# Eigenvalue Distribution of Time and Frequency Limiting 

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## 1. Introduction

Let $S$ and $T$ represent sets, each the union of a finite number ( $\mu$ and $v$, respectively) of fixed disjoint intervals. We shall be interested in the behavior, as the parameter $c \rightarrow \infty$, of the eigenvalues $\left\{\lambda_{k}(c)\right\}$ of the integral equation

$$
\begin{equation*}
\int_{y \in c S} h_{T}(x-y) \phi_{k}(y) d y=\lambda_{k}(c) \phi_{k}(x), \quad x \in c S \tag{1}
\end{equation*}
$$

in which the Fourier transform of $(2 \pi)^{1 / 2} h_{T}$ is $\chi_{T}(\omega)$, the characteristic function of $T$ :

$$
h_{T}(s) \equiv \frac{1}{2 \pi} \int_{\omega \in T} e^{i s s_{1}} d \omega
$$

The operator in (1) consists of restricting $\phi_{k}(y)$ to the set $c S$, restricting the Fourier transform of the function so obtained to the set $T$, and viewing the result again on $c S$. We can therefore represent it compactly as $A_{c} \equiv P(c S) Q(T) P(c S)$, where $P(\Omega)$ and $Q(\Delta)$ represent orthogonal projections in the Hilbert space $L^{2}(-\infty, \infty)$ onto the subspaces of those functions which vanish outside of $\Omega$, and those whose Fourier transform vanishes outside of $\Delta$. The interest of the problem lies in the fact that the eigenvalues are useful in describing the geometry of these subspaces, and thereby provide information about how the energy of a function can be distributed over time and frequency. From a comparison of $\sum_{k} \lambda_{k}(c)$ and
$\sum_{k} \lambda_{k}^{2}(c)$ it is easy to see that the eigenvalues (when arranged in nonincreasing order) are very close to 1 , then very close to 0 , the transition occurring in an interval of values of $k$ which is centered at $k=c \mathbf{m}(S) \mathbf{m}(T) / 2 \pi$, with $\mathbf{m}(\cdot)$ Lebesgue measure, and grows in width at the rate of only $\log c$. This fact has found application to certain questions concerning sampling $[4]$, to an extension of Szegö's theorem [5], and to the eigenvalues of Hankel matrices [7]. In a paper [6] devoted to asymptotic expressions for the eigenvalues and for the corresponding eigenfunctions, in the case that $S$ and $T$ are single intervals, Slepian conjectured that when $n=(2 \pi)^{-1} \mathbf{m}(S) \mathbf{m}(T) c+\pi^{-2} b \log c, \lambda_{n}(c) \rightarrow\left(1+e^{b}\right)^{-1}$. Here we prove this result, by showing it to stem from like behavior of Hankel operators [7], and generalize it to sets $S$ and $T$ which are finite unions of intervals. Specifically, letting $N(A, \alpha)$ denote the number of eigenvalues of an operator $A$ which exceed $\alpha$, we will show that, for $0<\alpha<1$,

$$
\begin{equation*}
N\left(A_{c}, \alpha\right)=\frac{\mathbf{m}(S) \mathbf{m}(T)}{2 \pi} c+\left(\frac{\mu v}{\pi^{2}} \log \frac{1-\alpha}{\alpha}\right) \log c+o(\log c) \tag{2}
\end{equation*}
$$

A traditional method for describing the distribution of eigenvalues for a difference kernel relies on finding the trace of iterates of the operator. The fact that the kernel of (1) is not absolutely integrable presents a severe impediment to that approach. We circumvent the difficulty here by considering, instead, the operator $A_{c}\left(1-A_{c}\right)$ whose eigenvalues, $\lambda_{k}(c)\left[1-\lambda_{k}(c)\right]$, are large precisely when $\lambda_{k}(c)$ lies in the intermediate range which we wish to study. We will be able to compute the trace of $\left[A_{c}\left(I-A_{c}\right)\right]^{n}$ and of $\left.A_{c} \mid A_{c}\left(I-A_{c}\right)\right]^{n}$. The resulting estimates, combined with the fact that, from (1),

$$
\begin{equation*}
\operatorname{tr} A_{c}=\int_{y \in c S} h_{T}(0) d y=\frac{\mathbf{m}(c S) \mathbf{m}(T)}{2 \pi} \tag{3}
\end{equation*}
$$

are sufficient to establish (2).
The discrete analogue of our results for $S$ a single interval, concerning the eigenvalues of Toeplitz matrices, follows easily from a result of Basor [1] on the asymptotics of their determinants. Her methods, although more complicated, also give the next term in the asymptotic expansion. They could undoubtedly be adapted to the present situation.

