, which updates \( \hat{f}_l(\cdot) \) each time upon receiving a new observation so as to continuously provide accurate estimation from the data observed thus far. Notice that this feature suits the setting of sequential sampling for multiple hypothesis testing very well.

For the \( l \)th effect, we obtain an empirical probability density function (EPDF), denoted by \( \hat{f}_l \), based on the set of elementary effects for factor \( \lfloor l/2 \rfloor \) via oKDE, for \( l = 1, 2, \ldots, 2k \). For notational simplicity, we drop the subscript \( l \) in the following discussion, as this approach can be applied to any of the \( 2k \) effects. In short, oKDE provides an initial estimate \( \hat{f} \) based on a sample of \( n_0 \) elementary effects collected in the pilot stage, and proceeds following three steps listed below upon receiving a new elementary effect in each subsequent sampling stage: (1) update \( \hat{f} \) with the new observation; (2) re-estimate the optimal bandwidth used by oKDE; and (3) refine and compress the estimate \( \hat{f} \). The necessity of the last step will be explained shortly.

Specifically, using a sample of \( n \) elementary effects for a given factor, say, \( \{d_i, i = 1, 2, \ldots, n\} \), the EPDF \( \hat{f} \) can be given by an \( n \)-component Gaussian mixture model as follows,

\[
\hat{f}_n(d) = \sum_{i=1}^{n} \omega_i \mathcal{N}_h(d - d_i),
\]

where \( \omega_i \) denotes the weight of the \( i \)th component, and \( \mathcal{N}_h(d - d_i) \) is a Gaussian kernel, which is defined as

\[
\mathcal{N}_h(d - d_i) = \left(2\pi h^2\right)^\frac{-1}{2} \exp\left(-\frac{(d-d_i)^2}{2h^2}\right),
\]

and \( d_i \) and \( h \) denote, respectively, the center and bandwidth of the \( i \)th Gaussian kernel, for \( i = 1, 2, \ldots, n \).

With the \( n_0 \) elementary effects obtained in the pilot stage, the kernel density estimate consisting of \( n_0 \) evenly weighted kernels constructed using an equal bandwidth \( h_{n_0} \) can be expressed as

\[
\hat{f}_{n_0}(d) = \frac{1}{n_0} \sum_{i=1}^{n_0} \mathcal{N}_h(d - d_i) = \frac{1}{n_0} \sum_{i=1}^{n_0} (2\pi h_{n_0}^2)^{-\frac{1}{2}} \exp\left(-\frac{1}{2h_{n_0}^2}(d-d_i)^2\right)
\]

We note that due to the absence of sufficient information about \( f \) in the pilot stage, an equal weight and identical bandwidth are used in all the components of \( \hat{f}_{n_0}(d) \).

Without loss of generality, suppose that we have obtained \( n - 1 \) elementary effects \( (n \geq n_0 + 1) \) and derived the EPDF \( \hat{f}_{n-1}(d) \). Upon receiving an additional elementary effect \( d_n \), the EPDF can be updated according to the following expression,

\[
\hat{f}_n(d) = \left(1 - \frac{1}{n-1}\right) \hat{f}_{n-1}(d) + \frac{1}{n} \mathcal{N}_h(d - d_n),
\]

where the new observation is assigned a weight of \( n^{-1} \). It is obvious that the estimation accuracy achieved by \( \hat{f}_n(d) \) depends heavily on the choice of bandwidth \( h_n \). By minimizing the asymptotic mean integrated squared error (Kristan et al. 2010), the optimal bandwidth can be obtained as

\[
h_n^* = \left(\frac{1}{2n\sqrt{\pi}} \int f''(x)^2 dx\right)^{\frac{1}{2}},
\]
where $f''(\cdot)$ denotes the second derivative of $f(\cdot)$. As $n$ becomes large, $f(d)$ can be approximated well by $\hat{f}_{n-1}(d)$. Therefore, an estimator $\hat{h}_n^2$ of $h_n^2$ can be obtained by substituting $\hat{f}_{n-1}(d)$ into (17); and $\hat{f}_n(d)$ can be obtained from (16) by using $\hat{h}_n^2$ in place of $h_n^2$.

