

NON PARAMETRIC METHODS FOR GENOMIC INFERENCE

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Large-scale statistical analysis of data sets associated with genome sequences plays an important role in modern biology. A key component of such statistical analyses is the computation of p-values and confidence bounds for statistics defined on the genome. Currently such computation is commonly achieved through ad hoc simulation measures. The method of randomization, which is at the heart of these simulation procedures, can significantly affect the resulting statistical conclusions. Most simulation schemes introduce a variety of hidden assumptions regarding the nature of the randomness in the data, resulting in a failure to capture biologically meaningful relationships. To address the need for a method of assessing the significance of observations within large scale genomic studies, where there often exists a complex dependency structure between observations, we propose a unified solution built upon a data subsampling approach. We propose a piecewise stationary model for genome sequences and show that the subsampling approach gives correct answers under this model. We illustrate the method on three simulation studies and two real data examples.

1. Introduction.

1.1. *Background.* This paper grew out of a number of examples arising in data coming from the ENCODE Pilot Project (Birney et al. (2007)), which is composed of multiple, diverse experiments performed on a targeted 1% of the human genome. Computational analyses of this data are aimed at revealing new insights about how the information coded in the DNA blueprint is turned into functioning systems in the living cell. Variations of some of the methods described here have been applied at various places in that paper, as well as in Margulies et al. (2007), for assessing significance and computing confidence bounds for statistics that operate along a genomic sequence. The background of these methods is described in cookbook form

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in the supplements to those papers, and it is the goal of this paper to present them rigorously and to develop some necessary refinements.

Essentially, we will argue that, in making inference about statistics computed from “large” stretches of the genome, in the absence of real knowledge about the evolutionary path which led to the genome in question, the best we can do is to model the genome by a piecewise stationary ergodic random process. The variables of this process can be base pair composition or some other local features, such as various annotated functional elements.

In the purely stationary case some of the types of questions that we will address, such as tests for independence of point processes, confidence bounds for expectations of local functions, goodness of fit of models, have been considered extensively. However, inference for piecewise stationary models appears not to have been investigated. With the advent of enormous amounts of genomic data all sorts of inferential questions have arisen. The proposed model may be the only truly nonparametric approach to the genome, although just as in ordinary nonparametric statistics there are many possible ways of carrying out inference.

Our methods are based on a development of the resampling schemes of Politis and Romano (1994), Politis, Romano, and Wolf (1999) and the block bootstrap methods of Künsch (1989). As we shall see, in many situations, Gaussian approximations can replace these schemes. But in these situations, as with the ordinary bootstrap, we believe that a subsampling approach is valuable for the following reasons:

- Letting the computer do the approximation is much easier.
- Some statistics, such as tests of the Kolmogorov Smirnov type, are functions of stochastic processes to which a joint Gaussian approximation applies. Then, limiting distributions can only be computed by simulation.
- The bootstrap distributions of our statistics show us whether the approximate Gaussianity we have invoked for the “true” distribution of these statistics is in fact warranted. This visual confirmation is invaluable.

This paper is organized as follows. We begin with some concrete examples from the ENCODE data as well as other types of genomic data in Section 1.2, and proceed with a motivated description of our model in Section 2. Our methods are discussed both qualitatively and mathematically in Sections 3 and 4. Section 5 contains results from simulation studies and real data analysis. Finally, an appendix with proofs of theorems stated in Sections 3 and 4 completes the paper.

The statistics and methods discussed in this paper have been implemented in several computing languages and are available for download at <http://encode.dyndns.org/>. Each of these implementations runs in $n\log(n)$ time, where n is the number of instances of the more frequent feature. On a desktop PC (Intel Core Duo 3Ghz and 2Gb RAM) the Python version takes over 1000 samples per second for features on the order of 10^4 instances.

1.2. *Motivating Examples.* We start with several fundamental questions that arise in genomic studies.

- **Association of functional elements in genomes.** In genomic analyses, a natural problem of interest is the association among the different functional sites/features annotated along the DNA sequence. Its biological motivation comes from the common belief that significant physical overlapping of functional sites in the genome suggests biological constraints or relationships. In the ENCODE project, to understand the possible functional roles of the evolutionarily constrained sequences that are conserved across multiple species, overlap between the constrained sequences and several experimental annotations, such as the 5'UTR, RxFragments, pseudogenes and coding sequences (CDSs), have been evaluated using the method discussed in this paper. It was found that the overlap of most experimental annotations with the constrained sequences are significantly different from random (Birney et al., 2007). An illustrative example from The ENCODE Project (Birney et al., 2007) is detailed in Section 5.1.
- **Cooperativity between transcription factor binding sites.** In some situations, it is interesting to study the associations between neighboring functional sites that do not necessarily overlap. For instance, it is known that transcription factors often work cooperatively and their binding sites (TFBS) tend to occur in clusters. Consequently, an effective method to identify interacting transcription factors has been to evaluate the significance of co-occurrences of their binding sites in a local genomic region (Kang et al. (2009); Yu et al. (2004)). This study has the same formulation as the above ENCODE examples given a functional site defined as follows: for a TFBS of length l at position i , we define the region $(i - m, i + l + m)$ as a functional site. Then two overlapping functional sites is equivalent to two neighboring TFBSs with interdistance less than $2m$, and the methods discussed in this paper for evaluating the significance of overlapping functional features can be applied, though we leave this application which involves considering statistics of the K-S type to a later paper.

- **Correlating DNA copy number with genomic content.** Recent technology has made it possible to assay DNA copy number variation at a very fine scale along the genome (For review, see Carter (2007)). Many studies, e.g. Redon et al. (2006), have shown that such variation in DNA copy number is a common type of polymorphism in the human genome. To what extent do these regions of copy number changes overlap with known genomic features, such as coding sequences? Redon et al. performed such an analysis and argued that copy number variant regions have a significant paucity for coding regions. The p-values supporting this claim were based on random permutations of the start locations of the variant segments. This assumes uniformity and stationarity of the copy number variants. The methods discussed in this paper for evaluating the significance of overlapping features, which assumes neither uniformity nor stationarity, can again be applied to this problem. Actually, the results from our method suggest a different conclusion on this problem (see section 5.5).

As we have seen in these examples, a common question asked in many applications is the following: Given the position vectors of two features in the genome, e.g. “conservation between species” and “transcription start sites”, and a measure of relatedness between features, e.g. base or region percentage overlap; how significant is the observed value of the measure? How does it compare with that which might be observed “at random”?

The essential challenge in the statistical formulation of this problem is the appropriate modeling of randomness of the genome, since we observe only one of the multitudes of possible genomes that evolution might have produced for our and other species.

How have such questions been answered previously? Existing methods employ varied ways to simulate the locations of features within genomes, but all center around the uniformity assumption of the features’ start positions: The features must occur homogeneously in the studied genome region, e.g. Blakesley et al (2004) and Redon et al. (2006). This assumption ignores the natural clumping of features as well as the non-stationarity of genome sequences. Clumping of features is quite common along the genome due to either the feature’s own characteristic, e.g. transcription factor binding sites (TFBSs) tend to occur in clusters, or the genome’s evolutionary constraints, e.g. conserved elements are often found in dense conservation neighborhoods. Ignoring these natural properties could result in misleading conclusions.

In this paper, we suggest a piecewise stationary model for the genome (see section 2), and based on it, propose a method to infer the relationships between features which we view as “nonparametric” as possible (see sections

4.2 and 4.4). The model is based on assumptions which we demonstrate in real data examples in Section 1.2 to be plausible.

2. The Piecewise Stationary Model.

2.1. *Genomic motivation.* We postulate the following for the observed genomes or genomic regions:

- They can be thought of as a concatenation of a number of regions, each of which is homogenous in a way we describe below.
- Features that are located very far from each other on the average have little to do with each other.
- The number of such homogeneous regions is small compared to the total length of the observed genomes that we consider.

These assumptions, which form the underpinning of our *block stationary model* for genomic features, are motivated by earlier studies of DNA sequences, which show that there are global shifts in base composition, but that certain sequence characteristics are locally un-changing. One such characteristic is the GC content. Bernardi et al. (1985) coined the term “isochore” to denote large segments (of length greater than 300 Kb) that have fairly homogeneous base composition, and especially, constant GC composition. Even earlier, evidence of segmental DNA structure can be found in chromosomal banding in polytene chromosomes in drosophila, visible through the microscope, that result from underlying physical and chemical structure. These banding patterns are stable enough to be used for the identification of chromosomes and for genetic mapping, and are physical evidence for a block stationarity model for the GC content of the genome.

The experimental evidence for segmental genome structure and the increasing availability of DNA sequence data have inspired attempts to computationally segment DNA into statistically homogeneous regions. The paper by Braun and Müller (1998) offers a review of statistical methods developed for detecting and modeling the inhomogeneity in DNA sequences. There have been many attempts to segment DNA sequences by both base composition (Fu and Curnow (1990), Churchill (1989,1992), Li et al (2002)) and chemical characteristics (Li et al. (1998)). Most of these computational studies concluded that a model that assumes block-wise stationarity gives a significantly better fit to the data than stationary models (see, for example, the conclusions of two very different studies by Fickett, Torney, and Wolf (1992) and Li et al. (1998)).

A subtle issue in the definition of “homogeneity” is the scale at which the genome is being analyzed. Inhomogeneity at the kilobase resolution, for

example, might be “smoothed out” in an analysis at the megabase level. The level of resolution is a modeling issue that must be considered carefully with the goal of the analysis in mind.

Implicit in our formulation is an “ergodic” hypothesis. We want probabilities to refer to the population of potential genomes. We assume that the statistics of the genome we have mimic those of the population of genomes. This is entirely analogous to the ergodic hypothesis that long term time averages agree with space averages for trajectories of dynamic systems.

2.2. Mathematical formulation. In mathematical terms, the block stationarity model assumes that we observe a sequence of random variables $\{X_1, \dots, X_n\}$ positioned linearly along the genomic region of interest. $X_k, k = 1, \dots, n$, may be base composition, or some other measurable feature. We assume that there exist integers $\tau = \tau^{(n)} = (\tau_0, \dots, \tau_U)$, where $0 = \tau_0 < \tau_1 < \dots < \tau_U = n$, such that the collections of variables, $\{X_{\tau_i}, \dots, X_{\tau_{i+1}}\}$ are separately stationary for each $i = 0, \dots, U - 1$. We let $n_i = \tau_i - \tau_{i-1}$ be the length of the i -th region, and let there be U such regions in total. For convenience, we introduce the mapping

$$\pi : \{1, \dots, n\} \rightarrow \{(i, j) : 1 \leq i \leq U, 1 \leq j \leq n_i\}$$

which relates the relabeled sequence, $\{X_{ij} : 1 \leq i \leq U, 1 \leq j \leq n_i\}$ to the original sequence $\{X_1, \dots, X_n\}$. We write $\pi = (\pi_1, \pi_2)$ with $\pi(k) = (i, j)$ if and only if $k = \tau_i + j$. We will use the notations X_{ij} and X_k interchangeably throughout this paper.

For any k_1, k_2 , let $\mathcal{F}_{k_1}^{k_2}$ be the σ -field generated by X_{k_1}, \dots, X_{k_2} . Define $m(k)$ to be the standard Rosenblatt mixing number (Dedecker et al, 2007),

$$m(k) = \sup\{|\mathbb{P}(AB) - \mathbb{P}(A)\mathbb{P}(B)| : A \in \mathcal{F}_1^l, B \in \mathcal{F}_{l+k}^n, 1 \leq l \leq n - k\}.$$

Then, assumptions 1-3 stated at the beginning of this section translate to the following:

- A1. The sequence $\{X_1, \dots, X_n\}$ is piecewise stationary. That is, $\{X_{ij} : 1 \leq j \leq n_i\}$ is a stationary sequence for $i = 1, \dots, U$.
- A2. There exists constants c and $\beta > 0$ such that $m(k) \leq ck^{-\beta}$ for all k .
- A3. $U/n \rightarrow 0$.

An immediate and important consequence of A1-A3 is that for any fixed small k , if we define $W_1 = (X_1, \dots, X_k), W_2 = (X_{k+1}, \dots, X_{2k}), \dots, W_m = (X_{n-k+1}, \dots, X_n)$, where $m = n/k$, then $\{W_1, \dots, W_m\}$ also obey A1-A3. This is useful, for example, in the region overlap example considered in the next section.

The remarkable feature of these assumptions, which are more general to our knowledge than any made heretofore in this context, is that they still allow us to conduct most of the statistical inference of interest. Not surprisingly, these assumptions lead to more conservative estimates of significance than any of the previous methods.

3. Linear Statistics and Gaussian Approximation. We study the distribution of a class of linear statistics of interest under the above piecewise stationary model. As an illustration, we consider the ENCODE data examples, and suppose that we are interested in base pair overlap between Feature A and Feature B. We can represent base pair overlap by defining

$$\begin{aligned} I_k &= 1 \text{ if position } k \text{ belongs to Feature A and 0 otherwise,} \\ J_k &= 1 \text{ if position } k \text{ belongs to Feature B and 0 otherwise.} \end{aligned}$$

We can then define $X_k = I_k J_k$ to be the indicator that position k belongs to both Feature A and Feature B. Then, for the $n = 30$ Megabases of the ENCODE regions, the mean base pair overlap is equal to

$$\bar{X} = \sum_{k=1}^n X_k / n.$$

Another biologically interesting statistic is the (asymmetric) region overlap, defined as follows: Suppose the consecutive feature stretches are T_1, \dots, T_K with lengths τ_1, \dots, τ_K , and the corresponding non-feature stretches S_1, \dots, S_α with lengths $\rho_1, \dots, \rho_\alpha$. We assume here that the initial and final stretches consist of one feature and one non-feature stretch. The complementary situation, when both initial and final stretches are of the same type is dealt with similarly. Without loss of generality, suppose the initial stretch is non-feature. Then, $S_1 = 1, \dots, \rho_1$, $T_1 = \rho_1 + 1, \dots, \rho_1 + \tau_1$, $S_2 = \rho_1 + \tau_1 + 1, \dots, \rho_1 + \tau_1 + \rho_2$, etc. Using I_i, J_i as indicators of feature identity, we define the (unnormalized) region overlap of feature A stretches with feature B stretches as $\frac{1}{n} \sum_{k=1}^K V_k$ where $V_k = 1 - \prod \{1 - J_i : i \in T_k\}$. This statistic is not linear in terms of functions of single basepairs, but is linear in functions of blocks J_j of basepairs. These blocks are of random sizes, but it is consistent with our hypothesis of piecewise stationarity that, except for end effects due to feature instances crossing segment boundaries, the V_k are also stationary. If the lengths τ_1, \dots, τ_K are negligible compared to n , while K is of the order of n , we can expect that the mixing hypothesis is also valid.

In general, we focus our attention on statistics that can be expressed as a function of the mean of $\mathbf{g}(X_i)$, where \mathbf{g} is some well behaved d -dimensional vector function to be characterized in later sections. By the flexible definition of \mathbf{g} , this encompasses a wide class of situations.