## 2. Preliminaries

The main tool we employ is that of the trace norm $\|A\|_{1}$ of a compact operator $A$, defined as $\sum_{i=1}^{\infty} s_{i}(A)$, where $\left\{s_{i}^{2}\right\}$ are the eigenvalues of $A^{*} A$. This norm is finite whenever $A$ is the product $C_{1} C_{2}$ of two operators of

Hilbert-Schmidt type, i.e., ones for which the eigenvalues of $C_{i}^{*} C_{i}$ are summable. We will have occasion to use the following basic properties of the trace norm |3, Chap. 3|:

$$
\begin{equation*}
\|A+B\|_{1} \leqslant\|A\|_{1}+\|B\|_{1} \tag{4}
\end{equation*}
$$

If $B$ is a bounded operator, with (ordinary) norm $\|B\|$,

$$
\begin{align*}
& \|B A\|_{1} \leqslant\|B\|\|A\|_{1}  \tag{5}\\
& \|A B\|_{1} \leqslant\|B\|\|A\|_{1} .
\end{align*}
$$

For any orthonormal basis $\left\{\phi_{i}\right\}, \operatorname{tr} A \equiv \sum_{i=1}^{\infty}\left(A \phi_{i}, \phi_{i}\right)$ is independent of the choice of basis, and

$$
\begin{equation*}
|\operatorname{tr} A| \leqslant\|A\|_{1} \tag{6}
\end{equation*}
$$

Now let us turn specifically to the sorts of operators we will consider. The Hilbert space is $L^{2}(-\infty, \infty)$, in which the Fourier transform

$$
F \phi \equiv \frac{1}{(2 \pi)^{1 / 2}} \int_{-\infty}^{\infty} \phi(t) e^{-i \omega t} d t
$$

defines a unitary transformation. Let $P(S) \phi$ denote the orthogonal projection of $\phi$ onto the subspace of functions vanishing outside the set $S$; explicitly,

$$
P(S) \phi=\chi_{s}(t) \phi(t)
$$

Analogously, let $Q(T) \phi=F^{-1} P(T) F \phi$ denote the orthogonal projection of $\phi$ onto the subspace of functions whose Fourier transform vanishes outside $T$. $Q(T) \phi$ can be represented in the form

$$
Q(T) \phi=\int h_{T}(x-y) \phi(y) d y
$$

where

$$
F h_{T}=\frac{1}{(2 \pi)^{1 / 2}} \chi_{T}(\omega)
$$

If $S$ or $T$ is an interval $[a, b]$, we will often write $P(a, b)$ or $Q(a, b)$ in place of $P(S)$ or $Q(T)$. In this situation,

$$
h_{T}(s)=\frac{e^{i b s}-e^{i a s}}{2 \pi i s}
$$

Moreover, in an operator of the form $P(J) Q(a, b) P(K)$, with $J$ and $K$ intervals having disjoint interiors, we can pass to the limit as $a \rightarrow-\infty$ or $b \rightarrow \infty$, obtaining an integral operator with kernel $\chi_{K}(y) \chi_{J}(x)\left(e^{i b(x-y)} / 2 \pi i(x-y)\right) \quad$ or $\quad \chi_{K}(y) \chi_{J}(x)\left(-e^{i a(x-y)} / 2 \pi i(x-y)\right)$ for $P(J) Q(-\infty, b) P(K)$ or $P(J) Q(a, \infty) P(K)$, respectively.

Suppose now that $J, K, L, M, N$ are intervals, finite or semi-infinite, and consider the operator $R=P(J) Q(M) P(K) Q(N) P(L)$. Let us say that $R_{1}$ and $R_{2}$ are unitarily equivalent, written $R_{1} \sim R_{2}$, if $R_{1}=U R_{2} U^{-1}$ for some unitary transformation $U$. By definition, both $\operatorname{tr} R$ and $\|R\|_{1}$ are unchanged under this equivalence. By choosing $U$ to be, successively, rescaling $\left[f(t) \rightarrow|\gamma|^{-1 / 2} f(\gamma t)\right]$, translation $[f(t) \rightarrow f(t+\gamma)]$, shift in frequency $\left[f(t) \rightarrow e^{i v t} f(t)\right]$, complex conjugation $[f(t) \rightarrow \overline{f(t)}]$, and the Fourier transform $[f \rightarrow F f]$, we find, for any scalar $\gamma_{j}$,

$$
\begin{align*}
R & \sim P\left(J / \gamma_{1}\right) Q\left(\gamma_{1} M\right) P\left(K / \gamma_{1}\right) Q\left(\gamma_{1} N\right) P\left(L / \gamma_{1}\right)  \tag{7}\\
& \sim P\left(J+\gamma_{2}\right) Q\left(M+\gamma_{3}\right) P\left(K+\gamma_{2}\right) Q\left(N+\gamma_{3}\right) P\left(L+\gamma_{2}\right)  \tag{8}\\
& \sim P(J) Q(-M) P(K) Q(-N) P(L)  \tag{9}\\
& \sim Q(J) P(M) Q(K) P(N) Q(L) . \tag{10}
\end{align*}
$$

Moreover, if $J$ and $K$ intersect in a set of measure 0 , and $A$ and $B$ are complementary, then since $Q(A)=I-Q(B)$, we have

$$
\begin{equation*}
P(J) Q(A) P(K)=-P(J) Q(B) P(K) \tag{11}
\end{equation*}
$$

We will be interested in conditions under which $\|R\|_{1}$ is uniformly bounded, independently of the choice of the intervals $J, M, K, N, L$ within a certain class, a property we will denote by $R=O(1)$. To discuss this, let us observe that an integral operator in $L^{2}$ whose kernel is $p(x) q(y)$ has rank 1 and trace norm $\|p\|\|q\|$, where $\|f\|^{2}=\int_{-\infty}^{\infty}|f(x)|^{2} d x$. Consequently, by (6) and (4), if an operator $A$ is defined by a kernel $\int p(x, z) q(y, z) d z$,

$$
\begin{equation*}
|\operatorname{tr} A| \leqslant\|A\|_{1} \leqslant \int\|p(\cdot, z)\|\|q(\cdot, z)\| d z \tag{12}
\end{equation*}
$$