We observe from (16) that as the number of components in $\hat{f}_n(d)$ increases linearly with the sample size of elementary effects $n$, the computational complexity approximately scales as $O((n^2 - n)/2 + n)$, which is mainly attributed to the calculation of bandwidth given by (17) (Kristan et al. 2011). To alleviate this computational burden, a handful of methods have been proposed to reduce (or compress) the number of components in the estimator given (López-Rubio and de Lazcano-Lobato 2008; Deng et al. 2008; Kristan et al. 2010; Kristan et al. 2011). In this paper, we adopt the approach proposed by Kristan et al. (2011) to maintain a comparatively stationary model scale (i.e., use a compressed model). The underlying idea of compression is to identify $M$ ($\leq n$) clusters of similar components in the original kernel estimator, such that components belonging to the $m$th ($m = 1, 2, \ldots, M$) cluster can be well characterized by a single component centered at $\bar{d}_m$ in the compressed estimator. Specifically, this approach can help reduce the original $n$-component model given in (16) to the following one comprised of only $M$ components:

$$
\hat{f}(d) = \sum_{m=1}^{M} \hat{\alpha}_m \mathcal{X}_{\hat{h}_m}(d - \bar{d}_m),
$$

where the compressed parameters $\hat{\alpha}_m = \sum_{i \in \pi(m)} \alpha_i$, $\bar{d}_m = \hat{\alpha}_m^{-1} \sum_{i \in \pi(m)} \alpha_idx_i$ and $\hat{h}_m = \hat{\alpha}_m^{-1} \sum_{i \in \pi(m)} \alpha_i(h_i + d^2_i) - \bar{d}_m^2$, and $\pi(m)$ denotes the collection of $m$th disjoint (compressed) set of indices, for $m = 1, 2, \ldots, M$.

The key of implementing the compression is to seek an appropriate clustering allocation $\{\pi(m)\}_{m=1}^{M}$, together with the minimum number of clustering number $M$ to use. For the sake of brevity, we omit the details of the compression procedure and refer the interested reader to Algorithm 1 of Kristan et al. (2011).

### 3.4 Construction of Nonparametric Test Statistics

The task of obtaining a nonparametric estimator of $\Lambda_l(n)$ (given in (9)) to test the $l$th pair of hypotheses boils down to providing two separate EPDFs under the null and alternative hypotheses, namely, $\hat{f}_l(\cdot | \theta_l \in H_l)$ and $\hat{f}_l(\cdot | \theta_l \in G_l)$. We note that an estimator $\hat{f}_l$ directly obtained based on the original sample of elementary effects for factor $j$, $\mathbf{d}_j = \{d_{j:1}, d_{j:2}, \ldots\}$ with $j = [1/2]$, may be an appropriate estimator for neither $f_l(\cdot | \theta_l \in H_l)$ nor $f_l(\cdot | \theta_l \in G_l)$.

Inspired by Shi et al. (2016), we transform the original sample of elementary effects $\mathbf{d}_j$ to a sample that complies with a parameter setting that is consistent with a given hypothesis of the $l$th test. For testing either the main or interaction effect, we apply two transformations to $\mathbf{d}_j$, such that the transformed samples, $\bar{d}_l^{(s)}$ for $s = 0$ and 1, complies with the $H_l$ (respectively, $G_l$) when $s = 0$ (resp., $s = 1$). We note that the superscripts (0) and (1) correspond to the null and alternative hypotheses, respectively.

When testing the significance of the main effect of factor $j$ via the $l$th hypothesis test with $l = 2j - 1$ for $j = 1, 2, \ldots, k$, the following transformation equation is applied to $\mathbf{d}_j$ to obtain $\bar{d}_l^{(s)}$ for $s = 0$ and 1:

$$
d_{l;i}^{(s)} = d_{j:i} - \bar{\mu}_j(n) + \Delta_{\text{EE}}^{(s)}, \quad i = 1, 2, \ldots, n, \tag{19}
$$

where $n$ denotes the sample size of $\mathbf{d}_j$ and $\bar{\mu}_j(n)$ denotes the sample mean of the $n$ original elementary effects obtained for factor $j$; and $\Delta_{\text{EE}}^{(s)}$ is as defined in (5). It is clear that the transformed sample $\bar{d}_l^{(s)} = (d_{l;1}^{(s)}, d_{l;2}^{(s)}, \ldots, d_{l;n}^{(s)})$ has its mean equal to $\Delta_{\text{EE}}^{(s)}$ for $s = 0, 1$ and its variance being fixed at $\sigma^2_f$.

When testing the significance of the interaction effect of factor $j$ via the $l$th hypothesis test with $l = 2j$ for $j = 1, 2, \ldots, k$, the following transformation equation is applied to obtain $\bar{d}_l^{(s)}$ for $s = 0$ and 1:

$$
d_{l;i}^{(s)} = \frac{d_{j:i}}{\sigma_j(n)} \Delta_{\text{EE}}^{(s)}, \quad i = 1, 2, \ldots, n, \tag{20}
$$
where $\hat{\sigma}_j(n)$ denotes the sample standard deviation of the elementary effects in $d_j$, and $\Delta_{\text{IE}}^{(s)}$ is as defined in (7). We see that the transformed sample $\tilde{d}_t^{(s)}$ has its mean equal to zero with its variance equal to $\Delta_{\text{IE}}^{(s)}$ for $s = 0, 1$.