First, we consider vector linear statistics of the form $\mathbf{T}_n(\mathbf{X}) = n^{-1} \sum_{k=1}^n \mathbf{g}(X_k)$. We introduce the following notation:

$$E[\mathbf{T}_n] \equiv \boldsymbol{\mu} \equiv \sum_{i=1}^U f_i \boldsymbol{\mu}_i,$$

where

$$\begin{aligned} \boldsymbol{\mu}_i &\equiv E[\mathbf{g}(X_{i1})], \\ f_i &\equiv n_i/n, \end{aligned}$$

and

$$(3.1) \quad \Sigma_n \equiv \text{Var}(n^{\frac{1}{2}} \mathbf{T}_n) = \sum_{i=1}^U f_i C_i(n f_i),$$

where

$$C_i(m) = C_{i0} + 2 \sum_{\ell=1}^m C_{i\ell} \left(1 - \frac{(\ell-1)}{m} \right)$$

and

$$(3.2) \quad C_{i0} \equiv \text{Var} \mathbf{g}(X_1), \quad C_{i\ell} \equiv \text{Cov}(\mathbf{g}(X_{i1}), \mathbf{g}(X_{i(\ell+1)})) .$$

In Theorem 3.1 below, we prove asymptotic Gaussianity of \mathbf{T}_n given a few more technical assumptions:

A4. $\frac{1}{n} \sum_{i:n_i \leq l} n_i \rightarrow 0$ for all $l < \infty$.

A5. $\forall i, \|\mathbf{g}\|_\infty \leq C < \infty$.

A6. $0 < \varepsilon_0 \leq \|\Sigma_n\| \leq \varepsilon_0^{-1}$, for all n , where $\|\cdot\|$ is a matrix norm.

In particular, A4 implies that the contribution of “small regions” to the overall statistic must not be too large.

Theorem 3.1. *Under conditions A1-A6,*

$$(3.3) \quad n^{\frac{1}{2}} \Sigma_n^{-\frac{1}{2}} (\mathbf{T}_n - \boldsymbol{\mu}) \Rightarrow \mathcal{N}(\mathbf{0}, J)$$

where J is the $d \times d$ identity.

If we have estimates $\hat{\tau}$ of τ which are consistent in a suitably uniform sense, then estimates of $C_{i\ell}$, $C_i(m)$ using $\hat{\tau}$ in place of τ are also consistent. However, simply plugging these estimates into (3) does not yield consistent estimates of σ^2 if our approach were to compute confidence intervals by Gaussian approximation. This is well known for the stationary case. Some regularization is necessary. We do not pursue this direction but prefer to approach the inference problem from a resampling point of view – see next section.

In many cases, the statistics of interest are not linear. For example, in the analysis of the ENCODE data a more informative statistic is the %bp overlap defined as

$$(3.4) \quad B \equiv \frac{\bar{X}}{D},$$

where

$$D = \sum_{k=1}^n I_k$$

is the total base count of feature A .

More conceptually, the region overlap is

$$(3.5) \quad R \equiv \frac{1}{W_I} \sum_{k=1}^K V_k,$$

where $W_I = \sum_{i=1}^n I_{i-1}(1 - I_i)$, the number of Feature A instances.

A standard delta method computation shows that the standard error of B can be approximated as follows: Let $\mu(D)$ and $\mu(\bar{X})$ be respectively the expectation of D and \bar{X} . Then,

$$\frac{\bar{X}}{D} - \frac{\mu(\bar{X})}{\mu(D)} \approx \frac{\bar{X} - \mu(\bar{X})}{\mu(D)} - \mu(\bar{X}) \frac{(D - \mu(D))}{\mu^2(D)},$$

and hence we can approximate $\frac{\bar{X}}{D}$ by a Gaussian variable with mean $\frac{\mu(\bar{X})}{\mu(D)}$ and variance

$$(3.6) \quad \sigma^2(B) \approx \frac{\sigma^2(\bar{X})}{\mu^2(D)} + \frac{\mu^2(\bar{X})}{\mu^4(D)} \sigma^2(D) - 2 \frac{\mu(\bar{X})}{\mu^3(D)} \text{cov}(\bar{X}, D),$$

where $\sigma^2(B)$, $\sigma^2(\bar{X})$, $\sigma^2(D)$ are the corresponding variances and $\text{Cov}(\bar{X}, D)$ denotes the covariance. In doing inference, we can use the approximate Gaussianity of B with $\sigma^2(B)$ estimated using the above formula with regularized sample moments replacing the true moments.

We also note that goodness of fit or equality of population test statistics, such as Kolmogorov-Smirnov and many others, can be viewed as functions of empirical distributions, which themselves are infinite dimensional linear statistics, and we expect, but have not proved, that the methods discussed in this paper and the underlying theories apply to those cases as well, under suitable assumptions.

4. Subsampling Based Methods. Here we propose a subsampling based approach, in particular, a combined segmentation-block subsampling method to conduct statistical inference under the piecewise stationary model, which we call “segmented block subsampling”. In our method, the segmentation parameters governing scale are chosen first and then the size of the subsample is chosen based on stability criteria. The segmentation procedure, as we discussed, is motivated by the heterogeneity of large-scale genomic sequences. The insight of the block subsampling approach is to take into account the local genomic structure, such as natural clumping of features, when conducting statistical inference. We have explicitly demonstrated the advantages of using segmentation and block subsampling by simulation studies in section 5.

4.1. Stationary Block Subsampling.

4.1.1. *Review of Results for the Case of $U = 1$.* For completeness, we recall the following basic algorithm of Politis and Romano (1994) to obtain an estimate of the distribution of the statistic $\mathbf{T}_n(X_1, \dots, X_n)$ under the assumption that the sequence X_1, \dots, X_n is stationary (i.e., $U = 1$).

Algorithm 4.1. a) Given $L \ll n$ choose a number N uniformly at random from $\{1, \dots, n - L\}$.

b) Given the statistic \mathbf{T} , as above, compute

$$\mathbf{T}_L(X_{N+1}, \dots, X_{N+L}) \equiv \mathbf{T}_{L1}^* .$$

c) Repeat B times without replacement to obtain $\mathbf{T}_{L1}^*, \dots, \mathbf{T}_{LB}^*$.

d) Estimate the distribution of $\sqrt{n}(\mathbf{T}_n - \mu)$ by the empirical distribution \mathcal{L}_B^* of

$$\left\{ \sqrt{\frac{n}{L}} [\mathbf{T}_{Lj}^* - \mathbf{T}_n(X_1, \dots, X_n)], \quad 1 \leq j \leq B \right\} .$$

By Theorem 4.2.1 of Politis, Romano and Wolf (1999),

$$(4.1) \quad \mathcal{L}_B^* \implies \mathcal{N}_d(\mathbf{0}, \Sigma) .$$

in probability if (3.3) holds and if $\frac{L}{n} \rightarrow 0$. As usual, convergence of \mathcal{L}_B^* in law in probability simply means that if ρ is any metric for weak convergence on R^d then $\rho(\mathcal{L}_B^*, \mathcal{L}) \xrightarrow{P} 0$.

Since all variables we deal with are in L_2 we take ρ to be the Mallows metric,

$$\rho_M^2(F, G) = \min \{E_P(U - V)^2 : P \text{ such that } U \sim F, V \sim G\} .$$

Useful properties of ρ_M are:

- a) $\rho_M^2(\sum \pi_i F_i, \sum \pi_i G_i) \leq \sum \pi_i \rho_M^2(F_i, G_i)$ for all $\pi_i \geq 0$, $\sum \pi_i = 1$ and
- b) If $F = F_1 * \dots * F_m$, $G = G_1 * \dots * G_m$, that is F and G are distributions of sums of m independent variables, then $\rho_M^2(F, G) \leq \sum_{i=1}^m \rho_M^2(F_i, G_i)$.

For convenience, when no confusion is possible we will write $\rho_M(V, W)$ for $\rho_M(F, G)$ for random variables $V \sim F$, $W \sim G$.

4.1.2. *Performance of the block subsampling method in the piecewise stationary model when $U > 1$.* We turn to the analogue of Theorem 4.2.1 in Politis, Romano and Wolf (1999) for $U > 1$. We consider a vector linear statistic, for which the true distribution was described in Section 3. Here, we ask how Algorithm 4.1, which assumes stationarity, performs in this nonstationary context. We show that in general, it does not give correct confidence bounds but is conservative, sometimes exceedingly so. The results depend on L , the subsample size, which is a crucial parameter in Algorithm 4.1. We sketch these issues in Theorem 4.2 below, for simplicity letting g be the one dimensional identity function $g(x) = x$. Let

$$\begin{aligned} \tau^2 &= U^{-1} \sum_{i=1}^U (\mu_i - \bar{\mu})^2 \\ \bar{X}_i &\equiv n_i^{-1} \sum_{j=1}^{n_i} X_{ij} \quad \bar{X} \equiv n^{-1} \sum_{k=1}^n X_k = \sum_{i=1}^U f_i \bar{X}_i \end{aligned}$$

Also let

$$n_i^* \equiv \text{Cardinality of } S_i \equiv \{k : k \in [N, N + L], \pi_1(k) = i\}$$

and

$$\begin{aligned} \bar{X}_i^* &= 1(n_i^* \neq 0) \sum_j \{X_{ij} : j \in S_i\} / n_i^* , \\ \bar{X}_L^* &= \sum_{i=1}^U f_i^* \bar{X}_i^* , \quad \text{where } f_i^* \equiv \frac{n_i^*}{L} . \end{aligned}$$

We introduce one assumption that is obviously needed for any analysis of the block or segmented resampling bootstraps

A7. $L \rightarrow \infty$

and two other assumptions which are used in different parts of Theorem 4.2 but not in the rest of the paper, and are thus given a different numbering.

B1. $\frac{L}{n} \rightarrow 0$
 B2. $\frac{\sum_{i=1}^U f_i}{n} \rightarrow 0$

Theorem 4.2. *Let \mathcal{L}_n be the distribution which assigns mass f_i to $(\mu_i - \mu)$, $1 \leq i \leq U$ and write C_i for $C_i(nf_i)$. Suppose assumptions A1-A5, and A7 hold.*

- (i) *If B2 holds, $\rho_M(\bar{X}_L^* - \bar{X}, \mathcal{L}_n) \xrightarrow{P} 0$*
 (ii) *If*

$$(4.2) \quad \sum_{i=1}^U f_i (\mu_i - \mu)^2 = o(L^{-1})$$

and B1 holds, then

$$\rho_M[\sqrt{L}(\bar{X}_L^* - \bar{X}), \sum_{i=1}^U f_i \mathcal{N}(0, C_i)] \xrightarrow{P} 0$$

- (iii) *If (4.2) and B1 hold and*

$$(4.3) \quad \sum_{i=1}^U f_i 1(|\Sigma_n - C_i| \geq \varepsilon) \rightarrow 0$$

for all $\varepsilon > 0$, then

$$\rho_M(\sqrt{L}(\bar{X}_L^* - \bar{X}), \mathcal{N}(0, \Sigma_n)) \xrightarrow{P} 0.$$

The implications of Theorem 4.2 are as follows. If equation (4.2) doesn't hold then $\bar{X}_L^* - \bar{X}$ doesn't converge in law at scale $L^{-\frac{1}{2}}$ so that it doesn't reflect the behaviour of $L^{\frac{1}{2}}(\bar{X}_L - \mu)$ at all. This is a consequence of inhomogeneity of the segment means. Evidently in this case, confidence intervals of the percentile type for μ , $[\bar{X} + c_n(\alpha), \bar{X} + c_n(1 - \alpha)]$ where $c_n(\alpha)$ is the α quantile of the distribution of $\bar{X}_L^* - \bar{X}$, will have coverage probability tending to 1, since $c_n(\alpha)$ and $c_n(1 - \alpha)$ do not converge to 0 at rate $L^{-\frac{1}{2}}$ as they should by Theorem 4.2. If B2 does not hold we have to consider the possibility that $[N, N+L]$ covers K_N consecutive segments, whose total length is

$o(n)$, such that the average over all such blocks is close to μ . However, in the absence of a condition such as (4.2) or mutual cancellation of μ_i^* the scale of \bar{X}_L^* will be larger than $L^{-1/2}$. These issues will be clarified by the proof of Theorem 4.2 in the appendix. We note also that (4.2) can be weakened to requiring that the mean of blocks of consecutive segments whose total length is small compared to n are close to μ to order $o(L^{-1/2})$. But our statement makes the issues clear. Finally, note that *B2* holds automatically if the number of segments U is bounded and if *B1* holds.

If (4.2) does hold but (4.3) doesn't, then $\sqrt{L}(\bar{X}_L^* - \bar{X})$ converges in law to the Gaussian mixture $\sum_{i=1}^U f_i \mathcal{N}(0, C_i)$. The mixture of Gaussians is more dispersed in a rough sense than a Gaussian with the same variance, which is,

$$\sigma_n^2 = \sum_{i=1}^U f_i C_i .$$

See Andrews and Mallows(1974). Especially note that, if W has the mixture distribution and V is the Gaussian variable with the same variance, then

$$Ee^{tW} = \sum f_i e^{-\frac{t^2}{2} c_i} \geq e^{-\frac{t^2}{2}} \sum f_i C_i = Ee^{tV}$$

by Jensen's inequality. This suggests that the tail probabilities will also be overestimated. The overdispersion here, which leads to conservativeness that is not as extreme as in case (i), is due to inequality of the variances from segment to segment. Finally, if (4.3) holds then the segments have essentially the same mean and variance and stationary block subsampling works.

A mark of either (4.2) or (4.3) failing is a lack of Gaussianity in the distribution of $\bar{X}_L^* - \bar{X}$. This was in fact observed at some scales in the ENCODE project, which led us to crudely segment on biological grounds with reasonable success. However, the correct solution, which we now present in this paper, is to estimate the segmentation and appropriately adjust the subsampling procedure.

4.2. A Segmentation Based Block Subsampling Method. We saw in the previous section that the naïve block subsampling method that was designed for the stationary case breaks down when the sequence follows a piecewise stationary model. We propose a stratified block subsampling strategy, which stratifies the subsample based on a “good” segmentation of the sequence which is estimated from the data. We first state the block subsampling method, and then in Section 4.2.3 give minimal conditions on the estimated segmentation for its consistency. In Section 4.3 we discuss possible segmentation methods.

4.2.1. *Description of Algorithm.* Assume that we are given a segmentation $\mathbf{t} = (0 = t_0, t_1, \dots, t_{m+1} = n)$, where m is the number of regions in \mathbf{t} . Assume that the total size L of the subsample is pre-chosen. We define a stratified block subsampling scheme as follows.

Algorithm 4.3. For $i = 1, \dots, m$, let $\lambda_i = \lambda_i(t) = \lceil (t_i - t_{i-1})L/n \rceil$. We use the notation $X_{i;l}$ to denote the block of length l starting at i :

$$X_{i;l} = (X_i, \dots, X_{i+l-1}).$$

Then, for each subsample,

Draw integers $\mathbf{N} = \{N_1, \dots, N_m\}$, with N_i chosen uniformly from $\{(t_{i-1} + 1, \dots, t_i - \lambda_i(t) + 1)\}$, and let

$$X^* = (X_1^*, \dots, X_m^*) = (X_{N_1; \lambda_1(t)}, \dots, X_{N_m; \lambda_m(t)}).$$

Repeat the above B times to obtain B subsamples: $X^{*,1}, \dots, X^{*,B}$.