As an immediate consequence we obtain the following criterion.
Lemma. Suppose that the interiors of $J$ and $L$ lie in the complement of $K$. Then $R=O(1)$ under each of the conditions:
(L1) If $\mathrm{m}(K)$ is uniformly bounded, and either $J$ and $L$ are on opposite sides of $K$, or one of $J, L$ is uniformly separated from $K$.
(L2) If $K=(-\infty, 0], \mathrm{m}(J)$ or $\mathrm{m}(L)$ is uniformly bounded, and $J$ or $L$ is uniformly separated from 0 .
(L3) If $K=(-\infty, 0]$, the finite endpoints of $M$ and $N$ are uniformly separated from one another, and either $J$ or $L$ is uniformly separated from 0.
(L4) If $K=(-\infty, 0], \mathbf{m}(N)$ or $\mathbf{m}(M)$ is uniformly bounded, and $\mathbf{m}(J)$ or $\mathbf{m}(L)$ is uniformly bounded.

Proof. In L1 we may, by (8) and (7), suppose that $K=[0,1]$. The operator $R$ is then given by a kernel of the form

$$
\int_{0}^{1} d u \frac{e^{i n_{2}(u-y)}-e^{i n_{1}(u-y)}}{2 \pi i(u-y)} \frac{e^{i m_{2}(x-u)}-e^{i m_{1}(x-u)}}{2 \pi i(x-u)}
$$

with $x \in J, y \in L$; if $M$ or $N$ is semi-infinite, the exponentials corresponding to the infinite endpoint are omitted. Let us write the integrand as a sum of terms of the form

$$
\frac{e^{i n_{j}(u-y)}}{u-y} \frac{e^{i m_{k}(x-u)}}{x-u}
$$

If $J \subset[1, \infty)$ and $L \subset(-\infty, 0]$, we find from (12), with $\gamma$ an appropriate constant,

$$
\|R\|_{1} \leqslant \gamma \int_{0}^{1}(1-u)^{-1 / 2} u^{-1 / 2} d u=O(1)
$$

Likewise, if $J$ and $L$ are on the same side of 0 or 1 , but $L$ is uniformly separated from $K$, say $L \subset[1+\delta, \infty)$, then

$$
\|R\|_{1} \leqslant \gamma \int_{0}^{1}(1-u)^{-1 / 2}(1+\delta-u)^{-1 / 2} d u=O(1)
$$

For L2, arguing analogously with $J=[a, b], 0 \leqslant a<b$, and $L \subset[\delta, \infty]$, we find

$$
\|R\|_{1} \leqslant \gamma \int_{-\infty}^{0} d u\left[(a-u)^{-1 / 2}-(b-u)^{-1 / 2}\right](\delta-u)^{-1 / 2}=O(1)
$$

In L3, $R$ is given by a sum of at most four kernels of the form $e^{i\left(\theta_{1} x-\theta_{2} y\right)} \gamma \int_{-\infty}^{0} d u\left(e^{i a u} /(u-x)(u-y)\right)$ with $|\alpha|$ bounded uniformly away from $0, x \in J, y \in L$. Integration by parts converts this to

$$
\gamma e^{i\left(\theta_{1} x-\theta_{2} y\right)}\left[\frac{1}{i a} \frac{1}{x y}-\frac{1}{i \alpha} \int_{-\infty}^{0} d u \frac{e^{i \alpha u}}{(u-x)(u-y)}\left(\frac{1}{u-x}+\frac{1}{u-y}\right)\right]
$$

If both $J$ and $L$ are uniformly separated from 0 , say $J, L \subset[\delta, \infty)$, the first component, of rank 1 , has trace norm $\gamma / a \delta$, while the second has uniformly
bounded trace norm by (12). If only one of $J, L$ is separated from 0 , say $J=[\delta, \infty), L=[0, \infty)$, then by (8)

$$
\begin{aligned}
& P(J) Q(M) P(K) Q(N) P(L) \sim P(2 \delta, \infty) Q(M) P(-\infty, \delta) Q(N) P(\delta, \infty) \\
&= P(2 \delta, \infty) Q(M) P(-\infty, 0) Q(N) P(\delta, \infty) \\
& \quad+P(2 \delta, \infty) Q(M) P(0, \delta) Q(N) P(\delta, \infty)
\end{aligned}
$$

Now the first operator on the right is $O(1)$ by the preceding reasoning, while the second is $O(1)$ by L1.

In L4 we may again suppose $N=[0,1]$, whereupon the kernel is

$$
\int_{-\infty}^{0} d u \frac{e^{i(u-y)}-1}{2 \pi i(u-y)} \frac{e^{i m_{2}(x-u)}-e^{i m_{1}(x-u)}}{2 \pi i(x-u)}
$$

Since $\left|e^{-i s}-1\right| /|s| \leqslant \gamma /(1+|s|)$ we see that the integrand is bounded by $\gamma /(1+y-u)(x-u)$. If either $J$ or $L \subset[0, \infty)$ is uniformly bounded in length, boundedness of $\|R\|_{1}$ again follows from (12). This completes the proof of the lemma.

