# MAXIMUM SPREAD OF GRAPHS AND BIPARTITE GRAPHS 

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#### Abstract

Given any graph $G$, the spread of $G$ is the maximum difference between any two eigenvalues of the adjacency matrix of $G$. In this paper, we resolve a pair of 20-year-old conjectures of Gregory, Hershkowitz, and Kirkland regarding the spread of graphs. The first states that for all positive integers $n$, the $n$-vertex graph $G$ that maximizes spread is the join of a clique and an independent set, with $[2 n / 3\rfloor$ and $[n / 3\rceil$ vertices, respectively. Using techniques from the theory of graph limits and numerical analysis, we prove this claim for all $n$ sufficiently large. As an intermediate step, we prove an analogous result for a family of operators in the Hilbert space over $\mathscr{L}^{2}[0,1]$. The second conjecture claims that for any fixed $m \leq n^{2} / 4$, if $G$ maximizes spread over all $n$-vertex graphs with $m$ edges, then $G$ is bipartite. We prove an asymptotic version of this conjecture. Furthermore, we construct an infinite family of counterexamples, which shows that our asymptotic solution is tight up to lower-order error terms.


## 1. Introduction

The spread $s(M)$ of an arbitrary $n \times n$ complex matrix $M$ is the diameter of its spectrum; that is,

$$
s(M):=\max _{i, j}\left|\lambda_{i}-\lambda_{j}\right|,
$$

where the maximum is taken over all pairs of eigenvalues of $M$. This quantity has been well studied in general (see [11, 16, 22, 33] for details and additional references). Most notably, Johnson, Kumar, and Wolkowitz [16, Theorem 2.1] obtained the lower bound

$$
s(M) \geq\left|\sum_{i \neq j} m_{i, j}\right| /(n-1)
$$

for normal matrices $M=\left(m_{i, j}\right)$, and Mirsky [22, Theorem 2] determined the upper bound

$$
s(M) \leq \sqrt{2 \sum_{i, j}\left|m_{i, j}\right|^{2}-(2 / n)\left|\sum_{i} m_{i, i}\right|^{2}}
$$

for any $n \times n$ matrix $M$, which is tight for any normal matrix with $n-2$ of its eigenvalues all equal and equal to the arithmetic mean of the other two.

The spread of a matrix has also received interest in some particular cases. Consider a simple undirected graph $G=(V(G), E(G))$ of order $n$. The adjacency matrix $A$ of $G$ is the $n \times n$ matrix whose rows and columns are indexed by the vertices of $G$, with entries satisfying $A_{u, v}=1$ if $\{u, v\} \in E(G)$ and $A_{u, v}=0$ otherwise. This matrix is real and symmetric, and so its eigenvalues are real and can be ordered $\lambda_{1}(G) \geq \lambda_{2}(G) \geq \cdots \geq$

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$\lambda_{n}(G)$. When considering the spread of the adjacency matrix $A$ of some graph $G$, the spread is simply the distance between $\lambda_{1}(G)$ and $\lambda_{n}(G)$, denoted by

$$
s(G):=\lambda_{1}(G)-\lambda_{n}(G) .
$$

In this instance, $s(G)$ is referred to as the spread of the graph $G$.
In [13], Gregory, Hershkowitz, and Kirkland investigated a number of properties regarding the spread of a graph, determined upper and lower bounds on $s(G)$, and made two key conjectures. Let us denote the maximum spread over all $n$-vertex graphs by $s(n)$, the maximum spread over all $n$-vertex graphs of size $m$ by $s(n, m)$, and the maximum spread over all $n$-vertex bipartite graphs of size $m$ by $s_{b}(n, m)$. The join of two vertex-disjoint graphs $G \vee H$ is the graph containing both $G$ and $H$, and all edges between the two. The complement of a graph $\bar{G}$ is the graph on $V(G)$ containing only the edges not in $E(G)$. Let $K_{k}$ be the clique of order $k$ and $G(n, k):=K_{k} \vee \overline{K_{n-k}}$ be the join of the clique $K_{k}$ (the graph with all possible edges included) and the independent set $\overline{K_{n-k}}$ (the graph with no edges). We say a graph is spread extremal if it has spread $s(n)$. The conjectures addressed in this article are as follows.