To obtain a confidence interval for $\boldsymbol{\mu}$, we assume that the statistic \mathbf{T}_n has approximately a $N(\boldsymbol{\mu}, \Sigma_n/n)$ distribution as in the previous section. For each subsample drawn as described in Algorithm 4.3, compute the statistic $\mathbf{T}_L^{*,b} = \mathbf{T}_L^{*,b}(\mathbf{t}) = \mathbf{T}_L(X^{*,b})$. Form the sampling estimate of variance,

$$(4.4) \quad \widehat{\Sigma}_n \equiv \frac{L}{B} \sum_{b=1}^B (\mathbf{T}_L^{*,b} - \bar{\mathbf{T}}_L^*)' (\mathbf{T}_L^{*,b} - \bar{\mathbf{T}}_L^*),$$

where $\bar{\mathbf{T}}_L^* \equiv \sum_{b=1}^B \mathbf{T}_L^{*,b} / B$. We can now proceed to estimate the confidence interval for \mathbf{T}_n in a sequence of standard ways. For example, in the univariate case where $\sigma_n^2 \equiv \Sigma_n$ is a scalar:

- a). Use $\bar{X} \pm z_{1-\alpha/2} \frac{\widehat{\sigma}_n}{\sqrt{n}}$, where z_β is the β th quantile of $N(0, 1)$, for a $1 - \beta$ confidence interval.
- b). Efron's percentile method: Let $\bar{X}_{(1)}^* < \dots < \bar{X}_{(B)}^*$ be the ordered $\bar{X}^{*,b}$, then use $[\bar{X}_{(\lfloor B\alpha/2 \rfloor)}^*, \bar{X}_{(\lfloor B(1-\alpha/2) \rfloor)}^*]$ as a $1 - \alpha$ confidence interval.
- c). Use a studentized interval (Efron (1981)) or Efron's (1987) *BCA* method, see Hall (1992) for an extensive discussion.

Although the theory for (c) giving the best coverage approximation has not been written down, as it has been for the ordinary bootstrap, we expect it to continue to hold. Evidently, these approaches can be applied not only to vector linear statistics like T_n but also to smooth functions of vector linear statistics.

This algorithm assumes a given segmentation \mathbf{t} , which should be set to some good estimate $\hat{\boldsymbol{\tau}}^{(n)} = \{0 = \hat{t}_0, \hat{t}_1, \dots, \hat{t}_m = n\}$ of the true change points $\boldsymbol{\tau}^{(n)}$. In order for the algorithm to perform well, a good segmentation is critical unless the sequence is already reasonably homogeneous. In Section 4.2.2 below we state the result that the algorithm is consistent if the given segmentation equals the true changepoints. Then, in Section 4.2.3, we state a few assumptions on the data determined segmentation $\hat{\boldsymbol{\tau}}^{(n)}$ which would enable us to act as if the segmentation were known and state Theorem 4.5 to that effect.

4.2.2. Consistency with True Segmentation. Under the hypothetical situation where the segmentation \mathbf{t} assumed in Algorithm 4.3 is the true segmentation, then the algorithm can be easily shown to be consistent. Here we state the result, which will be proved in Appendix.

First, we state a stronger version of the assumption B1, which requires that the square of the subsample size $L = L_n$ to be $o(n)$:

$$\text{A8. } L_n^2/n \rightarrow 0.$$

Then, the consistency of Algorithm 4.3 given the true segmentation is given in the following theorem.

Theorem 4.4. *If assumptions A1-A8 hold, then*

$$(4.5) \quad L_n^{1/2} \Sigma_n^{-1/2} [T_{L_n}^*(\tau_n) - T_n] \Rightarrow N(0, J)$$

in probability.

4.2.3. Consistency with Estimated Segmentation. Let $\hat{\boldsymbol{\tau}} = \hat{\boldsymbol{\tau}}^{(n)} = (\hat{\tau}_1^{(n)}, \dots, \hat{\tau}_{\hat{U}_n}^{(n)})$ be a segmentation of the sequence X_1, \dots, X_n , which is determined from the data, and which is intended to estimate the true change-points $\boldsymbol{\tau} = \boldsymbol{\tau}^{(n)}$. We will state conditions on $\hat{\boldsymbol{\tau}}$ such that the statistic obtained from Algorithm 4.3 based on $\hat{\boldsymbol{\tau}}$ is close to the statistic obtained from the same algorithm based on the true segmentation $\boldsymbol{\tau}$. This can be stated formally as follows. For any segmentation \mathbf{t} , let $\mathbf{X}^*(\mathbf{t})$ be a subsample drawn according to Algorithm 4.3 based on \mathbf{t} . Let $F_{n,\mathbf{t}}^*(\cdot)$ be the distribution of $\sqrt{L}\{T[\mathbf{X}^*(\mathbf{t})] - E^*T[\mathbf{X}^*(\mathbf{t})]\}$ conditioned on X_1, \dots, X_n and \mathbf{t} . Then, the desired property on the estimated segmentation $\hat{\boldsymbol{\tau}}$ is that

$$(4.6) \quad \rho_M^2[F_{n,\hat{\boldsymbol{\tau}}^{(n)}}^*, F_{n,\boldsymbol{\tau}^{(n)}}^*] \rightarrow_p 0, \quad \text{as } n \rightarrow \infty$$

where ρ_M^2 is the Mallows' metric described in Section 4.1.1. That is, for inferential purposes, $T[\mathbf{X}^*(\hat{\boldsymbol{\tau}})]$ has approximately the same distribution as

$T[\mathbf{X}^*(\boldsymbol{\tau})]$. Then, since we have shown in Section 4.2.2 that

$$\rho_M^2[F_{n,\boldsymbol{\tau}^{(n)}}^*, \Phi(\Sigma_n)] \rightarrow_p 0,$$

where $\Phi(\Sigma_n)$ is the Gaussian distribution with mean 0 and variance Σ_n , (4.6) implies that

$$\sqrt{L_n} \Sigma_n^{-1} \{T[\mathbf{X}^*(\hat{\boldsymbol{\tau}}^{(n)})] - E^*T[\mathbf{X}^*(\boldsymbol{t})]\} \rightarrow N(0, J).$$

Let $\hat{n}_i = \hat{\tau}_{i+1}^{(n)} - \hat{\tau}_i^{(n)}$. We now state conditions on the estimated segmentation which guarantee (4.6).

- A9. $\frac{\hat{U}_n}{n} \rightarrow 0$,
- A10. $\frac{1}{n} \sum_{i:\hat{n}_i \leq k} \hat{n}_i \rightarrow 0$ for all $k < \infty$,
- A11. $\frac{L_n}{n} \sum_{i=1}^{U_n} \min_{1 \leq j \leq \hat{U}_n} |\tau_i - \hat{\tau}_j| \rightarrow_p 0$.

Assumptions A9 and A10 for $\hat{\boldsymbol{\tau}}^{(n)}$ mirror assumptions A3 and A4 for $\boldsymbol{\tau}^{(n)}$. Assumption A11 is a consistency criterion: As the data set grows, the total discrepancy in the estimation of $\boldsymbol{\tau}^{(n)}$ by $\hat{\boldsymbol{\tau}}^{(n)}$ must be small.

Theorem 4.5. *Under assumptions A1-A11, (4.6) holds.*

The proof is given in the appendix. There are trivial extensions of this theorem to smooth functions of vector means, which are, in fact, needed but simply cloud the exposition.

Theorem 4.5 implies that confidence intervals based on subsamples

$$\{\mathbf{X}^{*,j}(\hat{\boldsymbol{\tau}}^{(n)}) : j = 1, \dots, B\}$$

constructed by Algorithm 4.3 conditional on $\hat{\boldsymbol{\tau}}^{(n)}$ are consistent, as long as $\hat{\boldsymbol{\tau}}^{(n)}$ satisfies A9-A11. Here is the formal statement of this fact in the one dimensional case, where $\hat{\sigma}_n^2$ replaces $\hat{\Sigma}_n$ and \mathbf{g} is the identity.

Corollary 4.6. *Under assumptions A1-A11,*

1. *Let $\hat{\sigma}_n^2$ be the block subsampling estimate of variance defined in (4.4), then*

$$P(\bar{X} - z_{1-\alpha/2} \hat{\sigma}_n / \sqrt{n} < \mu < \bar{X} + z_{1-\alpha/2} \hat{\sigma}_n / \sqrt{n}) \rightarrow_p 1 - \alpha.$$

2. *Confidence intervals estimated by Efron's percentile method are consistent. That is,*

$$P([\bar{X}_{([n\alpha/2])}^* < \mu < \bar{X}_{([n(1-\alpha/2)])}^*]) \rightarrow_p 1 - \alpha.$$

4.3. *Segmentation Methods.* The objective of the segmentation step is to divide the original data sequence X_1, \dots, X_n into approximately homogeneous regions so that the variance estimated in Algorithm 4.3 approximates the true variance of T_n . A segmentation into regions of constant mean is sufficient for guaranteeing that Algorithm 4.3 gives consistent variance estimates. Therefore, we focus here on the segmentation of X into regions of constant mean.

In our simulation and data analysis, we use the dyadic segmentation approach, which we motivate and describe here using the simple case of g identity. First consider the simple case where X_1, \dots, X_n are independent with variance 1. In testing the null hypothesis

$$H_0 : E[X_i] = \mu,$$

versus the alternative H_A that there exists $1 < \tau < n$ such that $E[X_i] = \mu_1$ for $i < \tau$ and $E[X_i] = \mu_2 \neq \mu_1$ for $i \geq \tau$, one can show that the following is the generalized likelihood ratio test:

$$\text{Reject } H_0 \text{ if } \max_{1 < j < n} nM(j) > c,$$

where

$$(4.7) \quad M(j) = \frac{j}{n} (\bar{X}_{1:j} - \bar{X}_{1:n})^2 + \frac{n-j}{n} (\bar{X}_{j+1:n} - \bar{X}_{1:n})^2.$$

The maximum likelihood estimate of the change-point parameter τ is the value that maximizes $M(j)$.

Our proof of Theorem 4.5 in the appendix shows that, in the case where there is one true change in mean at τ , the increase in the variance estimated by block subsampling with block length L given no segmentation (i.e. $\mathbf{t}^{(n)} = \{0, n\}$) over the variance estimated by Algorithm 4.3 conditioned on a change-point at τ is $LM(\tau) + o_p(1)$. Subsampling conditioned on any segmentation $t \neq \tau$ would give an even larger variance estimate. Hence, segmenting at $\hat{\tau} = \arg \max_j M(j)$ is optimal in the sense that $\hat{\tau}$ is the change-point estimate that minimizes the asymptotic error of the block subsampling variance estimate. This fact does not require the assumption of independence observations, and is true for any stationary sequence. Thus, if we knew that there were only one change-point, and if the goal of the segmentation is to obtain the best stratified variance estimate, then the best place to segment is t . The block subsampling variance estimate, given the

segmentation $\{0, t, n\}$, would be

$$(4.8) \quad \begin{aligned} V(t) &= \left(\frac{t}{n^2}\right) \sum_{i=1}^{t-tL/n} (\bar{X}_{i:i+tL/n} - \bar{X}_{1:t})^2 \\ &+ \left(\frac{n-t}{n^2}\right) \sum_{i=t+1}^{n-(n-t)L/n} (\bar{X}_{i:i+(n-t)L/n} - \bar{X}_{t+1:n})^2. \end{aligned}$$

The Dyadic segmentation procedure recursively applies the above logic, as described below.

Algorithm 4.7. Fix minimum region length $0 < L_s < n$ and threshold $b > 0$. Initialize $\mathbf{t} = \{t_0 = 0, t_1 = n\}$. Repeat:

1. For $i = 1, \dots, |\mathbf{t}| - 1$, let $M^{(i)}(j)$ and $V^{(i)}(j)$ be respectively the processes (4.7) and (4.8) computed on the subsequence $X_{t_{i-1}+1}, \dots, X_{t_i}$. If $t_i - t_{i-1} > 2L_s$ then let $t'_i = \arg \max_{t_{i-1}+L_s < j < t_i-L_s} M^{(i)}(j)$, $B_i = M^{(i)}(t'_i)$, and $V_i = V^{(i)}(t'_i)$. Otherwise, let $B_i = 0$, $V_i = \infty$.
2. let $\lambda_i = L(t_i - t_{i-1})/n$, and

$$J_i = 1 \left(\frac{(t_i - t_{i-1})B_i}{\sqrt{V_i \hat{\lambda}_i}} > b \right).$$

If $\prod_i J_i = 0$, stop and return \mathbf{t} .

3. Let $i^* = \arg \max_i B_i$, and $t^{new} = t'_{i^*}$.
4. Let $\mathbf{t} = \mathbf{t} \cup t^{new}$, reordered so that t_i is monotonically increasing in i .

Each step of the recursion in Algorithm 4.7 proceeds as follows: In step 1, $M^{(i)}(j)$, the between group sum of squares, and $V^{(i)}(j)$, the block subsampling variance estimates, are computed for each segment $[t_{i-1} + 1, t_i]$ of the current segmentation. B_i is the maximum squared difference in mean for segment i , t'_i is the change-point estimate that achieves this maximum, and $\hat{\lambda}_i V_i$ is the estimate of variance given a change-point at t'_i . In computing B_i and V_i we do not allow break points that create a region with length less than L_s . In step 2, we normalize the statistic $(t_i - t_{i-1})B_i$ by the estimated standard deviation $\sqrt{\hat{\lambda}_i V_i}$. If this normalized statistic is below the boundary b for every subsegment, then the recursion stops and returns the current segmentation. Otherwise, in step 3, the optimal location to segment $t^{(new)}$ is chosen to be the cut that maximizes the decrease in error of the block subsampling variance estimate, conditioned on the fact that it had passed the thresholding in step 2. In step 4, this new change-point $t^{(new)}$ is added to the current segmentation \mathbf{t} .

The normalization by V_i in step 2 requires an appropriate choice $L = L_b$ of the block subsampling sample size. Often it is easier to choose this parameter after the segmentation is given. If a ball park value of L_b is not available, then the normalization by V_i can be omitted, in which case the parameter b in step 3 should be set to 0. This would be equivalent to stopping the segmentation only when the next optimal cut will violate the minimum region length L_s . In the examples of Section 5.1 we set $b = 0$, thus decoupling the choice of L_s from that of L_b .

The two parameters required by Algorithm 4.7 are L_s and b . The choice for L_s is discussed in Section 4.5. The choice of b can be guided by the fact that under the null hypothesis, if L were chosen appropriately, then $(t_i - t_{i-1})M^{(i)}(j)/[V^{(i)}(j)\hat{\lambda}_i]^{1/2}$ is a pivot with approximate distribution χ_1^2 . Asymptotic approximations for the family-wise error rate have been derived in the case of independent sequences (James et al., 1987). In the case of dependent sequences a Bonferroni adjustment can be applied to adjust for multiple testing. We also used the formulas given in James et al. (1987) to get a crude cutoff, which seems to work in practice.

Algorithm 4.7 belongs to the class of dyadic segmentation algorithms for detection of change-points, the consistency of which were studied by Vostrikova (1981). These algorithms are greedy procedures that avoid the search over all possible segmentations, which would be computationally intractable. They have been applied successfully to various settings in biology, including segmentation of GC content (Li et al., 2002) and the analysis of DNA copy number data (Olshen et al., 2004).

The consistency of Algorithm 4.3 requires conditions A9-A11 to be satisfied by the estimated segmentation. The key condition is A11 which defines a consistency criterion on the segmentation. Consistency of dyadic segmentation has been proved in Vostrikova (1981) for sequences that satisfy the following conditions:

1. Let $\epsilon_t = X_t - E[X_t]$, then $\|\epsilon_t\|^2$ is a submartingale and $E\|\epsilon_t\|^2 < ct^\beta$, $c > 0$, $\beta < 2$.
2. The number of regions is fixed and the region sizes are of order n , i.e.

$$\tau_n = (nr_1, \dots, nr_U), \quad 0 < r_1 < \dots < r_U.$$

It is easy to verify that condition 1 is satisfied by the piecewise stationary model due to the mixing condition A2. Condition 2 is more stringent than our assumptions A3 and A4, under which U_n is allowed to increase with n . The consistency of dyadic segmentation for the case of $U_n \rightarrow \infty$ has been explored in Venkatraman (1992), who gave asymptotic conditions on the rejection threshold and on the sizes of the regions to ensure consistency under

the assumption of an independent Gaussian sequence. However, these conditions are hard to verify in practice, and we believe that for our applications in genomics the more stringent condition of Vostrikova (1981) is sufficient. Previous studies on segmenting the genome based on features such as the GC content (Fu and Curnow (1990), Li et al. (2002)) have used this finite regions assumption to achieve reasonable results.