## 3. Results

We now return to (1), restricting for the moment to the case that $S$ and $T$ consist of a single interval. By (7) and (8) we can renormalize so that $T^{\prime}=[0,1]$ and $S^{\prime}=[0, \mathbf{m}(S) \mathbf{m}(T)]$, whereupon $c S^{\prime}=[0, c \mathbf{m}(S) \mathbf{m}(T)]$. On rescaling $c$ by $c^{\prime}=c \mathbf{m}(S) \mathbf{m}(T)$, we see that to establish (2) it is sufficient to prove it for the case $\mathbf{m}(S)=\mathbf{m}(T)=1$. Accordingly, let $A_{c} \equiv$ $P(0, c) Q(0,1) P(0, c)$.

Theorem 1. For $0<\alpha<1$,

$$
N\left(A_{c}, \alpha\right)=\stackrel{c}{2 \pi}+\left(\frac{1}{\pi^{2}} \log \frac{1-\alpha}{\alpha}\right) \log c+o(\log c)
$$

Proof. Since $P$ and $Q$ are orthogonal projections, $P^{2}=P$ and $Q^{2}=Q$. Thus $A_{c}^{2}=P(0, c) Q(0,1) P(0, c) Q(0,1) P(0, c)$, hence

$$
\begin{aligned}
A_{c}-A_{c}^{2}= & P(0, c) Q(0,1)[P(-\infty, 0)+P(c, \infty)] Q(0,1) P(0, c) \\
= & P(0, c) Q(0,1) P(-\infty, 0) Q(0,1) P(0, c) \\
& +P(0, c) Q(0,1) P(c, \infty) Q(0,1) P(0, c)
\end{aligned}
$$

The product of the two operators on the right contains the factor $P(-\infty, 0) Q(0,1) P(0, c) Q(0,1) P(c, \infty)$, which by (7) is equivalent to $P(-\infty, 0) Q(0, c) P(0,1) Q(0, c) P(1, \infty)$, and so is $O(1)$ by L1. The remaining factors, consisting of projections, have norm bounded by 1 , so the product is $O(1)$ by (5). Thus for each $n$,

$$
\begin{align*}
\left.\mid A_{c}\left(I-A_{c}\right)\right]^{n}= & {[P(0, c) Q(0,1) P(-\infty, 0) Q(0,1) P(0, c)]^{n} } \\
& +[P(0, c) Q(0,1) P(c, \infty) Q(0,1) P(0, c)]^{n}+O(1) \tag{13}
\end{align*}
$$

If we now replace the first $P(0, c)$ by $P(1, c)$ in $P(0, c) Q(0,1) P(-\infty, 0)$ $Q(0,1) P(0, c)$, the difference is $P(0, c) Q(0,1) P(-\infty, 0) Q(0,1) P(0,1)$, which is $O(1)$ by L4. Thus

$$
\begin{aligned}
& P(0, c) Q(0,1) P(-\infty, 0) Q(0,1) P(0, c) \\
& \quad=P(1, c) Q(0,1) P(-\infty, 0) Q(0,1) P(1, c)+O(1)
\end{aligned}
$$

Continuing, let us write $Q(0,1)=Q(0, \infty)-Q(1, \infty)$ in the right-hand operator and expand. The terms involving both $Q(0, \infty)$ and $Q(1, \infty)$, for example, $P(1, c) Q(0, \infty) P(-\infty, 0) Q(1, \infty) P(1, c)$, are each $O(1)$ by L3. Thus

$$
\begin{aligned}
P(0, c) & Q(0,1) P(-\infty, 0) Q(0,1) P(0, c) \\
= & P(1, c) Q(0, \infty) P(-\infty, 0) Q(0, \infty) P(1, c) \\
& +P(1, c) Q(1, \infty) P(-\infty, 0) Q(1, \infty) P(1, c)+O(1)
\end{aligned}
$$

The product of the two operators on the right contains the factor $P(-\infty, 0)$ $Q(0, \infty) P(1, c) Q(1, \infty) P(-\infty, 0)$, which is equivalent to $P(1, \infty) Q(-\infty, 0)$ $P(-c, 0) Q(-\infty,-1) P(1, \infty)$ by (8) and (7), and so is $O(1)$ by L3. The second operator on the right-hand side of (13) is unitarily equivalent to the first; on applying the analogous chain of argument to it, we finally obtain

$$
\begin{align*}
{\left[A_{c}(I-\right.} & \left.\left.A_{c}\right)\right]^{n} \\
= & {[P(1, c) Q(0, \infty) P(-\infty, 0) Q(0, \infty) P(1, c)]^{n} } \\
& +[P(1, c) Q(1, \infty) P(-\infty, 0) Q(1, \infty) P(1, c)]^{n}  \tag{14}\\
& +[P(0, c-1) Q(0, \infty) P(c, \infty) Q(0, \infty) P(0, c-1)]^{n} \\
& +[P(0, c-1) Q(1, \infty) P(c, \infty) Q(1, \infty) P(0, c-1)]^{n}+O(1)
\end{align*}
$$