Conjecture 1 ([13, Conjecture 1.3]). For any positive integer $n$, the graph of order $n$ with maximum spread is $G(n,\lfloor 2 n / 3\rfloor)$; that is, $s(n)=\left\lfloor(4 / 3)\left(n^{2}-n+1\right)\right\rfloor^{1 / 2}$ and this value is attained only by $G(n,[2 n / 3\rfloor)$.
Conjecture 2 ([13, Conjecture 1.4]). If G is a graph with nvertices and $m \leq\left\lfloor n^{2} / 4\right\rfloor$ edges attaining the maximum spread $s(n, m)$, then $G$ must be bipartite. That is, $s_{b}(n, m)=$ $s(n, m)$ for all $m \leq\left\lfloor n^{2} / 4\right\rfloor$.

Conjecture 11 is referred to as the Spread Conjecture, and Conjecture 2 is referred to as the Bipartite Spread Conjecture. Much of what is known about Conjecture 11 is contained in [13], but the reader may also see [29] for a description of the problem and references to other work. In this paper, we resolve both conjectures. We prove the Spread Conjecture for all $n$ sufficiently large, prove an asymptotic version of the Bipartite Spread Conjecture, and provide an infinite family of counterexamples to illustrate that our asymptotic version of the Bipartite Spread Conjecture is tight up to lower-order error terms. These results are given by Theorems 1.1, 1.2, and 1.3.
Theorem 1.1. There exists a constant $N$ so that the following holds: Suppose $G$ is a graph on $n \geq N$ vertices with maximum spread. Then $G$ is the join of a clique on $\lfloor 2 n / 3\rfloor$ vertices and an independent set on $\lceil n / 3\rceil$ vertices.

Theorem 1.2. For all $n, m \in \mathbb{N}$ satisfying $m \leq\left\lfloor n^{2} / 4\right\rfloor$,

$$
s(n, m)-s_{b}(n, m) \leq \frac{1+16 m^{-3 / 4}}{m^{3 / 4}} s(n, m) .
$$

Theorem 1.3. For any $\varepsilon>0$, there exists some $n_{\varepsilon}$ such that

$$
s(n, m)-s_{b}(n, m) \geq \frac{1-\varepsilon}{m^{3 / 4}} s(n, m)
$$

for all $n \geq n_{\varepsilon}$ and some $m \leq\left\lfloor n^{2} / 4\right\rfloor$ depending on $n$.
The proof of Theorem 1.1 constitutes the main subject of this work. The general technique consists of showing that a spread-extremal graph has certain desirable properties, solving an analogous problem for graph limits, and then using this result to
say something about the Spread Conjecture for sufficiently large $n$. For the interested reader, we state the analogous graph limit result in the language of functional analysis (see [18] for an introduction to graph limits).

Theorem 1.4. Let $W:[0,1]^{2} \rightarrow[0,1]$ be a Lebesgue-measurable function such that $W(x, y)=W(y, x)$ for a.e. $(x, y) \in[0,1]^{2}$, and let $A=A_{W}$ be the kernel operator on $\mathscr{L}^{2}[0,1]$ associated with $W$. For all unit functions $f, g \in \mathscr{L}^{2}[0,1]$,

$$
\langle f, A f\rangle-\langle g, A g\rangle \leq \frac{2}{\sqrt{3}}
$$

Moreover, equality holds if and only if there exists a measure-preserving transformation $\sigma$ on $[0,1]$ such that for a.e. $(x, y) \in[0,1]^{2}$,

$$
W(\sigma(x), \sigma(y))= \begin{cases}0, & \text { if }(x, y) \in[2 / 3,1] \times[2 / 3,1] \\ 1, & \text { otherwise }\end{cases}
$$

The bound of $2 / \sqrt{3}$ in Theorem 1.4 is the scaled limit of the spread $s(G(n,\lfloor 2 n / 3\rfloor))=$ $\left[(4 / 3)\left(n^{2}-n+1\right)\right\rfloor^{1 / 2}$, and the function $W$ (up to a set of measure zero) corresponds to a scaled limit of $G(n,\lfloor 2 n / 3\rfloor)$. In addition to being a key ingredient in the proof of Theorem 1.1, Theorem 1.4 also immediately implies a result for arbitrary symmetric nonnegative matrices.