The dyadic segmentation procedure uses information from the entire sequence to call the first change, and then recursively uses all of the information from each subsegment to call each successive change in that segment. An alternative is to use pseudo-sequential procedures, which are sequential (online) schemes that have been adapted for change-point detection when the entire sequence of a fixed length is completely observed. The basic idea of this class of methods is to do a directional scan starting at one end of the sequence. Every time a change-point is called, the observations prior to the change-point are ignored and the process starts over to look for the next change after the previously detected change-point. Specifically, let $\hat{\tau}_0 = 0$, and given $\hat{\tau}_1, \dots, \hat{\tau}_k$,

$$\hat{\tau}_{k+1} = \inf\{l > \hat{\tau}_k : S(X_{\hat{\tau}_k}, X_{\hat{\tau}_{k+1}}, \dots, X_{\hat{\tau}_l}) > b\},$$

where S is a suitably defined change-point statistic and b is a pre-chosen boundary. The estimates from pseudo-sequential schemes depend on the direction in which the data is scanned. Thus, while they may be suitable for, say, timeseries data, they may not be natural for segmentation of genomic data, which in most cases do not have an obvious directionality. The consistency of pseudo-sequential procedures has been studied by Venkatraman (1992), who gave conditions on $b = b_n$ and $\hat{\tau}^{(n)}$ for consistency of $\hat{\tau}^{(n)}$ under the setting that X_i are independent Gaussian with changing means.

4.4. *Testing the Null Hypothesis of No Associations.* Here we extend the results in section 4.2 to testing the null hypothesis of no association using non-linear statistics. As we discussed in Section 1.2, the inference problem typically posed in high-throughput genomics is that of association of two features. In terms of our framework we have two 0-1 processes $\{I_k\}_{k=1, \dots, n}$ and $\{J_k\}_{k=1, \dots, n}$ both defined on a segment of length n of the genome. We assume that the joint process $\{I_k, J_k\}$ is piecewise stationary and mixing and want to test the hypothesis that the two point processes $\{I_k\}_{k=1, \dots, n}$ and $\{J_k\}_{k=1, \dots, n}$ are independent. We have studied two fairly natural test statistics in ENCODE, the “percent basepair overlap”,

$$B_n = \frac{\sum_{k=1}^n I_k J_k}{\sum_{k=1}^n I_k},$$

and the ‘‘regional overlap,’’ R_n , which we define in Section 3, with large values of these statistics indicating dependence. The major problem we face in constructing a test is what critical values $o_{n\alpha}, r_{n\alpha}$ we should specify so that

$$(4.9) \quad P_{H_0}[B_n \geq o_{n\alpha}] \approx \alpha,$$

where H_0 is the hypothesis that the vectors $(I_1, \dots, I_n)^T$ and $(J_1, \dots, J_n)^T$ are independent, and the corresponding $r_{n\alpha}$ for R_n .

We aim for statistics based on B_n, R_n (respectively) which are asymptotically Gaussian with mean 0 under H_0 . In general, we have to be careful about our definition of independence. If we interpret H_0 as we stated, simply as independence of the vectors $(I_1, \dots, I_n)^T$ and $(J_1, \dots, J_n)^T$, then

$$E_{H_0}(B_n) \approx \frac{\sum_{i=1}^U \sum_{k=1}^n n_i E_{H_0}(I_{ik}) E_{H_0}(J_{ik})}{\sum_{i=1}^U \sum_{k=1}^n n_i E_{H_0}(I_{ik})}$$

where I_{ik} and J_{ik} refer to the k th basepair in the i th segment, and hence we have,

$$(4.10) \quad E_{H_0}(B_n) \approx \frac{\sum_{i=1}^U \lambda_i E_{H_0}^{(i)}(I) E_{H_0}^{(i)}(J)}{\sum_{i=1}^U \lambda_i E_{H_0}^{(i)}(I)}$$

The natural estimate of this expectation is then,

$$\frac{1}{\bar{I}} \sum_{i=1}^U \lambda_i \bar{I}_i \bar{J}_i$$

where $\lambda_i \equiv \frac{n_i}{n}$, \bar{I}_i is the average of I_{ik} , \bar{J}_i is the average of J_{ik} , and \bar{I} is the grand average. We assume the correct segmentation.

We proceed with construction of a test statistic and estimation of the null distribution. In view of (4.10) our test statistic based on B_n is

$$(4.11) \quad T_n^O \equiv n^{\frac{1}{2}}(B_n - \tilde{J}_n)$$

where

$$\tilde{J}_n \equiv \left(\sum_{i=1}^{\hat{U}} \hat{\lambda}_i \hat{I}_i \hat{J}_i \right) / \frac{1}{n} \sum_{j=1}^n \hat{I}_j$$

$$(4.12) \quad \text{where } \hat{\lambda}_i = \lambda_i(\hat{\mathbf{t}}), \quad \hat{I}_i = n_i^{-1}(\hat{\mathbf{t}}) \sum_{j=\hat{t}_{i-1}+1}^{\hat{t}_i} I_j$$

with \hat{J}_i similarly defined. Here is the algorithm based on this statistic.

Algorithm 4.8. We do the following.

1. Pick at random without replacement two starting points, K_1 and K_2 , of blocks of length L from $\{1, \dots, n - L\}$.
2. Let $(I_{K_1+1}, \dots, I_{K_1+L})^T$ and $(J_{K_1+1}, \dots, J_{K_1+L})^T$, $(I_{K_2+1}, \dots, I_{K_2+L})^T$ and $(J_{K_2+1}, \dots, J_{K_2+L})^T$ be the two sets of two feature indicators.
3. Form

$$\overline{IJ}_{nL}^{*1} \equiv \frac{1}{L} \sum_{l=1}^L I_{K_1+l} J_{K_2+l}$$

$$\bar{I}_{nL}^{*1} \equiv \frac{1}{L} \sum_{l=1}^L I_{K_1+l}$$

$$\overline{IJ}_{nL}^{*2} \equiv \frac{1}{L} \sum_{l=1}^L I_{K_2+l} J_{K_1+l}$$

and define \bar{I}_{nL}^{*2} , \bar{J}_{nL}^{*1} , \bar{J}_{nL}^{*2} analogously. Let

$$F_{nL}^* \equiv \frac{1}{2} \left(\frac{\overline{IJ}_{nL}^{*1}}{\bar{I}_{nL}^{*1}} + \frac{\overline{IJ}_{nL}^{*2}}{\bar{I}_{nL}^{*2}} \right)$$

$$T_{nL}^* \equiv F_{nL}^* - \bar{J}_{nL}^*$$

where

$$\bar{J}_{nL}^* = \frac{1}{2} (\bar{J}_{nL}^{*1} + \bar{J}_{nL}^{*2})$$

and \bar{I}_{nL}^* is defined analogously. Let F_{nLb}^* , \overline{IJ}_{nLb}^{*1} etc. be obtained by choosing (K_{1b}, K_{2b}) , $b = 1, \dots, B$ independently as usual.

4. We use the following $c_{nL\alpha}$ as a critical value for B_n at level α ,

$$c_{nL\alpha} = \bar{J}_n + \left(\frac{2L}{n} \right)^{\frac{1}{2}} T_{nL(B(1-\alpha))}^*$$

where $T_{nL(1)}^* \leq \dots \leq T_{nL(B)}^*$ are the ordered T_{nLb}^* and $[\cdot]$ denotes integer part and $\bar{J}_n = \frac{1}{n} \sum_{k=1}^n J_k$.

5. If the sequence is piecewise stationary with estimated segments $j = 1, \dots, \hat{U}_n$ as in Section 4.3, we draw independently B sets of starting points, $K_{11}^{(j)}, \dots, K_{1B}^{(j)}$ and $K_{21}^{(j)}, \dots, K_{2B}^{(j)}$, of blocks of length $\hat{\lambda}_j L$ from each segment $i = 1, \dots, j$ when each pair is drawn at random without replacement. Here $\sum_{i=1}^{\hat{U}_n} \hat{\lambda}_i = 1$ and $\hat{\lambda}_i$ is proportional to the length of

estimated segment i . Then piece T_{nLb}^* together as follows. Let

$$\begin{aligned}\bar{I}J_{nLb}^{*1i} &= \frac{1}{L\hat{\lambda}_i} \sum_{l=1}^{\hat{\lambda}_i} I_{iK_{1b+l}} J_{iK_{2b+l}} \\ \bar{I}_{nLb}^{*1i} &= \frac{1}{L\hat{\lambda}_i} \sum_{l=1}^L I_{iK_{1b+l}} \\ &\text{etc} \\ \bar{F}_{nLb}^* &= \sum_{i=1}^{\hat{U}} \hat{\lambda}_i \left(\frac{\bar{I}J_{nLb}^{*1i}}{\bar{I}_{nLb}^{*1i}} + \frac{\bar{I}J_{nLb}^{*2i}}{\bar{I}_{nLb}^{*2i}} \right).\end{aligned}$$

Then,

$$T_{nLb}^* = F_{nLb}^* - \tilde{J}_{nLb}^*,$$

where

$$\tilde{J}_{nLb}^* = \frac{\sum_{i=1}^{\hat{U}} (\bar{I}_{nLb}^{*i})(\bar{J}_{nLb}^{*i})\hat{\lambda}_i}{\sum_{i=1}^{\hat{U}} (\bar{I}_{nLb}^{*i})\hat{\lambda}_i}$$

with $\bar{I}_{nLb}^{*i} = \bar{I}_{nLb}^{*1i} + \bar{I}_{nLb}^{*2i}$. The critical value is,

$$\tilde{J}_n + \left(\frac{2L}{n}\right)^{\frac{1}{2}} T_{nL(B(1-\alpha))}^*,$$

as before.

We can apply this principle more generally to statistics which are functions of sums of products of I 's and J 's evaluated at the same positions.

The proof of the following theorem is given in the Appendix.

Theorem 4.9. *If \mathcal{L}_0, P_0 denote distributions under the hypothesis of independence and (A1)-(A11) hold, then*

1. $\mathcal{L}_0(T_n^O) \implies \mathcal{N}(0, \sigma_0^2)$
2. *With probability tending to 1,*

$$\mathcal{L}_0^*(T_{n,L}^{O*}) \implies \mathcal{N}(0, \sigma_0^2)$$

3. $P_0[T_n^O \geq \left(\frac{2L}{n}\right)^{\frac{1}{2}} \hat{q}_{1-\alpha}^0] \rightarrow \alpha$ where $\hat{q}_{1-\alpha}^0$ is the $[(1-\alpha)B]$ th of T_{nLb}^{O*} , $1 \leq b \leq B$.

In practice, this definition of independence makes our statistic in effect reflect *conditional independence* of I_i and J_i given the segment to which i belongs. This can be unsatisfactory in practice, for instance, when the

features are concentrated in small segments such that large, sparse, segments swamp the inference.

We define *independence irrespective of segment identity* as saying that the average over all permutations of the segments of the joint distribution of the point process features are independent. Formally, if (P_1, \dots, P_U) , (Q_1, \dots, Q_U) denote the marginal distributions of $\{\{I_{ik} : k = 1, \dots, n_i\} : i = 1, \dots, U\}$ and $\{\{J_{ik} : k = 1, \dots, n_i\} : i = 1, \dots, U\}$, and (R_1, \dots, R_n) correspond to the joint distribution of $\{(I_{ik}, J_{ik}) : 1 \leq k \leq n\}$, then let $(\bar{P}_1, \dots, \bar{P}_U) = \frac{1}{U!} \sum (P_{\pi_1}, \dots, P_{\pi_U})$ where π ranges over all permutations of $1, \dots, U$. Define $(\hat{Q}_1, \dots, \hat{Q}_U)$ and $(\hat{R}_1, \dots, \hat{R}_U)$ similarly. Then, our hypothesis is

$$(4.13) \quad H_1 : \hat{R} = \hat{P} \times \hat{Q}$$

This is simply saying that independence is not conditional on relative genomic position of segments.

It is easy to see that we should now define

$$(4.14) \quad T_n^{\tilde{O}} = n^{\frac{1}{2}}(B_n - \hat{J}_n)$$

where $\hat{J}_n = \frac{1}{n} \sum_{i=1}^n J_i$.

The reason for this is that

$$(4.15) \quad E_{\hat{R}}(B_n) \approx \frac{E_{\hat{R}}(\frac{1}{n} \sum_{i=1}^n I_i J_j)}{E_{\hat{R}}(\hat{I})}$$

Under H_1 ,

$$E_{\hat{R}}(\frac{1}{n} \sum_{i=1}^n I_i J_j) = E_{\hat{P}}(\hat{I}) E_{\hat{Q}}(\hat{J})$$

and

$$E_{\hat{R}}(I) = E_{\hat{P}}(\hat{I})$$

so that the statistic simplifies to the $U = 1$ form, as above.

It is clear that the conclusion of (4.9) continues to hold when applied to $T_n^{\tilde{O}}$. Note that the form of the bootstrap is unchanged, since $T_n^{\tilde{O}}$ is invariant under permutation of the segments.

We now turn to R_n as defined in Section 3. We assume that $V_k : k = 1, \dots, K$ are strongly mixing and stationary. If we assume H_0 we have no closed form

for $E_{H_0}(\frac{1}{W} \sum_{k=1}^K V_k)$ by which to center R_n . To estimate this quantity we apply a version of the double bootstrap (Beran (1988); Hall (1992); Letson and McCullough (1998)).

Consider $\frac{1}{n} \sum_{k=1}^K V_k$ under H_1 . We draw B_1 pairs of large blocks of length mL , and we compute the % false region overlap, call it R_b^* , $b = 1, \dots, B$, in each pair of “large” blocks, where mL is still negligible compared to segment size, but $m \rightarrow \infty$. Define

$$(4.16) \quad \hat{E}_{H_1}(R_n) = \frac{1}{2B} \sum_{b=1}^{B_1} R_b^*$$

and

$$(4.17) \quad \tilde{T}_n^{(R)} = n^{\frac{1}{2}}(R_n - \hat{E}_{H_1}(R_n))$$

Note that we again want to consider independence irrespective of segment identity, so that R_b^* above are computed without any segmentation beyond the natural segmentation, e.g. chromosomes. Now compute the empirical distribution of $\tilde{T}_n^{(R)}$ using the size L segmented block subsampling and proceed as usual. We can define $\tilde{T}_n^{(R)}$ corresponding to H_0 in the same way, though we now have to cut up our mL blocks in proportion to segment sizes to center. We do not pursue this since the H_1 hypothesis gives stable results while H_0 does not.

We have not proved a result justifying the use of the double bootstrap in this way, but simulations suggest that it behaves as expected, see Section 5.3.

4.5. *Choice of Segment Size L_s .* Two tuning parameters appear in our procedure in addition to b appearing in the segmentation scheme. L_s is the smallest allowed size of a “stationary” piece after segmentation. It essentially determines the scale of the segmentation, which we view as an application context dependent quantity that users need to control. The reason is that stationarity is a matter of scale. To put it concretely, consider the situations where I_j , $j = 1, \dots, n$ are simply the base pair nucleotides A, C, G, T and consider the scale of a large gene of length n . Then, it seems natural that the exons and introns correspond to consecutive stationary regimes. However, suppose we now move our scale to a gene rich genomic region of length N . Now, it is the genes themselves and the intergenic regions which correspond to an at least initial segmentation.