Now by suitably applying (7)-(9) we see that each of the operators on the right-hand side of (14) is unitarily equivalent to the first. Thus, setting

$$
K_{c} \equiv P(1, c) Q(0, \infty) P(-\infty, 0) Q(0, \infty) P(1, c)
$$

we find, by (6),

$$
\begin{equation*}
\operatorname{tr}\left[A_{c}\left(I-A_{c}\right)\right]^{n}=4 \operatorname{tr} K_{c}^{n}+O(1) \tag{15}
\end{equation*}
$$

Next, let us apply $A_{c}$ to both sides of (14), obtaining $A_{c}\left[A_{c}\left(I-A_{c}\right)\right]^{n}$ as the sum of four operators plus $O(1)$. By (8) with $\gamma_{2}=0, \gamma_{3}=1$, (9), and (11) we see that the second, $P(0, c) Q(0,1)[P(1, c) Q(1, \infty) P(-\infty, 0) Q(1, \infty)$ $P(1, c)]^{n}$, is unitarily equivalent to the first, and by analogous arguments, so are the third and fourth. We conclude that

$$
\begin{equation*}
\operatorname{tr} A_{c}\left[A_{c}\left(I-A_{c}\right)\right]^{n}=4 \operatorname{tr} P(0, c) Q(0,1) K_{c}^{n}+O(1) \tag{16}
\end{equation*}
$$

Further, on replacing $Q(0,1)$ in (16) by $Q(0, \infty)$, the difference contains as factor the operator $Q(1, \infty) P(1, c) Q(0, \infty) P(-\infty, 0) Q(0, \infty)=$ $-Q(1, \infty) P(1, c) Q(-\infty, 0) P(-\infty, 0) Q(0, \infty)$ by (11). But by (10) the latter is equivalent to $-P(1, \infty) Q(1, c) P(-\infty, 0) Q(-\infty, 0) P(0, \infty)$, which is $O(1)$ by L 3 . It follows that

$$
\begin{equation*}
\operatorname{tr} A_{c}\left[A_{c}\left(I-A_{c}\right)\right]^{n}=4 \operatorname{tr} P(0, c) Q(0, \infty) K_{c}^{n}+O(1) \tag{17}
\end{equation*}
$$

Finally, by (11), $K_{c}=P(1, c) Q(-\infty, 0) P(-\infty, 0) Q(-\infty, 0) P(1, c)$ so that, by (9),

$$
\operatorname{tr} P(0, c) Q(0, \infty) K_{c}^{n}=\operatorname{tr} P(0, c) Q(-\infty, 0) K_{c}^{n}
$$

But as $Q(0, \infty)+Q(-\infty, 0)=I$, this shows that

$$
\operatorname{tr} P(0, c) Q(0, \infty) K_{c}^{n}=\frac{1}{2} \operatorname{tr} P(0, c) K_{c}^{n}=\frac{1}{2} \operatorname{tr} K_{c}^{n}
$$

so that by (17)

$$
\begin{equation*}
\operatorname{tr} A_{c}\left[A_{c}\left(I-A_{c}\right)\right]^{n}=2 \operatorname{tr} K_{c}^{n}+O(1)=\frac{1}{2} \operatorname{tr}\left[A_{c}\left(I-A_{c}\right)\right]^{n}+O(1) \tag{18}
\end{equation*}
$$

Now $K_{c}$ is given by the kernel $\left(1 / 4 \pi^{2}\right) \int_{0}^{\infty} d u /(u+x)(u+y)$, with $1 \leqslant x$, $y \leqslant c$, and so resembles a Hankel operator. Proceeding then as in [7], we apply the change of variable $x=e^{2 \sigma}, y=e^{2 t}, u=e^{2 \zeta}$, which transforms $K_{c}$ unitarily into the integral operator on $L^{2}[0,(\log c) / 2]$ defined by the difference kernel $k(\tau-\sigma)$, with

$$
k(s)=\frac{1}{4 \pi^{2}} \int_{-\infty}^{\infty} \operatorname{sech} r \operatorname{sech}(s-r) d r
$$

a rapidly decreasing function of $s$. In accordance with Szegö's theorem
[ 2,5$]$, the eigenvalue distribution of $K_{c}$ can be determined from the Fourier transform of $k$. Specifically, since

$$
\begin{aligned}
\int_{-\infty}^{\infty} d v e^{-i v s} \operatorname{sech} v & =\pi \operatorname{sech} \frac{\pi}{2} s \\
N\left(K_{c}, \alpha\right) & =\frac{\log c}{2} \frac{1}{2 \pi} \mathbf{m}\left\{s \left\lvert\, \operatorname{sech}^{2} \frac{\pi}{2} s>4 \alpha\right.\right\}+o(\log c) \\
& =\frac{\log c}{4 \pi} \cdot 2 \cdot \frac{2}{\pi} \operatorname{sech}^{-1}(4 \alpha)^{1 / 2}+o(\log c)
\end{aligned}
$$

Thus, by definition of $N$,

$$
\begin{aligned}
\operatorname{tr} K_{c}^{n} & =\int_{0+}^{1 / 4} x^{n} d_{x}\left[-N\left(K_{c}, x\right)\right] \\
& =\frac{1}{2 \pi^{2}} \log c \int_{0+}^{1 / 4} x^{n} \frac{d x}{x(1-4 x)^{1 / 2}}+o(\log c)
\end{aligned}
$$

and letting $x=t(1-t), 0<t \leqslant \frac{1}{2}$, we find

$$
\begin{align*}
\operatorname{tr} K_{c}^{n} & =\frac{1}{2 \pi^{2}} \log c \int_{0}^{1 / 2}[t(1-t)]^{n} \frac{d t}{t(1-t)}+o(\log c)  \tag{19}\\
& =\frac{1}{4 \pi^{2}} \log c \int_{0}^{1}[t(1-t)]^{n} \frac{d t}{t(1-t)}+o(\log c)
\end{align*}
$$