Corollary 1.5. Let $A=\left(a_{i, j}\right)$ be an $n \times n$ symmetric nonnegative matrix. Then

$$
\lambda_{1}(A)-\lambda_{n}(A) \leq \frac{2 n}{\sqrt{3}} \max _{i, j} a_{i, j}
$$

and

$$
\max _{\substack{\|u\|=\|v\|=1 \\\langle u, v\rangle=0}}|\langle u, A v\rangle| \leq \frac{n}{\sqrt{3}} \max _{i, j} a_{i, j} .
$$

Corollary 1.5, paired with the lower bound provided by graph $G(n,\lfloor 2 n / 3\rfloor)$, implies that $(2 n-1) / \sqrt{3}<s(n) \leq 2 n / \sqrt{3}$ for all $n>1$, i.e. the spread conjecture is true for all $n$ up to an additive $1 / \sqrt{3}$ factor. Furthermore, Corollary 1.5 implies that the maximum spread of a symmetric $0-1$ matrix is exactly $2 n / \sqrt{3}$ for $n \equiv 0 \bmod 3$ and is within $1 /(2 \sqrt{3} n)$ of $2 n / \sqrt{3}$ for $n \not \equiv 0 \bmod 3$. The second part of Corollary 1.5 gives a bound on the magnitude of off-diagonal entries of a nonnegative matrix under a unitary change of basis, and is also tight for $n \equiv 0 \bmod 3($ and tight up to $O(1 / n)$ for $n \not \equiv 0 \bmod 3)$.

The proof of Theorem 1.1 can be found in Sections 2 through 6, with certain technical details reserved for Appendices A and B. We provide an in-depth overview of the proof of Theorem 1.1 in Subsection 1.1. By comparison, the proofs of Theorems 1.2 and 1.3 are surprisingly short, making use of the theory of equitable decompositions and a well-chosen class of counterexamples. Their proof can be found in Section 7. Finally, in Section 8, we discuss further questions and possible future avenues of research.
1.1. High-level outline of spread proof. Here, we provide a concise, high-level description of our asymptotic proof of the Spread Conjecture. The proof itself is quite involved, making use of interval arithmetic and a number of fairly complicated symbolic calculations, but conceptually, it is quite intuitive. It consists of four main steps.

Step 1 (Graph-theoretic results). In Section 2 , we observe a number of important structural properties of any graph that maximizes spread for a given order $n$. In particular, we show that

- any graph that maximizes spread must be the join of two threshold graphs (Lemma 2.1),
- both graphs in this join have order linear in $n$ (Lemma 2.2),
- any unit eigenvectors $\mathbf{x}$ and $\mathbf{z}$ corresponding to $\lambda_{1}(A)$ and $\lambda_{n}(A)$ have infinity norms of order $n^{-1 / 2}$ (Lemma 2.3),
- the quantities $\lambda_{1} \mathbf{x}_{u}^{2}-\lambda_{n} \mathbf{z}_{u}^{2}, u \in V$, are all nearly equal, up to a term of order $n^{-1}$ (Lemma 2.4).
This last structural property serves as the backbone of our proof. In addition, we note that, by a tensor argument, an asymptotic upper bound for $s(n)$ implies a bound for all $n$.