This dependence of segmentation on scale has a natural intuitive consequence. Consider a statistic such as base pair overlap of two features. As one

increases the region size n in which one wishes to declare significant overlap, the standard deviation of the statistic, which is $O(n^{-1/2})$, decreases, and p-values decrease. However, if, as one would expect, the region over which n increases becomes homogeneous on a larger scale, coarser segmentation would then be called for. This, as we have noted, necessarily increases the standard deviation of the statistic, and from that point of view significance becomes more difficult to achieve.

Put another way, it is not impossible to think of the whole genome itself as being stationary on a large scale, but that we can hierarchically segment the genome in many ways so that each large subsegment is stationary, but the segments are not identically distributed, even where they are of equal length. For instance, a natural initial segmentation is to chromosomes.

Finally, we argue in mathematical terms going the other way from inhomogeneity to homogeneity. Start with a sequence of independent (say) Bernoulli variables X_1, X_2, \dots, X_n , with X_j being Bernoulli(p_j). If the p_j are arbitrary, the only segmentation we can perform is the useless trivial one, where each X_j is its own segment. But, now if suppose that we tell ourselves that p_j , $1 \leq j \leq n/2$ are drawn i.i.d. from $U(0, 1/2)$ and for $n/2 + 1 \leq j \leq n$ from $U(1/2, 1)$, we suddenly just have two segments to consider.

Thus, L_s in our view needs to be treated as the smallest scale on which homogeneity is expected. Note that these considerations are not limited to testing. They also govern confidence intervals, as discussed in section 4.2.3.

4.6. *Choice of L_b , the subsample size.* We believe that the best way to choose L_b , after segmentation has been estimated, is so that the resulting subsampling distribution of the statistics is as stable as possible and L_s is large but $\ll n$. We also formally consider Gaussianity of the distribution and, if possible, maximizing that feature as well. This does not necessarily mean segment more – since A10 and A11 may then fail. We advocate but do not analyze further the following proposal put forward in m -out-of- n subsampling by Bickel, Götze, and van Zwet (1997) and analyzed in detail by Götze and Rackauskas (2001) and Bickel and Sakov (2005):

1. Let $\bar{X}_n^*(L)$ be the statistic computed from the sample drawn with blocks of length L . Compute the block subsampling distribution \mathcal{L}_{L_k} for the statistic

$$\sqrt{L_k}(\bar{X}_n^*(L_k) - \bar{X}_n)$$

and $L_k = \rho^k n$, where $\rho < 1$ and $k = 1, 2, \dots, K$.

2. Compute a “distance” $d^*(k)$ between \mathcal{L}_{L_k} and $\mathcal{L}_{L_{k-1}}$.
3. Choose $L_b = L_{k_0^*}$, where $k_0^* = \arg \min d^*(k)$.

In practice, we use for $d^*(k)$ the pseudometric $|\sqrt{\frac{L_{k-1}}{L_k}} IQR(\mathcal{L}_{L_k}) - IQR(\mathcal{L}_{L_{k-1}})|$, where $IQR(\mathcal{L})$ is the interquartile range of \mathcal{L} .

In continuing work with Götze, van Zwet, we are in the process of trying to show that, under mild conditions, as $n \rightarrow \infty$ we have $L_b \rightarrow \infty$, $L_b/n \rightarrow 0$. More significantly, we expect that in a fashion analogous to Götze and Rackauskas (2001) and Bickel and Sakov (2005), under restrictive conditions and for suitable choice of distance, L_b yields an estimate which is as good as possible in the following sense: If \mathcal{L}_m is the actual distribution of $\sqrt{m}(\bar{X}_m - \mu)$, $d(k)$ is the distance between \mathcal{L}_m and \mathcal{L}_{L_k} , and $k_0 = \arg \min_k d(k)$, then

$$\frac{d(k_0^*)}{d(k_0)} \rightarrow_p c.$$

Thus, $L_{k_0^*} = \rho^{k_0^*} n$ yields performance of the same order as $\rho^{k_0} n$.

5. Simulation and Data Studies.

5.1. *Simulation Study I.* In this section, we perform a simple simulation study to demonstrate the power of our block-subsampling method in the situation where features are naturally clustered. We simulated a binary sequence x_1, \dots, x_n with $n = 10,000$ by the following Markovian model:

$$(5.1) \quad P(x_1 = 1) = \frac{p_0}{2}, P(x_i = 1) = \frac{p_0}{2} + (1 - p_0) \frac{\sum_{k=i-w}^{i-1} x_k}{w} \text{ for } i = 2, \dots, n,$$

where w is the order of the Markov model, or intuitively, the size of the dependency window, and p_0 indicates the level of dependency. The smaller p_0 is, the stronger the dependence between the neighboring positions is. We define the following two types of features at position i in the sequence:

- Feature I: the occurrence of sequence 11100 starting at position i
- Feature II: the occurrence of more than six 1's in the next 10 consecutive positions including the current position i .

From model (5.1), the feature II will occur in clusters in the sequence. The overlap between the two types of features can be measured by the statistic

$$S = \frac{\sum_{k=1}^n I_k J_k}{\sum_{k=1}^n I_k}$$

with I_k, J_k being binary and indicating the occurrences of sites of Type I and II respectively.

Figure 1 shows the distribution of S estimated through different ways:

- The true distribution is the empirical distribution of estimated S from 10000 random sequences generated under model (5.1)
- The Ordinary Bootstrap distribution is derived by performing a base-by-base uniform sampling of the sequence x_1, \dots, x_n to construct 10000 sequences of length n .
- The Feature Randomization distribution is derived by keeping features of type I fixed and randomizing uniformly the start positions of the features of type II to construct 10000 sequences of length n .
- The block subsampling distribution is derived by drawing independent samples of blocks of length $L = 40$ and stringing the blocks together to construct 10000 sequences of length n .

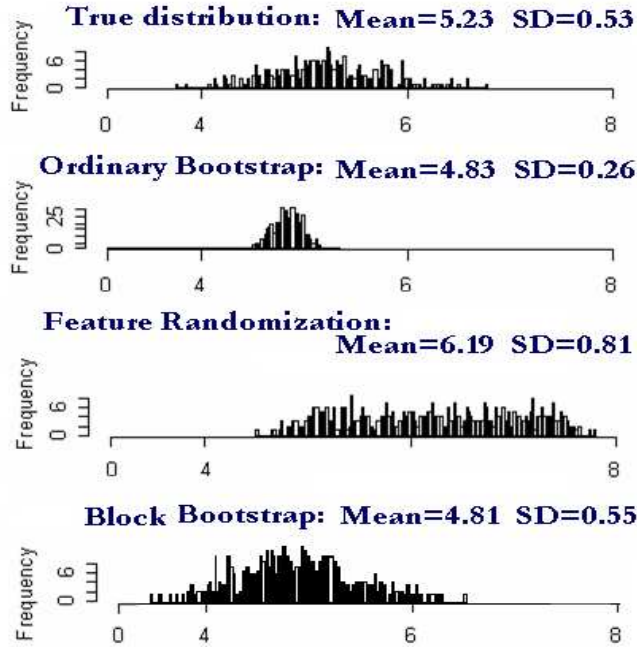


Fig 1. Comparison of different subsampling schemes

We see that block subsampling produces more reliable estimates of the variance of S compared to the naive methods: ordinary bootstrapping and feature randomization. Both naive methods ignore the dependence between positions and thus fail to take into account the natural clumps of the feature II. This is the key reason for the poor performance of the two naive methods.

5.2. *Simulation Study IIa.* Our second simulation study examines the case where the sequence is generated from a piecewise stationary model where there is more than one homogeneous region. As before, we consider the problem of estimating the percentage of base pair overlap between two features, and compare the performance of four strategies:

1. feature randomization,
2. naive block subsampling from unsegmented sequence,
3. block subsampling from sequence segmented using the true change-points, and
4. block subsampling from sequence segmented using the change-points estimated by binary segmentation.

In our simulation model, we generate X_t, Y_t independently from a Neyman-Scott process characterized as follows:

1. Cluster centers occur along the sequence according to a Poisson process of rate λ_i in region i .
2. The number of features in each cluster follows Poisson distribution with mean α .
3. The start of features are located at a geometric distance (mean μ) from the cluster center.
4. The features are generated with length that is geometric with mean β .
5. Overlap between features generated using steps 1-4 are ignored.

For simplicity we let there be only 2 homogeneous regions, each of length $T = 10000$. Consider the setting where the parameters for the two regions have the following values: $(\lambda_1, \alpha_1, \mu_1, \beta_1) = (0.01, 10, 10, 5)$ and $(\lambda_2, \alpha_2, \mu_2, \beta_2) = (0.02, 10, 10, 5)$. Figure 2 shows a simulated example, where features A and B are plotted as well as their overlap. Figure 2 also shows the cumulative sum and the segmentation. Figure 3 shows respectively the histograms of the estimated distribution of the overlap statistic \bar{X}^* centered and scaled. It is clear that the feature randomization underestimates the standard deviation, whereas naive block subsampling without segmentation gives a mixture distribution with long tails. Strategy 3, which subsamples assuming the true changepoint at τ is known, gives the correct distribution as expected. Strategy 4, which uses the estimated change-point, reassuringly gives a very similar distribution to Strategy 3. Table 1 gives the standard deviation estimates.

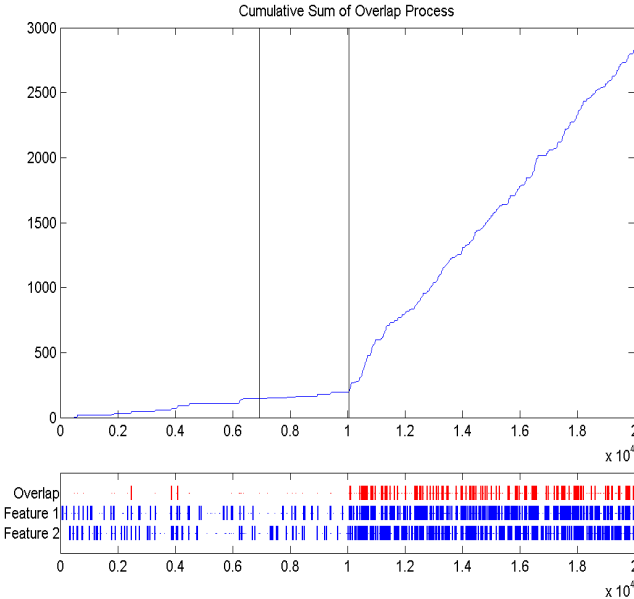


Fig 2. Example of one instance from simulation model 2. Top plot shows cumulative sum and estimated segmentation.

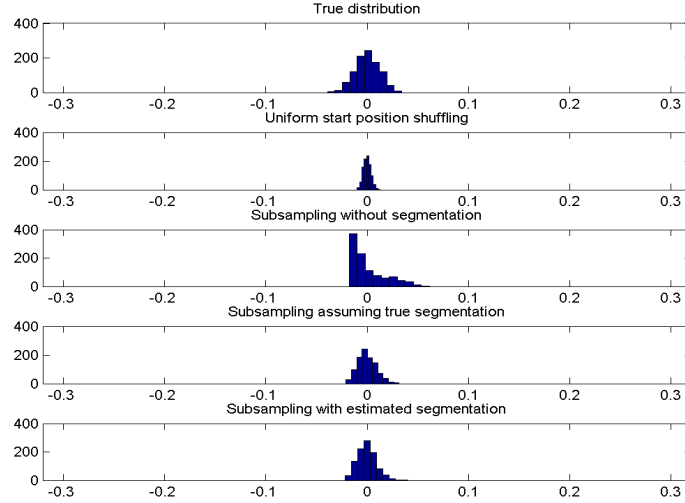


Fig 3. Comparison of different subsampling schemes

Method	Standard Error Estimate	Fold change from true value
True value	1.2e-002	–
Uniform shuffle	3.6e-003	0.3
Subsample, no segmentation	1.7e-002	1.4
Subsample, true segmentation	1.1e-002	0.91
Subsample, estimated segmentation	1.0e-002	0.83

Table 1. Estimates of standard error by four sampling strategies in simulation study 2a.

5.3. *Simulation Study IIb.* We utilized the Neyman-Scott process described in Simulation Study IIa to study the consistency of the double bootstrap method described in Section 4.4 for estimating the distribution of R_n . We consider the simple case where there is one homogeneous region. We utilized a larger region and a parameterization of the process that results in more and longer feature instances than we consider in the study above. We use $T = 5\text{Mb}$ and $(\lambda_1, \alpha_1, \mu_1, \beta_1) = (\lambda_2, \alpha_2, \mu_2, \beta_2) = (0.05, 10, 100, 75)$. This yields a pair of feature-sets with around 5,000 instances, where each feature-set covers around 17% of the 5Mb region. We simulated 20,000 pairs

of feature-sets from this process, and found that the mean of region-overlap between pairs, R_n , was 0.293, and the standard error was 0.0072. We subsampled 1,000 sets of 10,000 draws from this distribution, each of which yielded the mean above (to 3 significant digits), and the standard errors ranged from 0.0071 to 0.0073, which corresponds almost exactly to the theoretical 95% confidence interval for the standard error of the standard error of a Gaussian with standard deviation 0.0072 after 10,000 draws. No surprisingly, the distribution of R_n was Gaussian, as indicated by the Lilliefors and the Shapiro-Wilk test, which did not reject the hypothesis of Gaussianity at a significance level 0.05 with the full sample of 20,000 observations.

In order to test the capacity of segmented subsampling with a version of the double block bootstrap to discover this distribution based on only a single pair of observations, we selected the most extreme pair found during simulation, for which R_n was 0.321, corresponding to a z-score of 3.87. Since the number of feature instances is itself a random quantity, the job of block subsampling is particularly difficult: when R_n is far to the right of expectation, the feature-sets tend to contain more feature instances than those closer to the center. The pair we chose was no exception. The results are given in Table 2. Hence, it is not surprising that our subsampling procedure tends to over-estimate the mean. The Lilliefors test fails to reject the Gaussianity of any of the resulting distributions with sample sizes up to 1,000 at a significance level of 0.05, but does reject it for several of the smaller block-sizes when the sample size is pushed up to 10,000. The Shapiro-Wilk test, however, detects departures from Gaussianity for many of the distributions at a significance level of 0.05 for samples larger than 500. This is because R_n is predicated on relatively small counts of feature-instance overlaps, and hence the distributions tend to have heavy tails.

We note that the global minimum of the Inter-Qantile (IQ) statistic was found at $L_b/L_r = 0.15$ and $L_r/n = 0.06$. That is, 0.9% of the 5Mb region, or 45Kb, were included in each block sample. This block sample size is certainly sufficient to capture multiple feature-clusters, since the parameterized Neyman-Scott process above yields an average inter-cluster distance of about 1Kb.

To corroborate our hypothesis that the mean was overestimated because the feature-sets we chose were more dense than most, we applied our method with learned parametrization, $L_b/L_r = 0.15$ and $L_r/n = 0.06$, for a pair of feature-sets with $R_n = 0.293$, the population average. Indeed, the mean was estimated, after 10,000 samples, to be 0.293, and $\hat{\sigma}_n$ was 0.0072.