We observe also that, by symmetry of the integrand in this expression,

$$
\int_{0}^{1} t[t(1-t)]^{n} \frac{d t}{t(1-t)}=\int_{0}^{1}(1-t)[t(1-t)]^{n} \frac{d t}{t(1-t)}
$$

so that

$$
\begin{equation*}
\int_{0}^{1} t[t(1-t)]^{n} \frac{d t}{t(1-t)}=\frac{1}{2} \int_{0}^{1}[t(1-t)]^{n} \frac{d t}{t(1-t)} \tag{20}
\end{equation*}
$$

It is now easy to convert (15) and (18) into information concerning $N\left(A_{c}, \alpha\right)$. For by definition of $N$,

$$
\begin{aligned}
\operatorname{tr}\left[A_{c}\left(I-A_{c}\right)\right]^{n} & =\int_{0+}^{1}[t(1-t)]^{n} d_{t}\left[-N\left(A_{c}, t\right)\right] \\
\operatorname{tr} A_{c}\left[A_{c}\left(I-A_{c}\right)\right]^{n} & =\int_{0+}^{1} t[t(1-t)]^{n} d_{t}\left[-N\left(A_{c}, t\right)\right]
\end{aligned}
$$

and so by (15), (19), (18), and (20)

$$
\begin{aligned}
& \int_{0}^{1}[t(1-t)]^{n} d_{t}\left[-N\left(A_{c}, t\right)\right]=\frac{\log c}{\pi^{2}} \int_{0}^{1}[t(1-t)]^{n} \frac{d t}{t(1-t)}+o(\log c) \\
& \int_{0}^{1} t[t(1-t)]^{n} d_{t}\left[-N\left(A_{c}, t\right)\right]=\frac{\log c}{\pi^{2}} \int_{0}^{1} t[t(1-t)]^{n} \frac{d t}{t(1-t)}+o(\log c)
\end{aligned}
$$

In consequence, for each fixed polynomial $P$ which vanishes at 0 and at 1 ,

$$
\int_{0}^{1} P(t) d_{t}\left[-N\left(A_{c}, t\right)\right]=\frac{\log c}{\pi^{2}} \int_{0}^{1} P(t) \frac{d t}{t(1-t)}+o(\log c) .
$$

We can eliminate the second restriction by writing $P(t)=t P(1)+$ $[P(t)-t P(1)]$ and applying (3) to the first component. We obtain, for every $P$ vanishing at 0 ,

$$
\begin{aligned}
& \int_{0}^{1} P(t) d_{t}\left[-N\left(A_{c}, t\right)\right] \\
& \quad=P(1) \frac{c}{2 \pi}+\frac{\log c}{\pi^{2}} \int_{0}^{1}[P(t)-t P(1)] \frac{d t}{t(1-t)}+o(\log c)
\end{aligned}
$$

and this relation can be extended by approximation from polynomials to any function $F(t)$ for which $[F(t)-t F(1)] / t(1-t)$ is Riemann-integrable on $[0,1]$. In particular, taking $F(t)$ to be the characteristic function of the interval $0<\alpha<t<1$, we find

$$
\begin{aligned}
N\left(A_{c}, \alpha\right) & =\int_{\alpha}^{1} d_{t}\left[-N\left(A_{c}, t\right)\right] \\
& =\frac{c}{2 \pi}+\frac{\log c}{\pi^{2}}\left[-\int_{0}^{\alpha} \frac{d t}{1-t}+\int_{\alpha}^{1} \frac{d t}{t}\right]+o(\log c) \\
& =\frac{c}{2 \pi}+\frac{\log c}{\pi^{2}} \log \frac{1-\alpha}{\alpha}+o(\log c)
\end{aligned}
$$

This completes the proof of Theorem 1.
The same argument is sufficient to prove the more general case.
Theorem 2. Let $S$ and $T$ each be the union of a finite number, $\mu$ and $v$, respectively, of fixed disjoint closed intervals. Let $A_{c} \equiv P(c S) Q(T) P(c S)$ denotre the operator of (1). Then, for $0<\alpha<1$,

$$
N\left(A_{c}, \alpha\right)=\frac{\mathbf{m}(S) \mathbf{m}(T)}{2 \pi} c+\left(\frac{\mu v}{\pi^{2}} \log \frac{1-a}{\alpha}\right) \log c+o(\log c) .
$$