Step 2 (Graphons and a finite-dimensional eigenvalue problem). In Sections 3 and 4, we make use of graphons to understand how spread-extremal graphs behave as $n$ tends to infinity. Section 3 consists of a basic introduction to graphons, and a translation of the graph results of Step it the graphon setting. In particular, we prove the graphon analogue of the spread-extremal graph properties that

- vertices $u$ and $v$ are adjacent if and only if $\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}>0$ (Lemma 3.6),
- the quantities $\lambda_{1} \mathbf{x}_{u}^{2}-\lambda_{n} \mathbf{z}_{u}^{2}, u \in V$, are all nearly equal (Lemma 3.7).

Next, in Section 4, we show that the spread-extremal graphon for our problem takes the form of a particular stepgraphon with a finite number of blocks (Theorem 4.1). Through an averaging argument, we note that the spread-extremal graphon takes the form of a stepgraphon with a fixed structure of symmetric $7 \times 7$ blocks, illustrated below (black equals one, white equals zero).


The lengths $\alpha=\left(\alpha_{1}, \ldots, \alpha_{7}\right), \alpha^{T} \mathbf{1}=1$, of each row and column in the spread-extremal stepgraphon are unknown. For any choice of lengths $\alpha$, we can associate a $7 \times 7$ matrix whose spread is identical to that of the associated stepgraphon pictured above. Let $B$ be the $7 \times 7$ matrix with $B_{i, j}$ equal to the value of the above stepgraphon on block $i, j$, and $D=\operatorname{diag}\left(\alpha_{1}, \ldots, \alpha_{7}\right)$ be a diagonal matrix with $\alpha$ on the diagonal. Then the matrix $D^{1 / 2} B D^{1 / 2}$ has spread equal to the spread of the associated stepgraphon.


FIGURE 1. Contour plots of the spread for some choices of $\alpha$. Each point ( $x, y$ ) of Plot (a) illustrates the maximum spread over all choices of $\alpha$ satisfying $\alpha_{3}+\alpha_{4}=x$ and $\alpha_{6}+\alpha_{7}=y$ (and therefore, $\alpha_{1}+\alpha_{2}+$ $\alpha_{5}=1-x-y$ ) on a grid of step size $1 / 100$. Each point $(x, y)$ of Plot (b) illustrates the maximum spread over all choices of $\alpha$ satisfying $\alpha_{2}=\alpha_{3}=\alpha_{4}=0, \alpha_{5}=y$, and $\alpha_{7}=x$ on a grid of step size $1 / 100$. The maximum spread of Plot (a) is achieved at the two black "x"s, and implies that, without loss of generality, $\alpha_{3}+\alpha_{4}=0$, and therefore $\alpha_{2}=0$ (indices $\alpha_{1}$ and $\alpha_{2}$ can be combined when $\alpha_{3}+\alpha_{4}=0$ ). Plot (b) treats this case when $\alpha_{2}=\alpha_{3}=\alpha_{4}=0$, and the maximum spread is achieved on the black line. This implies that either $\alpha_{5}=0$ or $\alpha_{7}=0$. In both cases, this reduces to the block $2 \times 2$ case $\alpha_{1}, \alpha_{7} \neq 0$ (or, if $\alpha_{7}=0$, then $\alpha_{1}, \alpha_{6} \neq 0$ ).

Step 3 (Computer-assisted proof of a finite-dimensional eigenvalue problem). In Section 5 , we show that the optimizing choice of $\alpha$ is, without loss of generality, given by $\alpha_{1}=2 / 3, \alpha_{7}=1 / 3$, and all other $\alpha_{i}=0$ (Theorem 5.1). This is exactly the limit of the conjectured spread-extremal graph as $n$ tends to infinity. The proof of this fact is extremely technical, and relies on a computer-assisted proof using both interval arithmetic and symbolic computations. This is the only portion of the proof (of Theorem 1.1) that requires the use of interval arithmetic. Though not a proof, in Figure 1 iwe provide intuitive visual justification of this result. In this figure, we provide contour plots resulting from numerical computations of the spread of the above matrix for various values of $\alpha$. The numerical results suggest that the $2 \times 2$ block stepgraphon with lengths $2 / 3$ and $1 / 3$ is indeed optimal. See Figure 1 and the associated caption for details.