Although the purpose of this simulation was merely to check the consistency of our version of the double block bootstrap for data not unlike

actual genomic data, e.g. ChIP-seq “broad-peaks”, we decided to also check the performance of feature-start site shuffling for the same pair of feature-sets used above. In the case of B_n , the basepair overlap statistic, feature start-site shuffling correctly estimates the mean, but can (in the stationary case), radically underestimate the standard deviation. The same is not true in the case of R_n . Start-site shuffling is not assured (under our model) to provide an unbiased estimate of the mean or the standard deviation. We drew 10,000 samples from the distribution under shuffling, and found the mean to be 0.337, and the standard deviation to be 0.0070, which indicates that the pair of feature-sets under study in fact overlap slightly less than expected at random ($p \approx 0.011$). The fact that this conclusion is actually in the wrong direction in this relatively easy, stationary example, should make us skeptical of studies that rely upon start-site shuffling to draw conclusions about statistics that cannot be defined locally, such as R_n .

Our discussion of this simulation and the following real data examples exhibit the subtleties inherent in our approach. Subtleties appearing whenever inference follows regularization.

(L_r/n)	$L_b/L_r = 0.1$			$L_b/L_r = 0.15$			$L_b/L_r = 0.2$		
	$d^*(k)$	Mean	$\hat{\sigma}_n$	$d^*(k)$	Mean	$\hat{\sigma}_n$	$d^*(k)$	Mean	$\hat{\sigma}_n$
0.10	0.0240	0.3032	0.0074	0.0231	0.3029	0.0073	0.0181	0.3034	0.0072
0.09	0.0116	0.3042	0.0072	0.0064	0.3035	0.0072	0.0092	0.3035	0.0071
0.08	0.0070	0.3040	0.0072	0.0067	0.3038	0.0072	0.0049	0.3041	0.0071
0.07	0.0032	0.3039	0.0071	0.0042	0.3043	0.0072	0.0017	0.3042	0.0071
0.06	0.0003	0.3042	0.0071	0.00004	0.3044	0.0072	0.0002	0.3040	0.0071
0.05	0.0036	0.3042	0.0072	0.0070	0.3040	0.0071	0.0031	0.3042	0.0071
0.04	0.0115	0.3045	0.0071	0.0031	0.3044	0.0071	0.0057	0.3043	0.0071
0.03	0.0194	0.3047	0.0072	0.0178	0.3045	0.0072	0.0169	0.3047	0.0070
0.02	0.0571	0.3045	0.0071	0.0498	0.3047	0.0071	0.0396	0.3049	0.0071
0.01	N/A	0.3047	0.0071	N/A	0.3047	0.0071	N/A	0.3047	0.0071

Table 2: Comparison of block subsampling distributions

5.4. *Association of Non-coding ENCODE annotations and Constrained Sequences..* Here we present a real example of the study of association between “constrained sequences” and “non-exonic annotations” from the ENCODE project, limited to the 1.87Mbp ENCODE Pilot Region ENm001, also known as the CFTR locus. The constrained sequences are those highly conserved between human and the 14 mammalian species studied and sequenced by the ENCODE consortium. Enrichment of evolutionary constraint at the “non-exonic annotations” sites implies that the biochemical

assays employed by the ENCODE consortium are capable of identifying biologically functional elements. We tested the association of non-coding annotations and constrained elements using the base pair overlap statistic B_n in Section 4.3, using the conditional formulation. We interpret the lack of association as, given sequence composition and the distribution of each feature along the genome as observed, the assignments (by Nature) of feature A and feature B to individual bases are made independently. We derive the significance of the observed statistic under this null hypothesis following the method proposed in Section 4.3.

As we discussed, we have several issues to deal with:

- i) How do we segment? That is, what statistic(s) do we use for segmentation?
- ii) Is segmentation necessary or is the region sufficiently homogeneous?
- iii) If we segment, what L_s should we use?
- iv) Given a segmentation, what L_b is appropriate?

Here are our methods:

- a) The simplest choice for i) and the one we followed was to segment according to both numerator and denominator in B_n : intersect partitions and enforce an L_s bound. Given our theory, this should ensure homogeneity in the mean of B_n .
- b) Although strictly speaking ii) and iii) can be combined, we experimented a bit to also see if the theory of Section 4.1 was borne out in practice.
- c) We did not use the V statistic and thus only had to choose L_s . Again, we experimented with $L_s = 500Kb$ to preserve as much genomic structure as possible, and $L_s = 200Kb$ to ensure we had not undersegmented.
- d) We explored a variety of values of L_b , and studied the consistency between nearby values under the interquartile statistic (IQ statistic) discussed in Section 4.6. We draw conclusions based on the value of L_b that optimizes local consistency.

To segment the data, we applied the method in Section 4.3 to both features A and B, or in the language of Section 4, I and J , and then combined the segmentation. In segmenting each feature, we experimented with minimum segment lengths L_s of 200 and 500 Kb. Before subsampling, we combined the segmentations of A and B by taking a union of the change-points. This created regions with length less than L_s . However the total length of these regions comprise $< 0.1\%$ of the total Encode region, and were left out of the remaining analyses.

If the sequence were sufficiently homogeneous, we could forgo the initial segmentation step. Figure 4 shows an estimate of variance of B_n (with the appropriate renormalization) for a reasonable range of L_b , both before and after segmentation. Two trends are clearly evident. First, segmentation greatly reduces the estimated variance. As we discussed in Section 4.1.2, inhomogeneity of the sequence causes an inflated estimate of variance. If the data were homogeneous, segmentation should not change the variance estimate. Thus, the fact that the estimated variances drop after segmentation for such a large range of L_b 's suggests that the data is inhomogeneous. Secondly, and more importantly, the estimated variance of B_n increases sharply with increasing L_b in the unsegmented data. This is evidence of inhomogeneity in the mean of B_n across this ENCODE region: underlying shifts in mean, if ignored, can be mistaken for spurious long range autocorrelation, which also implicitly runs against our assumption. In either case, as Theorem 4.2.1 suggests, we would be overly conservative. Thus, a preliminary exploration of the data convinces us that this ENCODE region is inhomogeneous in I and/or J and segmentation is necessary.

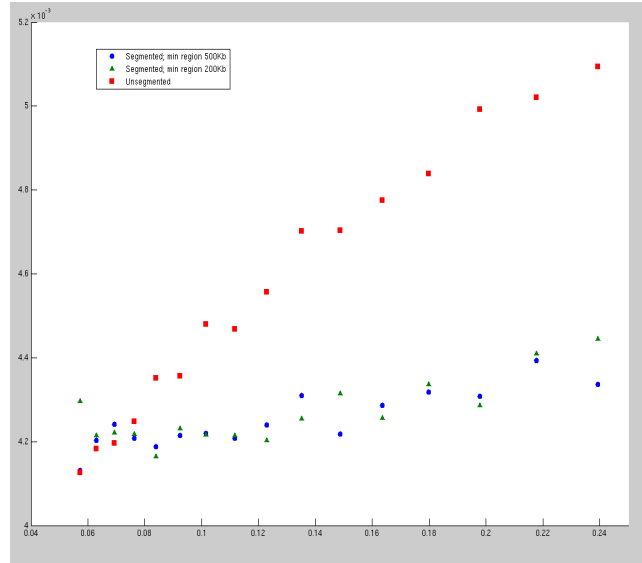


Fig 4. Estimated σ_n as a function of L_b for 10,000 samples

We found that 200Kb and 500Kb gave 5 and 3 segments respectively. Table 2 gives the results for 500Kb. What is fairly surprising, but reassur-

ing, is that over the whole broad range of L_b considered, the estimated SD of the statistic under the null was essentially flat after segmentation. Flat here means that variability was within a Monte Carlo SD for the 10,000 replications we used. We would expect longer values of L_b to include, in our estimate of σ , additional covariance between distant genomic positions captured by the extended block-length. The fact that this, by and large, does not appear to be happening is consistent with our hypothesis that the relevant mixing distance is indeed quite small compared to the size of approximately stationary regimes.

We found that there is still moderate deviation from Gaussianity in both the segmented and unsegmented case for $0.05 < L_b < 0.25$, both in the tails, as detected by the Shapiro-Wilk test, and in the body of the distribution under Lilliefors test. With a sample size of 100, neither test detects this departure, but at a sample size of only 500, it is detected under a number of parameterizations of L_b . As we discussed in Section 4.5, the definition of stationarity depends on the scale at which we view the genome. This suggests that our segmentation still does not take care of inhomogeneity in the variance. Hence, as we have mentioned, if we use the variance for the Gaussian approximation our results are still conservative.

The scientific conclusion of this example is that, indeed, there is strong association since the z value is over 9 SDs. We note that the effect of segmentation on our scientific conclusion is essentially non-existent. However, it is comforting to note that the change in (with and without segmentation) variance is in the correct direction.

L_b/n ($\rho = 0.99$)	B_n								
	unsegmented			$L_s = 500Kb$			$L_s = 250Kb$		
	$d^*(k)$	$\hat{\sigma}_n$	z-score	$d^*(k)$	$\hat{\sigma}_n$	z-score	$d^*(k)$	$\hat{\sigma}_n$	z-score
0.239	0.032	0.0051	8.47	0.043	0.0043	9.95	0.027	0.0044	9.71
0.217	0.010	0.0050	8.60	0.026	0.0044	9.82	0.015	0.0044	9.79
0.197	0.033	0.0050	8.65	0.008	0.0043	10.02	0.024	0.0043	10.07
0.179	0.010	0.0048	8.92	0.026	0.0043	10.00	0.015	0.0043	9.96
0.163	0.007	0.0048	9.04	0.017	0.0043	10.07	0.034	0.0043	10.14
0.148	0.017	0.0047	9.18	0.027	0.0042	10.23	0.016	0.0043	10.01
0.135	0.005	0.0047	9.18	0.017	0.0043	10.02	0.027	0.0044	10.14
0.122	0.019	0.0045	9.47	0.019	0.0042	10.18	0.015	0.0042	10.27
0.111	0.007	0.0045	9.66	0.011	0.0042	10.26	0.022	0.0042	10.24
0.101	0.005	0.0045	9.64	0.025	0.0042	10.23	0.022	0.0042	10.24
0.092	0.016	0.0044	9.91	0.006	0.0042	10.25	0.017	0.0042	10.20
0.083	0.006	0.0044	9.92	0.030	0.0042	10.31	0.015	0.0042	10.37
0.076	0.014	0.0042	10.16	0.011	0.0042	10.26	0.022	0.0042	10.24
0.069	0.004	0.0042	10.29	0.013	0.0042	10.18	0.014	0.0042	10.23
0.063	0.001	0.0042	10.32	0.013	0.0042	10.27	0.018	0.0042	10.24
0.057	N/A	0.00412	10.46	N/A	0.0041	10.45	N/A	0.0043	10.05

Table 3: Comparison of block subsampling distributions, $\rho^\beta n$ vs. $\rho^{\beta+1} n$ under the IQR statistic. Estimates $\hat{\sigma}_n$ and resulting z -scores of B_n shown.

5.5. *The association of Copy Number Variation with RefSeq annotated Exons in the human genome.* In this example, we re-analyzing a published dataset; this re-analysis leads to a different conclusion from the one made by the original paper. In 2006, Redon et al. published a set of 1,445 genomic regions with observed Copy Number Variation (CNVs) across individuals. These regions consist of both deletions and insertions, and more than half of them overlap genes. In the paper, the authors reported, amongst other things, a paucity of overlap with RefSeq genes at a significance level of 0.05. The statistic that they used is precisely our marginal formulation of R_n , but the null distribution to which they referred it is quite different. Their null was computed by randomly permuting both genes and CNVs, and hence treats the entire genome (or at least entire chromosomes), as homogeneous, and the distances between feature-instances as exponential. Thus, if feature-instance lengths were all 1bp, this would be a Poisson process. As discussed in Simulation Study IIb, under our model this procedure provides a non-biased estimate of the mean in the case of the B_n , but is unpredictable with respect to its estimate of the variance. In the case of R_n , it is unpredictable with respect to both the mean and the variance. Here, for comparison with the result of Redon et al. (2006), we examine only R_n .

Although we have attempted to replicate this portion of the Redon study, undoubtedly there are small differences between our efforts and those of Redon et al. (2006). For instance, we have masked all genomic repeats in the “Repeat Masker” track on the UCSC genome browser (genome.ucsc.edu). Redon et al. also considered patterns of repeats in their analysis, but may have utilized an at least slightly different map of genomic repeats. We find that 61.8% of the CNVs overlap RefSeq genes by at least 1 basepair. That is, we wish to assess the significance of our observed statistic $R_n = 0.618$.

The calibration of the subsampling procedure is non-trivial, especially in this application where we must consider the additional parameter L_r . Hence, in the following we provide complete detail regarding the calibration of our method for the data of Redon et al. (2006).

As before, our analysis begins with an assessment of the need for segmentation. In this case, we are dealing with whole human chromosomes, we expect that, in general, at least some segmentation is necessary. We segmented down to a minimum segment length of 10,000,000bps (10Mbs), letting $L_s = 10Mb$. The mean length of these CNVs is around 250Kb, and they are not uniformly distributed, so we are compelled not to segment down to regions much smaller than 10Mb by our desire to capture the appropriate spatial distribution of clusters of feature-instances. To assess the sufficiency of the resulting segmentation, we examine the Gaussianity of the segmented subsampling distributions. This examination is tied to our selection of block length.

To select an inner block length, L_b , and an outer block-length, L_r , we drew 10,000 samples for each of several lengths. We chose to use a linear, rather than exponential, scale for L_r/n : we selected 10 values from 0.01 to 0.10 in increments of 0.01. We chose three values of L_b/L_r , 0.05, 0.10, and 0.20. Each of these parameterizations yields several responses, including: an estimated z-score, $d^*(k)$, and measures of Gaussianity. In Figure 5, we plot the relationship between the estimated z-score, $d^*(k)$, L_r and L_b . Regarding the Gaussianity of the resulting distributions, at a significance level of 0.01 and a sample size of 5,000, neither the Shapiro-Wilk nor the Lilliefors test rejected the null hypothesis of Gaussianity for any of the 30 explored parameterizations. To supplement our biological intuition that segmentation is necessary when whole chromosomes are considered, we used the same 30 parameterizations with the unsegmented data, and performed the same tests to check the Gaussianity of the resulting distributions. Of the 30 parameterizations, 3 showed departures from Gaussianity under Lilliefors test, and 9 showed strong departures in the tails under the Shapiro-Wilks test. This indicates, as expected, that segmentation has substantially improved

the Gaussianity of the sample distributions. In practice, one might attempt a finer segmentation in hopes of further reducing the (conservative) bias in $\hat{\sigma}_n$. For this example we are satisfied with the current segmentation.

The global minimum of $d^*(k)$ occurs for $L_r/n = 0.09$ and $L_b/L_r = 0.20$. This parameterization yields an estimated z-score of 1.25, and therefore we conclude that we cannot corroborate the result of Redon et al. (2006). Under our model it appears that CNVs are, if anything, very slightly positively associated with genes ($p \approx 0.105$). We note that a few parameterizations, as shown in Figure 5, do produce z-scores greater than 2. However, these parameterizations correspond to large values of $d^*(k)$, and furthermore, significance is in the opposite direction reported by Redon et al. (2006). This highlights the need for carefully defined null distributions in genomic studies. We are not suggesting that the results presented necessarily invalidate the corresponding result of Redon et al. (2006), but rather we caution that scientific conclusions of this kind are predicated on how the researcher defines “at random”, and that this definition should be made to reflect, as much as possible, that which is known about the actual distribution of genomic elements. We presume that authors wish, in general, to err on the side of caution, and hence do not wish to report significant association when the association can be explained simply by a conservative choice of null.