Proof. We will show how to reduce the proof to that of Theorem 1. Let $\gamma_{i}: g_{i} \leqslant y \leqslant G_{i}$, with $g_{1}<G_{1}<\cdots<g_{\mu}<G_{\mu}$, denote the constituent intervals of $S$. The complement $S^{\prime}$ of $S$ now consists of $\varepsilon_{0}:-\infty<y<g_{1}$, $\varepsilon_{i}: G_{i}<y<g_{i+1}, i=1, \ldots, \mu-1$, and $\varepsilon_{\mu}: G_{\mu}<y<\infty$. Let $\delta_{i}: d_{i} \leqslant t \leqslant D_{i}$, with $d_{1}<D_{1} \cdots<d_{v}<D_{v}$, denote the constitutent intervals of $T$. As before, we consider

$$
\begin{align*}
A_{c}\left(I-A_{c}\right) & =P(c S) Q(T) P\left(c S^{\prime}\right) Q(T) P(c S) \\
& =\sum_{n, m, l, k, j} P\left(c \gamma_{n}\right) Q\left(\delta_{m}\right) P\left(c \varepsilon_{l}\right) Q\left(\delta_{k}\right) P\left(c \gamma_{j}\right) \tag{21}
\end{align*}
$$

Now by (7),

$$
P\left(c \gamma_{n}\right) Q\left(\delta_{m}\right) P\left(c \varepsilon_{l}\right) Q\left(\delta_{k}\right) P\left(c \gamma_{j}\right) \sim P\left(\gamma_{n}\right) Q\left(c \delta_{m}\right) P\left(\varepsilon_{l}\right) Q\left(c \delta_{k}\right) P\left(\gamma_{j}\right)
$$

and, by L 1 and L 2 , this is $O(1)$ unless $n=j$ and $\gamma_{j}$ is adjacent to $\varepsilon_{l}$. If $\varepsilon_{l}$ is a finite interval with $\gamma_{j}$ adjoining it one one side, and we extend $\varepsilon_{l}$ to $\infty$ on the other side, the difference, for example, $P\left(g_{2}, G_{2}\right) Q\left(c \delta_{m}\right) P\left(-\infty, G_{1}\right) Q\left(c \delta_{k}\right)$ $P\left(g_{2}, G_{2}\right)$, is $O(1)$ by L 2 . We conclude that we can reduce the sum in (21) to

$$
\begin{align*}
& A_{c}\left(I-A_{c}\right) \\
& \quad=\sum_{j, m, k} P\left(c \gamma_{j}\right) Q\left(\delta_{m}\right)\left[P\left(-\infty, c g_{j}\right)+P\left(c G_{j}, \infty\right)\right] Q\left(\delta_{k}\right) P\left(c \gamma_{j}\right)+O(1) \tag{22}
\end{align*}
$$

Moreover, by L4, as in Theorem 1

$$
\begin{aligned}
& P\left(c \gamma_{j}\right) Q\left(\delta_{m}\right) P\left(-\infty, c g_{j}\right) Q\left(\delta_{k}\right) P\left(c \gamma_{j}\right) \\
& \quad=P\left(c g_{j}+1, c G_{j}\right) Q\left(\delta_{m}\right) P\left(-\infty, c g_{j}\right) Q\left(\delta_{k}\right) P\left(c g_{j}+1, c G_{j}\right)+O(1)
\end{aligned}
$$

and by L3 the latter operator is $O(1)$ unless $m=k$. The same applies to the remaining operators of (22) and so we find

$$
\begin{aligned}
A_{c}(I- & \left.A_{c}\right) \\
= & \sum_{j, m} P\left(c g_{j}+1, c G_{j}\right) Q\left(\delta_{m}\right) P\left(-\infty, c g_{j}\right) Q\left(\delta_{m}\right) P\left(c g_{j}+1, c G_{j}\right) \\
& +P\left(c g_{j}, c G_{j}-1\right) Q\left(\delta_{m}\right) P\left(c G_{j}, \infty\right) Q\left(\delta_{m}\right) P\left(c g_{j}, c G_{j}-1\right)+O(1) .
\end{aligned}
$$

By taking powers on both sides and applying similar considerations to show that the cross-product terms are $O(1)$ we see that

$$
\begin{align*}
& {\left[A_{c}\left(I-A_{c}\right)\right]^{n}} \\
& \quad=\sum_{j, m}\left[P\left(c g_{j}+1, c G_{j}\right) Q\left(\delta_{m}\right) P\left(-\infty, c g_{j}\right) Q\left(\delta_{m}\right) P\left(c g_{j}+1, c G_{j}\right)\right]^{n} \\
& \quad+\left[P\left(c g_{j}, c G_{j}-1\right) Q\left(\delta_{m}\right) P\left(c G_{j}, \infty\right) Q\left(\delta_{m}\right) P\left(c g_{j}, c G_{j}-1\right)\right]^{n}+O(1) \\
& =\sum_{j, m}\left[P\left(c g_{j}+1, c G_{j}\right) Q\left(d_{m}, \infty\right) P\left(-\infty, c g_{j}\right) Q\left(d_{m}, \infty\right) P\left(c g_{j}+1, c G_{j}\right)\right]^{n} \\
& \quad+\left[P\left(c g_{j}+1, c G_{j}\right) Q\left(D_{m}, \infty\right) P\left(-\infty, c g_{j}\right) Q\left(D_{m}, \infty\right) P\left(c g_{j}+1, c G_{j}\right)\right]^{n} \\
& \quad+\left[P\left(c g_{j}, c G_{j}-1\right) Q\left(d_{m}, \infty\right) P\left(c G_{j}, \infty\right) Q\left(d_{m}, \infty\right) P\left(c g_{j}, c G_{j}-1\right)\right]^{n} \\
& \quad+\left[P\left(c g_{j}, c G_{j}-1\right) Q\left(D_{m}, \infty\right) P\left(c G_{j}, \infty\right) Q\left(D_{m}, \infty\right) P\left(c g_{j}, c G_{j}-1\right)\right]^{n} \\
& \quad+O(1) . \tag{23}
\end{align*}
$$