With the transformed sample $\tilde{d}_t^{(s)} = (\tilde{d}_{t,1}^{(s)}, \ldots, \tilde{d}_{n,p}^{(s)})$ for the $l$th effect, we obtain an EPDF using (18), say $\tilde{f}_i(\cdot|\tilde{d}_t^{(s)})$ for $s = 0, 1$. It is clear that $\tilde{f}_i(\cdot|\tilde{d}_t^{(s)})$ satisfies the corresponding hypothesis of interest specified in either (6) or (7). Finally, the kernel-based SPRT test statistic for the $l$th test can be constructed as follows:

$$\hat{\Lambda}_l(n) = \sum_{i=1}^{n} \left( \log \tilde{f}_i(d_{j,i}|\tilde{d}_t^{(1)}) - \log \tilde{f}_i(d_{j,i}|\tilde{d}_t^{(0)}) \right).$$

(21)

Compared to the original test statistic $\Lambda_l(n)$ that is available only if the exact distribution is known, the calculation of $\hat{\Lambda}_l(n)$ only requires a sample of elementary effects for factor $\lfloor l/2 \rfloor$ to obtain two kernel-based EPDFs that correspond to $H_i$ and $G_i$, respectively.

### 3.5 Specification of Stopping Rules

In this paper we adopt a common pair of stopping boundaries for testing the $2k$ effects, which can be regarded as a type of one-shot boundary stopping strategy. Specifically, let $a_l \equiv a$ and $b_l \equiv b$ be the upper and lower stopping boundaries for $l = 1, 2, \ldots, 2k$. The stopping time of the $l$th hypothesis test in (10) reduces to

$$T_l = \inf \left\{ n \geq n_0 : \hat{\Lambda}_l(n) \notin (b, a) \right\},$$

(22)

where $\hat{\Lambda}_l(n)$ is as given in (21). Subsequently, the stopping time of the entire procedure becomes

$$T = \inf \left\{ n \geq n_0 : \bigcap_{l=1}^{2k} \{ \hat{\Lambda}_l(n) \notin (b, a) \} \right\}.$$

(23)

Following the discussion given in Subsection 3.2, an immediate choice for the upper and lower boundaries can be given as

$$a = -\log \left( \frac{\alpha}{2k} \right) \quad \text{and} \quad b = \log \left( \frac{\beta}{2k} \right), \quad l = 1, 2, \ldots, 2k,$$

(24)

where we set $\alpha_l = \alpha/2k$ and $\beta_l = \beta/2k$ for $l = 1, 2, \ldots, 2k$, and $\alpha$ and $\beta$ are respectively the desired Type I and Type II familywise error levels specified in (8). It is easy to see that the two boundaries given by (24) are two horizontal lines in parallel with the horizontal axis that denotes the sample size $n$. The greater the values of $\alpha$ and $\beta$ are, the narrower the continue-sampling region of the $l$th test becomes and the faster the test terminates.

### 4 NUMERICAL EVALUATION

In this section we demonstrate the performance of CMM on a factor screening problem, which has also been studied by Morris (1991) and Pujo (2009). The simulation output at a given factor combination $x = (x_1, x_2, \ldots, x_{20})^\top \in [0, 1]^{20}$ is generated by

$$y = \beta_0 + \sum_{j=1}^{20} \beta_j w_j + \sum_{i<j}^{20} \beta_{i,j} w_i w_j + \sum_{\ell < i < j} \beta_{\ell,i,j} w_{\ell} w_i w_j + \sum_{s < \ell < i < j} \beta_{s,\ell,i,j} w_s w_{\ell} w_i w_j,$$

(25)

where $w_j \in [-1, 1]$ is transformed from $x_j \in [0, 1]$, according to the following two transformations: (1) the linear transformation $w_j = 2(x_j - 0.5)$, and (2) the nonlinear transformation $w_j = 2(1.1x_j/(x_j + 0.1) - 0.5)$, for $j = 1, 2, \ldots, 20$. The nonlinear transformation is applied to factors $j \in \{3, 5, 7\}$, and the linear
transformation is applied to the remaining 17 factors. We note that the value of $w_j$ is uniformly distributed in $[-1, 1]$ if the linear transformation is applied to $x_j$, whereas it is more likely to be in $[0, 1]$ if the nonlinear transformation is applied. The coefficients in (25) are specified as follows. Regarding the first and second-order coefficients, $\beta_j = 20$ for $j \in \{1, 2, \ldots, 10\}$ and $\beta_{i,j} = 15$ for $i,j \in \{1,2,\ldots,6\}$; the remaining $\beta_j$'s and $\beta_{i,j}$'s are independently sampled from standard normal distribution $\mathcal{N}(0,1)$. With respect to the third and fourth-order coefficients, $\beta_{\ell,i,j} = -10$ for $\ell,i,j \in \{1,2,\ldots,5\}$ and $\beta_{s,\ell,i,j} = 5$ for $s,\ell,i,j \in \{1,2,3,4\}$; the remaining $\beta_{\ell,i,j}$'s, $\beta_{s,\ell,i,j}$'s and $\beta_0$ are all set to zeros.