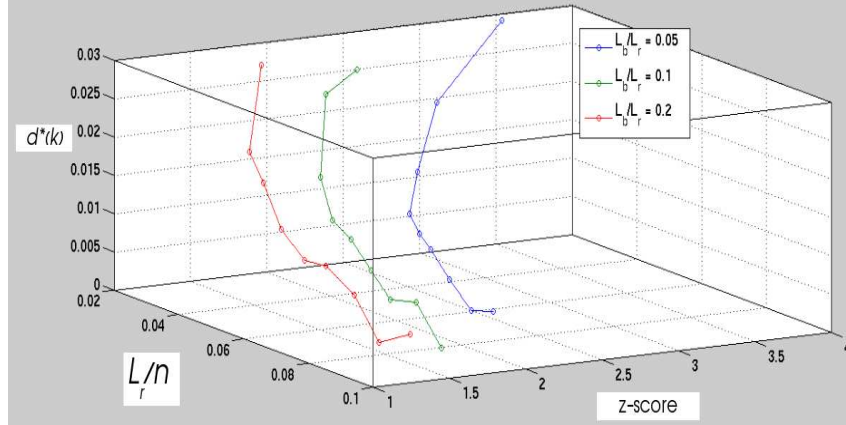


Fig 5. The relationship between the estimated z-score, $d^*(k)$, L_r and L_b . As L_r increases, our estimate of $\hat{\sigma}_n$ (not shown) increases, which drives the estimated z-score down. As L_r becomes too small, we lose the stability of our estimates, and $d^*(k)$ increases. For the smallest value of L_r shown here, the estimated z-score increases sharply, but the corresponding value of $d^*(k)$ indicates that this parameterization is unreliable. The ideal parameterization under $d^*(k)$ is given by $L_r/n = 0.09$ and $L_b/L_r = 0.20$.

6. Appendix.

6.1. *Proof of Theorem 3.1.* For simplicity we have $d = 1$ and \mathbf{g} the identity. The general case follows by the Cramér-Wold device. Write,

$$\text{Var}(S_n/\sqrt{n}) = A_n + B_n,$$

where

$$\begin{aligned} A_n &= \frac{1}{n} \sum_{a=1}^n \sum_{b=1}^n \text{Cov}[X_a, X_b] I(\pi_1(a) = \pi_1(b)), \\ B_n &= \frac{1}{n} \sum_{a=1}^n \sum_{b=1}^n \text{Cov}[X_a, X_b] I(\pi_1(a) \neq \pi_1(b)). \end{aligned}$$

We will show that

$$(6.1) \quad B_n = o(I)$$

The theorem then follows from, for example, Corollary 1, page 142 of Herndorff (1984).

PROOF OF (6.1): We first note that, by A2 we have the standard bound

$$(6.2) \quad |\text{Cov}[X_a, X_{a+k}]| \leq C^2 m(k),$$

and by A2, since $\beta > 1$,

$$\sum_{k \geq k_0} m(k) \leq c \sum_{k \geq k_0} k^{-\beta} \rightarrow 0 \quad \text{as } k_0 \rightarrow \infty.$$

Thus, for all $\epsilon > 0$, exists $k_1(\epsilon)$ such that for all $k \geq k_1(\epsilon)$,

$$(6.3) \quad \frac{1}{n} \sum_{a=1}^n \sum_{b=1}^n |\text{Cov}[X_a, X_b]| I(\pi_1(a) \neq \pi_1(b), |a-b| > k) \leq \epsilon.$$

Now, by A4, since

$$\sum_{a=1}^n I(\pi_1(a+k) > \pi_1(a) + 1) \leq \sum_{i:n_i \leq k} n_i = o(n),$$

by A2 and (6.2) we have,

$$\begin{aligned} & \frac{1}{n} \sum_{a=1}^n \sum_{b=1}^n |\text{Cov}[X_a, X_b]| I(\pi_1(a) \neq \pi_1(b), |a-b| \leq k) \\ (6.4) \quad & \leq \frac{C^2}{n} \sum_{a=1}^n \sum_{b=1}^n I(\pi_1(b) = \pi_1(a) + 1, |a-b| \leq k) + o(1). \end{aligned}$$

The first term on right hand side of (6.4) is bounded by

$$(6.5) \quad \frac{C^2}{n} \sum_{i=1}^U \sum_{j=n_i-k+1}^{n_i} 1 \leq 2C^2 U k/n.$$

Thus by A3, the above expression is $o(1)$. Combining (6.3-6.5) we obtain (6.1). Evidently, $A_n = \sigma_n^2$, and thus Theorem 4.2 follows.

6.2. Proof of Theorem 4.2..

(i) We use P^* throughout here for the randomization measure. By B2,

$$(6.6) \quad P^*[\pi_1(N) = \pi_1(N + L - 1)] \rightarrow 1$$

since the complementary event can happen iff $\pi_1(N) = i, \pi_2(N) = n_i - L + 1, \dots, n_i$ for some i and that probability is bounded by $\sum_{i=1}^U \frac{L}{n}$. But given $\pi_1(N) = \pi_1(N + L - 1)$, \bar{X}_L^* is then a draw from the finite population $\{L^{-1} \sum_{j=k}^{k+L-1} X_{ij} : 1 \leq k \leq n_i - L + 1\}$. By A5 and A7 and Dedeker et al. (2004)

$$(6.7) \quad \max\{\rho_M(L^{-1} \sum_{j=k}^{k+L-1} X_{ij}, \mu_i) : 1 \leq k \leq n_i - L + 1, 1 \leq i \leq U\} \xrightarrow{P} 0$$

and (i) follows from 6.6, 6.7 and property b) of ρ_M .

(ii) It is enough to show that

$$\rho_M\left(\mathcal{L}^*(\sqrt{L}(\bar{X}_L^* - \bar{X}) \mid \pi_1(N) = i), \mathcal{N}(0, C_i)\right) \xrightarrow{P} 0$$

for each i . But by Dedeker et al (2004),

$$\rho_M\left(L^{-\frac{1}{2}} \sum \left\{ (X_k - \mu_i) : \pi_1(k) = i, j \leq \pi_2(k) \leq j+L-1 \right\}, \mathcal{N}(0, C_i)\right) \rightarrow 0$$

for each i uniformly in j by the stationarity of $\{X_k : \pi_1(k) = i\}$ as $L \rightarrow \infty$. But

$$\sqrt{L}(\bar{X}_L^* - \bar{X}) = \sqrt{L}(\bar{X}_L^* - \mu) + o_P(1)$$

by B1. Further, as we have shown,

$$\rho_M(\sqrt{L}(\bar{X}_L^* - \mu_{\pi_1(N)}), \mathcal{N}(0, C_i)) \xrightarrow{P} 0$$

and finally,

$$L \sum_{i=1}^U f_i(\mu_i - \mu)^2 = o(1),$$

thus the result follows by Property b) of ρ_M .

(iii) It is enough to show that

$$\rho_M\left(\sum f_i \mathcal{N}(0, C_i), \mathcal{N}(0, \Sigma_n)\right) \rightarrow 0 .$$

But

$$\rho_M(\mathcal{N}(0, C_i), \mathcal{N}(0, \Sigma_n)) = O(|C_i - \Sigma_n|)$$

and again

$$\rho_M^2\left(\sum f_i \mathcal{N}(0, C_i), \sum f_i \mathcal{N}(0, \Sigma_n)\right) \leq \sum f_i \rho_M^2(\mathcal{N}(0, C_i), \mathcal{N}(0, \Sigma_n)) = o(1).$$

□

6.3. *Proof of Theorem 4.4.* Let $\bar{g}_i = \sum_{j=1}^{n_i} g(X_{ij})/n_i$. By Theorem 4.2.1 of Politis, Romano and Wolf (1999), if $f_i L \rightarrow \infty$, then

$$(f_i L)^{1/2} [T_{f_i L}(X_{N_i; f_i L}) - \bar{g}_i] \Rightarrow N(0, C_i(f_i L))$$

in law in probability. Since by the stratified block subsampling algorithm, the terms on the left above are independent across i under P^* , for every $\epsilon > 0$ there exists $U(\epsilon)$ independent of n such that

$$E^* \left| L^{1/2} \sum_{i=U(\epsilon)+1}^U f_i [T_{f_i L}(X_{N_i; f_i L}) - \bar{g}_i] \right|^2 \leq c^2 \sum_{i=U(\epsilon)+1}^U f_i \leq \epsilon$$

and hence,

$$E^* \left| L^{1/2} (T_L^* - T_n) - L^{1/2} [T_L^*(U(\epsilon)) - T_n(U(\epsilon))] \right|^2 \leq \epsilon$$

where

$$(6.8) \quad T_L^*(U(\epsilon)) = \sum_{i=1}^{U(\epsilon)} f_i T_{f_i L}(X_{N_i; f_i L})$$

$$(6.9) \quad T_n(U(\epsilon)) = \sum_{i=1}^{U(\epsilon)} f_i \bar{g}_i.$$

The result then follows if

$$\min\{nf_i : 1 \leq i \leq U(\epsilon)\} \rightarrow \infty$$

for all $\epsilon > 0$. This is implied by assumption A4.

6.4. *Proof of Theorem 4.5.* We first define some notation. Let $\hat{R}_i = \{\hat{\tau}_{i-1} + 1, \dots, \hat{\tau}_i\}$ be the region between $\hat{\tau}_{i-1}$ and $\hat{\tau}_i$, and

$$k_i = \sum_{j=1}^{U_n} I_{\{\hat{\tau}_{i-1} < \tau_j < \hat{\tau}_i\}} + 1$$

be the number of stationary regions within \hat{R}_i . Let $\tau_{i,0} = \hat{\tau}_i$ and for $j = 1, \dots, k_i - 1$, $\tau_{i,j} = \min\{j : \tau_j > \tau_{i,j-1}\}$. Thus $\tau_{i,j}$ is the j -th true change-point in region \hat{R}_i , and let

$$\tilde{\tau} = \tilde{\tau}^{(n)} = \tau^{(n)} \bigcup \hat{\tau}^{(n)} = \{\tau_{i,j} : i = 1, \dots, U_n, j = 1, \dots, k_i\}.$$

Define $R_{ij} = \{\tau_{i,j-1} + 1, \dots, \tau_{i,j}\}$ to be the j -th stationary region in \hat{R}_i . Then, define

$$\begin{aligned} f_{ij} &= |R_{ij}|/n, \\ \hat{f}_i &= |\hat{R}_i|/n, \\ \lambda_{ij} &= f_{ij}L, \\ \hat{\lambda}_i &= \hat{f}_iL. \end{aligned}$$

We use the notation $X_{s;t} = \{X_s, \dots, X_{s+t-1}\}$ to denote the block of size t starting at s , and $\bar{X}_{s;t}$ to be the mean of this block.

Consider a subsample $X^*(\tilde{\tau}) = \{X_{N_{i,j}, \lambda_{ij}} : i = 1, \dots, U_n; j = 1, \dots, k_i\}$ by Algorithm 4.3 conditioned on $\tilde{\tau}$. Thus for each i, j , $N_{i,j} \sim \text{Uniform}(R_{ij})$. Let

$$(6.10) \quad \tilde{\mu}_{n,i} = \sum_{j=1}^{k_i} \frac{f_{ij}}{\hat{f}_i} \bar{X}_{N_{i,j}, \lambda_{ij}}$$

$$(6.11) \quad \tilde{\mu} = \sum_i \hat{f}_i \tilde{\mu}_{n,i}.$$

Then, because of assumptions (A10-A11), $\tilde{\tau}$ satisfies all of the conditions of Theorem 4.4 and hence as a direct consequence of that theorem,

$$(6.12) \quad \sqrt{L_n} \sigma_n^{-1} [\tilde{\mu}_n - \bar{X}_n] \rightarrow N(0, 1).$$

Thus, proving Theorem 4.5 is equivalent to proving

$$(6.13) \quad \rho_M^2[F_n^*(\hat{\tau}^{(n)}), F_n^*(\tilde{\tau}^{(n)})] \rightarrow_p 0.$$

The subsample drawn by Algorithm 4.3 conditional on $\hat{\tau}$ is

$$X^*(\hat{\tau}) = \{X_{N_i, \hat{\lambda}_i} : i = 1, \dots, U_n\}.$$

We couple $X^*(\hat{\tau})$ to $X^*(\tilde{\tau})$ by letting

$$N_i = N_{ij} \quad \text{with probability } f_{ij}/\hat{f}_i, \quad j = 1, \dots, k_i$$

for each $i = 1, \dots, \hat{U}_n$. One can verify that by this construction, $N_i \sim \text{Uniform}(\hat{R}_i)$ as required by Algorithm 4.3. We let $\hat{\mu}_{n,i} = \bar{X}_{N_i; \hat{\lambda}_i}$ and $\hat{\mu}_n = \sum_{i=1}^{\hat{U}_n} f_i \hat{\mu}_{n,i}$.

Let $F_{n,i}^*(\hat{\tau})$ be the distribution of $\sqrt{L}(\hat{\mu}_{n,i} - \bar{X}_{\tau_{i-1}; n_i})$ and $F_{n,i}^*(\tilde{\tau})$ be the distribution of $\sqrt{L}(\tilde{\mu}_{n,i} - \bar{X}_{\tau_{i-1}; n_i})$ conditional on X_1, \dots, X_n and $\hat{\tau}^{(n)}$. Since under this conditioning, $\{\tilde{\mu}_{n,i}\}$ and $\{\hat{\mu}_{n,i}\}$ are each sets of independent random variables, by property (b) of the Mallows' metric,

$$\begin{aligned} \rho_M^2[F_n^*(\hat{\tau}), F_n^*(\tilde{\tau})] &\leq \sum_{i=1}^{\hat{U}_n} \hat{f}_i^2 \rho_M^2[F_{n,i}^*(\hat{\tau}), F_{n,i}^*(\tilde{\tau})] \\ &\leq L \sum_{i=1}^{\hat{U}_n} \hat{f}_i^2 \mathbb{E}^*[\hat{\mu}_{n,i} - \tilde{\mu}_{n,i}]^2 \\ &= L \sum_{i=1}^{\hat{U}_n} \hat{f}_i^2 \text{Var}^*[\hat{\mu}_{n,i} - \tilde{\mu}_{n,i}]. \end{aligned}$$

The second inequality above is due to the definition of Mallows' metric, and the following equality is due to the fact that $\mathbb{E}^*[\hat{\mu}_{n,i} - \tilde{\mu}_{n,i}] = o_p(1)$. Hence, to show (6.13) it is sufficient to show that

$$(6.14) \quad \sum_{i=1}^{\hat{U}_n} \hat{f}_i \hat{\lambda}_i \text{Var}^*(\tilde{\mu}_{n,i} - \hat{\mu}_{n,i}) \rightarrow_p 0.$$

We now study the terms $\hat{\lambda}_i \text{Var}^*(\tilde{\mu}_{n,i} - \hat{\mu}_{n,i})$. For clarity, we first derive an explicit formula for this term when $\hat{U}_n = 1$ and $U_n = 2$, and then generalize to the case $\hat{U}_n > 2$, $U_n > 1$. Under this simple scenario, the estimated segmentation contains no change-points, but there is one true change-point at τ . In this specific case we simplify our notation to let N_1 be uniformly drawn from $\{1, \dots, \tau\}$ and N_2 be uniformly drawn from $\{\tau + 1, \dots, n - L\}$, and let $f_1 = \tau/n$ and $f_2 = 1 - f_1$. Then, define

$$\begin{aligned} A_i^1 &= f_1 L \bar{X}_{i; f_1 L}, \\ A_i^2 &= f_2 L \bar{X}_{i; f_2 L}, \\ B_i^1 &= (1 - f_1) L \bar{X}_{i; (1-f_1)L}, \\ B_i^2 &= (1 - f_2) L \bar{X}_{i; (1-f_2)L}. \end{aligned}$$