As in Theorem 1, for given $j$ and $m$, each of these four operators is now unitarily equivalent to

$$
K_{c \mathbf{m}\left(\gamma_{j}\right)}=P\left(1, \mathrm{~cm}\left(\gamma_{j}\right)\right) Q(0, \infty) P(-\infty, 0) Q(0, \infty) P\left(1, c \mathrm{~m}\left(\gamma_{j}\right)\right),
$$

so that $\operatorname{tr}\left[A_{c}\left(I-A_{c}\right)\right]^{n}=\sum_{j, m} 4 \operatorname{tr} K_{c \mathrm{~m}\left(\gamma_{j}\right)}^{n}+O(1)$. But from (19), $\operatorname{tr} K_{c \mathrm{~m}\left(\gamma_{j}\right)}^{n}=$ $\operatorname{tr} K_{c}^{n}+O(1)$, consequently

$$
\begin{equation*}
\operatorname{tr}\left[A_{c}\left(I-A_{c}\right)\right]^{n}=4 \mu \nu \operatorname{tr} K_{c}^{n}+O(1) . \tag{24}
\end{equation*}
$$

Now let us apply $A_{c}=\sum_{a, b, d} P\left(c \gamma_{a}\right) Q\left(\delta_{b}\right) P\left(c \gamma_{d}\right)$ to (23). For the first component of (23), which we denote by $B(j, m)$, we have

$$
A_{c} B(j, m)=\sum_{a, b, d} P\left(c \gamma_{a}\right) Q\left(\delta_{b}\right) P\left(c \gamma_{d}\right) B(j, m)
$$

We now argue that the only significant contribution from this sum comes when $a=d=j$ and $b=m$. For if $d \neq j$, the corresponding operator vanishes. If $a \neq j$, the operator contains the factor

$$
\begin{aligned}
& P\left(c \gamma_{a}\right) Q\left(\delta_{b}\right) P\left(c g_{j}+1, c G_{j}\right) Q\left(d_{m}, \infty\right) P\left(-\infty, c g_{j}\right) \\
& \quad \sim P\left(\gamma_{a}\right) Q\left(c \delta_{b}\right) P\left(g_{j}+c^{-1}, G_{j}\right) Q\left(c d_{m}, \infty\right) P\left(-\infty, g_{j}\right)
\end{aligned}
$$

and the latter operator is $O(1)$ by L1. Finally, $P\left(c \gamma_{j}\right) Q\left(\delta_{b}\right) P\left(c \gamma_{j}\right) B(j, m)$ contains the factor
$Q\left(\delta_{b}\right) P\left(c g_{j}+1, c G_{j}\right) Q\left(d_{m}, \infty\right) P\left(-\infty, c g_{j}\right) Q\left(d_{m}, \infty\right) P\left(c g_{j}+1, c G_{j}\right)$.

If $b<m, \delta_{b}$ is in the complement of $\left[d_{m}, \infty\right)$. Then by (11), the operator of (25) coincides with

$$
-Q\left(\delta_{b}\right) P\left(c g_{j}+1, c G_{j}\right) Q\left(d_{m}, \infty\right) P\left(-\infty, c g_{j}\right) Q\left(-\infty, d_{m}\right) P\left(c g_{j}+1, c G_{j}\right)
$$

and, by (10), this contains a factor unitarily equivalent to $-P\left(\delta_{b}\right)$ $Q\left(c g_{j}+1, c G_{j}\right) P\left(d_{m}, \infty\right) Q\left(-\infty, c g_{j}\right) P\left(-\infty, d_{m}\right)$ which is $O(1)$ by L3. Similarly, if $b>m$, (25) coincides with $-Q\left(\delta_{b}\right) P\left(c g_{j}+1, c G_{j}\right) Q\left(-\infty, d_{m}\right)$ $P\left(-\infty, c g_{j}\right) Q\left(d_{m}, \infty\right) P\left(c g_{j}+1, c G_{j}\right)$, and again contains a factor equivalent to $-P\left(\delta_{b}\right) Q\left(c g_{j}+1, c G_{j}\right) P\left(-\infty, d_{m}\right) Q\left(-\infty, c g_{j}\right) P\left(d_{m}, \infty\right)$; as $\delta_{b}$ is now in the complement of $\left(-\infty, d_{m}\right]$, this is $O(1)$ by L 3 . We conclude that

$$
A_{c} B(j, m)=P\left(c \gamma_{j}\right) Q\left(\delta_{m}\right) B(j, m)+O(1)
$$

and similarly for each of the remaining terms of (23). Continuing as in Theorem 1, we find

$$
\begin{equation*}
\operatorname{tr} A_{c}\left[A_{c}\left(I-A_{c}\right)\right]^{n}=\frac{1}{2} \operatorname{tr}\left[A_{c}\left(I-A_{c}\right)\right]^{n}+O(1) \tag{26}
\end{equation*}
$$

The remaining argument of Theorem 1 now applies without change to (24) and (26), and yields Theorem 2.

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