We assess the efficacy and efficiency of CMM for detecting important main and interaction effects associated with the 20 factors. To implement the sequential procedure of CMM, we use an initial sample size $n_0 = 20$ for all factors, and set the target Type I and Type II familywise error levels respectively at $\alpha = \beta = 0.1$. The threshold parameters $\Delta_{EE}^{(0)}$, $\Delta_{EE}^{(1)}$, $\Delta_{IE}^{(0)}$, and $\Delta_{IE}^{(1)}$ are set to 20, 30, 40, and 60, respectively. The entire procedure of CMM is applied for 1000 independent macro-replications, and the efficacy of CMM is evaluated by the fraction of times a given effect is declared significant, denoted by $\widehat{P}\{\text{DI}\}$, which is a commonly used measure in factor screening (Wan et al. 2010; Shi et al. 2014; Shi et al. 2016): 

$$\widehat{P}\{\text{DI}\} = \frac{\#(\text{an effect is declared significant})}{\#\text{macro-replications}}.$$ 

The $\widehat{P}\{\text{DI}\}$'s obtained by CMM for testing the main and interaction effects of all 20 factors are shown in Figure 1. For each factor, two vertical bars are given that show the values of $\widehat{P}\{\text{DI}\}$ obtained for testing the significance of its corresponding main and interaction effects. We note that at the bottom of Figure 1 the 20 factors are also classified into 7 different groups “A–G” according to the coefficients associated with the terms involving each factor given in (25) and the transformation equations applied; see Shi et al. (2016) for details on how the groups are derived. The following observations can be made from Figure 1. First, the smaller $\mu_j$ (respectively, $\sigma_j$) is as compared to $\Delta_{EE}^{(0)}$ (resp. $\Delta_{IE}^{(0)}$), the closer the resulting $\widehat{P}\{\text{DI}\}$ for factor

![Figure 1: The resulting $\widehat{P}\{\text{DI}\}$ obtained by CMM for testing significance of the main and interaction effects of the 20 factors.](image-url)
Figure 2: Comparison of average sample sizes used by CMM for testing significance of the main and interaction effects of the 20 factors.

$j$ is to zero. Second, the greater $\mu_j$ (respectively, $\sigma_j$) is when compared to $\Delta_{EE}^{(1)}$ (resp. $\Delta_{IE}^{(1)}$), the closer the resulting $\hat{P}\{\text{DI}\}$ for factor $j$ gets to one. Third, for factors with $\mu_j \in (\Delta_{EE}^{(0)}, \Delta_{EE}^{(1)})$ or $\sigma_j \in (\Delta_{IE}^{(0)}, \Delta_{IE}^{(1)})$, the resulting $\hat{P}\{\text{DI}\}$ takes reasonable values in $[0, 1]$. Lastly, we note that CMM obtains almost identical $\hat{P}\{\text{DI}\}$ for factors within the same group. Therefore, we conclude that CMM provides a desired statistical performance guarantee for testing the significance of the main and interaction effects of each individual factor.

The computational efficiency of CMM is quantified by the average sample size of elementary effects used by the sequential procedure for testing a given effect across the 1000 macro-replications. The respective sample sizes used for testing the significance of the main and interaction effects of each factor are shown in Figure 2. We observe that CMM adapts the sample sizes used for the 20 factors according to the magnitudes of their respective main and interaction effects. In particular, for those factors with a high $\sigma_j$ and $\mu_j$ being close to $\Delta_{EE}^{(0)}$ or $\Delta_{EE}^{(1)}$, CMM typically uses a large sample size for testing the significance of the corresponding factor. It is intuitively clear that such an adaptive sampling strategy of CMM is more efficient than the default equal budget allocation rule adopted by standard MM.

5 CONCLUSIONS

In this paper, we propose the controlled Morris method (CMM) for factor screening that acts in a sequential manner to keep the computational effort down to a minimum. The SPRT-based multiple testing procedure adopted enables CMM to identify the factors with significant main and/or interaction effects while controlling Type I and Type II familywise error rates at desired levels. Though CMM is proposed in the context of factor screening, its distribution-free SPRT-based multiple testing procedure can be broadly applied to various settings beyond factor screening. Future research topics include extending CMM for factor screening in the stochastic simulation setting and enhancing the computational efficiency achieved by the sequential multiple testing procedure.
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