The subsample drawn by Algorithm 4.3 assuming $\hat{\tau}$ would simply be one block of length L with starting index N uniformly distributed on $\{1, \dots, n\}$. Since $L/n \rightarrow 0$, we can ignore the effects at the edges of the sequence and thus by the coupling of N to N_1 and N_2 ,

$$\hat{\mu} = [(A_{N_1}^1 + B_{N_1}^1)(1 - J) + (A_{N_2}^2 + B_{N_2}^2)J]/L + o(1),$$

where $J \sim \text{Bernoulli}(f_2)$ is the indicator variable for the event $\{N = N_2\}$. The coupled subsample which assumes knowledge of τ would have mean

$$\tilde{\mu} = (A_{N_1}^1 + A_{N_2}^2)/L + o(1),$$

and hence,

$$\hat{\mu} - \tilde{\mu} = [(B_{N_1}^1 - A_{N_2}^2)(1 - J) + (B_{N_2}^2 - A_{N_1}^1)J]/L.$$

Since

$$\mathbb{V}\text{ar}^*(\hat{\mu} - \tilde{\mu}) = \mathbb{E}^*[\mathbb{V}\text{ar}^*(\hat{\mu} - \tilde{\mu}|J)] + \mathbb{V}\text{ar}^*[\mathbb{E}^*(\hat{\mu} - \tilde{\mu}|J)],$$

we will examine each of the two terms in the above decomposition separately:

$$\begin{aligned} \mathbb{E}^*(\hat{\mu} - \tilde{\mu}|J = 0) &= \mathbb{E}^*(B_{N_1}^1 - A_{N_2}^2)/L \\ &= \frac{1}{L} \left[\sum_{N_1=1}^{\tau-L} \sum_{i=N_1+f_1L}^{N_1+L} X_i - \sum_{N_2=\tau+1}^{n-L} \sum_{i=N_2}^{N_2+f_2L} X_i \right] \\ &= f_2[\bar{X}_{1:\tau} - \bar{X}_{\tau+1:n}] + o(1). \end{aligned}$$

Similarly,

$$\mathbb{E}^*(\hat{\mu} - \tilde{\mu}|J = 1) = f_1[\bar{X}_{1:\tau} - \bar{X}_{\tau+1:n}].$$

It is easy to check by combining the above equations that $\mathbb{E}^*[\hat{\mu} - \tilde{\mu}] = 0$, and hence,

$$\begin{aligned} L\mathbb{V}\text{ar}^*[\mathbb{E}^*(\hat{\mu} - \tilde{\mu})] &= L(f_1f_2^2 + f_2f_1^2)[\bar{X}_{1:\tau} - \bar{X}_{\tau+1:n}]^2 \\ (6.15) \qquad \qquad \qquad &= f_1f_2L[\bar{X}_{1:\tau} - \bar{X}_{\tau+1:n}]^2. \end{aligned}$$

Now consider the term $\mathbb{E}^*[\mathbb{V}\text{ar}^*(\hat{\mu}_n - \tilde{\mu}_n|J)]$.

$$\begin{aligned} \mathbb{V}\text{ar}^*[\hat{\mu} - \tilde{\mu}|J = 0] &= \mathbb{V}\text{ar}^*(B_{N_1}^1 - A_{N_2}^2)/L^2 \\ &= [\mathbb{V}\text{ar}^*(B_{N_1}^1) + \mathbb{V}\text{ar}^*(A_{N_2}^2) - 2\mathbb{C}\text{ov}^*(A_{N_2}^2, B_{N_1}^1)]/L^2. \end{aligned}$$

By independence of N_1 and N_2 , $\mathbb{C}\text{ov}^*(A_{N_2}^2, B_{N_1}^1) = 0$.

$$\begin{aligned} \mathbb{V}\text{ar}^*(B_{N_1}^1) &= f_2L\hat{\sigma}_1^2(f_2L), \\ \mathbb{V}\text{ar}^*(A_{N_2}^2) &= f_2L\hat{\sigma}_2^2(f_2L), \end{aligned}$$

thus

$$\mathbb{V}\text{ar}^*[\hat{\mu} - \tilde{\mu}|J = 0] = f_2[\hat{\sigma}_1^2(f_2L) + \hat{\sigma}_2^2(f_2L)]/L,$$

and similarly,

$$\mathbb{V}\text{ar}^*[\hat{\mu} - \tilde{\mu}|J = 1] = f_1[\hat{\sigma}_1^2(f_1L) + \hat{\sigma}_2^2(f_1L)]/L,$$

and therefore

$$(6.16) \quad L\mathbb{E}^*[\mathbb{V}\text{ar}^*(\hat{\mu} - \tilde{\mu}|J)] = f_1f_2[\hat{\sigma}_1^2(f_1L) + \hat{\sigma}_1^2(f_2L) + \hat{\sigma}_2^2(f_1L) + \hat{\sigma}_2^2(f_2L)].$$

Now we generalize (6.15) and (6.16) to the case where $U_n > 2$, $\hat{U}_n = 1$. For $j = 1, \dots, U_n$, let N_j be uniformly distributed on $\{\tau_{j-1} + 1, \dots, \tau_j\}$ and

$$\begin{aligned} A_{N_i}^i &= f_i L \bar{X}_{N_i; f_i L} \\ B_{N_i}^i &= (1 - f_i) L \bar{X}_{N_i; (1-f_i)L}. \end{aligned}$$

Then

$$\hat{\mu} = \frac{1}{L} \sum_{i=1}^{U_n} I_{\{J=i\}} (A_{N_i}^i + B_{N_i}^i),$$

where J is multinomial with $P(J = i) = f_i$, $i = 1, \dots, U_n$. The corresponding coupled statistic is

$$\tilde{\mu} = \frac{1}{L} \sum_{i=1}^{U_n} A_{N_i}^i,$$

and hence

$$\hat{\mu} - \tilde{\mu} = \frac{1}{L} \sum_{i=1}^{U_n} I_{\{J=i\}} \left(B_{N_i}^i - \sum_{j \neq i} A_{N_j}^j \right).$$

Similar to the $U_n = 2$ case, by computing the first and second moments of $A_{N_i}^i$ and $B_{N_i}^i$ conditional on the observed sequence, we have the corresponding equations to (6.15) and (6.16) for the $U_n > 2$, $\hat{U}_n = 1$ case:

$$\begin{aligned} L\mathbb{V}\text{ar}^*[\mathbb{E}^*(\hat{\mu} - \tilde{\mu}|J)] &= \sum_{i=1}^{U_n} f_i \mathbb{E}^*(\hat{\mu} - \tilde{\mu}|J = i) \\ &= \sum_{i=1}^{U_n} f_i (1 - f_i)^2 L \left[\sum_{j \neq i} \frac{f_j}{1 - f_i} (\bar{X}_{R_j} - \bar{X}_{R_i}) \right]^2 + o_p(1). \end{aligned} \tag{6.17}$$

$$L\mathbb{E}^*[\text{Var}^*(\hat{\mu} - \tilde{\mu}|J)] = \sum_{i=1}^{U_n} f_i(1-f_i) \left[\hat{\sigma}_i^2((1-f_i)L) + \sum_{j \neq i} \frac{f_j}{1-f_i} \hat{\sigma}_j^2(f_i L) \right] + o_p(1) \quad (6.18)$$

Finally, generalizing (6.17) and (6.18) to the $\hat{U}_n > 1$ case, we have, for each region i , an independent variable J_i defined as for J above, taking values in $1, \dots, k_i$, giving us:

$$\hat{\lambda}_i \text{Var}^*[\mathbb{E}^*(\hat{\mu}_{n,i} - \tilde{\mu}_{n,i}|J_i)] = \hat{\lambda}_i \sum_{j=1}^{k_i} \frac{f_{ij}}{\hat{f}_i} \left(1 - \frac{f_{ij}}{\hat{f}_i}\right)^2 \left[\sum_{k \neq j} \frac{f_{ik}}{\hat{f}_i - f_{ij}} (\bar{X}_{R_{ik}} - \bar{X}_{R_{ij}}) \right]^2,$$

and

$$\begin{aligned} & \hat{\lambda}_i \mathbb{E}^*[\text{Var}^*(\hat{\mu}_{n,i} - \tilde{\mu}_{n,i}|J_i)] \\ &= \sum_{j=1}^{k_i} \frac{f_{ij}}{\hat{f}_i} \left(1 - \frac{f_{ij}}{\hat{f}_i}\right) \left[\hat{\sigma}_{ij}^2((\hat{f}_i - f_{ij})L) + \sum_{k \neq j} \frac{f_{ik}}{\hat{f}_i - f_{ij}} \hat{\sigma}_{ik}^2(f_{ij}L) \right], \end{aligned}$$

where for any $R \subseteq \{1, \dots, n\}$, $\bar{X}_R = \sum_{i \in R} X_i / |R|$ and for any i, j , and l $\hat{\sigma}_{ij}^2(l) = l \sum_{i \in R_{ij}} (\bar{X}_{i;l} - \bar{X}_{R_{ij}})^2 / n_{ij}$. By assumptions A5 and A6, and summing the above quantities over $i = 1, \dots, \hat{U}_n$, we have

$$\begin{aligned} \sum_{i=1}^{\hat{U}_n} f_i \hat{\lambda}_i \text{Var}^*(\hat{\mu}_{n,i} - \tilde{\mu}_{n,i}) &\leq C \sum_{i=1}^{\hat{U}_n} f_i \sum_{j=1}^{k_i} \left[\hat{\lambda}_i \frac{f_{ij}}{\hat{f}_i} \left(1 - \frac{f_{ij}}{\hat{f}_i}\right)^2 + \frac{f_{ij}}{\hat{f}_i} \left(1 - \frac{f_{ij}}{\hat{f}_i}\right) \right] \\ &\leq C \sum_{i=1}^{\hat{U}_n} \sum_{j=1}^{k_i} [L \min(f_{ij}, \hat{f}_i - f_{ij}) + \min(\hat{f}_i, \hat{f}_i - f_{ij})] \\ &\leq 2C \frac{L+1}{n} \sum_{i=1}^{U_n} \min_{1 \leq j \leq \hat{U}_n} |\tau_i - \hat{\tau}_j|. \end{aligned}$$

By assumption A11, the above converges in probability to 0, and thus (6.14) holds. \square

6.5. *Proof of Theorem 4.10.* We begin with the known stationary case. Evidently, we need only apply the delta method to

$$(\bar{I}\bar{J}_n, \bar{I}_n, \bar{J}_n)$$

where $\bar{I}\bar{J}_n \equiv \frac{1}{n} \sum_{k=1}^n I_k J_k$ and $\bar{I}_n \equiv \frac{1}{n} \sum_{k=1}^n I_k$, \bar{J}_n as defined. The result follows from Theorem 3.1. We continue with part 2. We require first a series of couplings.

Since we are working under H_0 , $\{I_k\}$ and $\{J_k\}$ are independent. Conditional on K_1, K_2 we construct for each M a joint distribution of (I_n^c, J_n^c) with the following properties.

1. Let $S_{11} \equiv \{k : k < K_1 < K_2, K_2 > K_1 + M + L\}$, $S_{12} = \{k : k > K_2 > K_1 + M + L\}$ and S_{21}, S_{22} with K_1, K_2 interchanged. Let $S_a = \cup_{b=1}^2 S_{ab}$. If $|K_1 - K_2| \leq M + L$ let $S_a \equiv S_b \equiv \phi$. Else let (I_k^c, J_k^c) , $k \leftarrow S_1$ be independent of (I_k^c, J_k^c) , $k \in S_2$.
2. The marginal distributions of $\{I_k^c\}$ and $\{J_k^c\}$ and $\{I_k\}$ and $\{J_k\}$ are the same respectively.
3. The conditional distribution of $\{(I_k^c, J_k^c) : k \notin S_1 \cup S_2\}$ given $\{(I_k^c, J_k^c) = (\varepsilon_{k1}, \varepsilon_{k2}) : k \in S_1 \cup S_2\}$ and $\{(I_k, J_k) : 1 \leq k \leq n\}$ is that of $\{(I_k, J_k) : k \notin S_1 \cup S_2\}$ given $(I_k, J_k) = (\varepsilon_{k1}, \varepsilon_{k2}) : k \leftarrow S_1 \cup S_2$. Since $\{I_k\}, \{J_k\}$ are independent $\{I_k^c\}$ and $\{J_k^c\}$ are also.
4. $P[I_k \neq I_k^c \text{ or } J_k \neq J_k^c \text{ for any } k \in S_1 \cup S_2] \leq m(M)$.
This construction is possible by Strassen's (1965) Theorem.
5. If $|K_1 - K_2| \leq M + L$, then make $\{I_k^c\}\{J_k^c\}$ be independent of each other and of $\{I_k, J_k\}$.

If we adjoin (K_1, K_2) we see that we have constructed a coupling which depends only on M and which has properties 1–4 given K_1, K_2 . We now define statistics $\bar{I}\bar{J}^{*1c}$, \bar{J}^{*1c} etc. (suppressing dependence on n, L, M) defined on (I^c, J^c) and bootstrap versions $\bar{I}\bar{J}^{*1c}$ etc. as well.

We want to argue generically that if $|g| \leq c < \infty$ and the assumptions of the theorem hold and

$$S_L \equiv \frac{1}{L} \sum_{l=1}^L [g(I_{K_1+l}, J_{K_2+l}) + g(I_{K_2+l}, J_{K_1+l})]$$

$$S_L^c = \frac{1}{L} \sum_{l=1}^L [g(I_{K_1+l}, J_{K_2+l}^c) + g(I_{K_2+l}, J_{K_1+l}^c)]$$

then,

$$|S_L - S_L^c| = o_P(L^{-\frac{1}{2}})$$

as $L, n \rightarrow \infty$, $L = o(n)$, $M = o(n)$ uniformly in g as above. But,

$$(6.19) \quad P[|S_L - S_L^c| \geq \varepsilon L^{-\frac{1}{2}}, \quad |K_1 - K_2| \leq M + L] \leq \frac{M + L}{n}$$

$$(6.20) \quad P[|S_L - S_L^c| \geq \varepsilon L^{-\frac{1}{2}}, \quad |K_1 - K_2| > M + L] \leq m(M)$$

Plugging in $g(I, J) = IJ, I, J$ we see that in the stationary case,

$$T_{nL}^* = T_{nL}^{*c} + o_P(L^{-\frac{1}{2}})$$

where T_{nL}^{*c} is defined by replacing I, J by I^c, J^c throughout in Algorithm 4.9. Now we can appeal to Theorem 4.1 to establish part 2 of Theorem 4.10 once we note that for T_{nL}^c given $K_1 - K_2 > M + L$, $\bar{I}\bar{J}^{*1c}, \bar{I}\bar{J}^{*2c}, \bar{I}^{*1c}, \bar{I}^{*2c}, \bar{J}^{*1c}, \bar{J}^{*2c}$ are mutually independent so that the asymptotic covariance of the two components of T_{nL}^{*c} and J_{nL}^{*c} is 0.

The generalization to the case of a fixed known segmentation satisfying the assumptions of Theorem 4.3 is straightforward by coupling the component corresponding to each segment and computing the variance of the difference using the independence of the bootstrap statistics for different segments. Finally the case of unknown segmentation can be dealt with by reducing the unknown segmentation case to the known one using Theorem 4.4 and then applying the result for known segmentation, Theorem 4.5. The result follows. \square

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