The actual proof of this fact consists of the following steps:

- we reduce the possible choices of nonzero $\alpha_{i}$ from $2^{7}$ to 17 different cases (Lemma A.2),
- using eigenvalue equations, the graphon versions of $\lambda_{1} \mathbf{x}_{u}^{2}-\lambda_{n} \mathbf{z}_{u}^{2}$ all nearly equal, and interval arithmetic, we prove that, of the 17 cases, only the cases
- $\alpha_{1}, \alpha_{7} \neq 0$

$$
-\alpha_{4}, \alpha_{5}, \alpha_{7} \neq 0
$$

can produce a spread-extremal stepgraphon (Lemma 5.2),

- prove that the $3 \times 3$ case cannot be spread extremal, using basic results from the theory of cubic polynomials and computer-assisted symbolic calculations (Lemma 5.4).
This proves that the spread-extremal graphon is a $2 \times 2$ stepgraphon that, without loss of generality, takes value zero on the block $[2 / 3,1]^{2}$ and one elsewhere (Theorem 1.4/Theorem 5.1).

Step 4 (From graphons to an asymptotic proof of the spread conjecture). Finally, in Section 6, we convert our result for the spread-extremal graphon to a statement for graphs. This process consists of two main parts:

- using our graphon theorem (Theorem 5.1), we show that any spread-extremal graph takes the form $\left(K_{n_{1}} \dot{\cup} \overline{K_{n_{2}}}\right) \vee \overline{K_{n_{3}}}$ for $n_{1}=(2 / 3+o(1)) n, n_{2}=o(n)$, and $n_{3}=(1 / 3+o(1)) n$ (Lemma 6.2), i.e. any spread-extremal graph is equal up to a set of $o(n)$ vertices to the conjectured optimal graph $K_{[2 n / 3]} \vee \overline{K_{[n / 3]}}$,
- we show that, for $n$ sufficiently large, the spread of $\left(K_{n_{1}} \dot{\cup} \overline{K_{n_{2}}}\right) \vee \overline{K_{n_{3}}}, n_{1}+n_{2}+$ $n_{3}=n$, is maximized when $n_{2}=0$ (Lemma 6.3).
Together, these two results complete our proof of the spread conjecture for sufficiently large $n$ (Theorem 1.1).


## 2. Properties of spread-extremal graphs

In this section, we review what has already been shown about spread-extremal graphs ( $n$-vertex graphs with spread $s(n)$ ) in [13]. We then prove a number of properties of spread-extremal graphs and properties of the eigenvectors associated with the maximum and minimum eigenvalues of a spread-extremal graph.

Let $G$ be a graph, and let $A$ be the adjacency matrix of $G$, with eigenvalues $\lambda_{1} \geq \cdots \geq$ $\lambda_{n}$. For unit vectors $\mathbf{x}, \mathbf{z} \in \mathbb{R}^{n}$, we have

$$
\lambda_{1} \geq \mathbf{x}^{T} A \mathbf{x} \quad \text { and } \quad \lambda_{n} \leq \mathbf{z}^{T} A \mathbf{z}
$$

Hence (as observed in [13]), the spread of a graph can be expressed as

$$
\begin{equation*}
s(G)=\max _{\mathbf{x}, \mathbf{z}} \sum_{u \sim v}\left(\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}\right), \tag{1}
\end{equation*}
$$

where the maximum is taken over all unit vectors $\mathbf{x}, \mathbf{z}$. This maximum is attained only for $\mathbf{x}, \mathbf{z}$ orthonormal eigenvectors corresponding to the eigenvalues $\lambda_{1}, \lambda_{n}$, respectively. We refer to such a pair of vectors $\mathbf{x}, \mathbf{z}$ as extremal eigenvectors of $G$. For any two vectors $\mathbf{x}, \mathbf{z}$ in $\mathbb{R}^{n}$, let $G(\mathbf{x}, \mathbf{z})$ denote the graph for which distinct vertices $u, v$ are adjacent if and only if $\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v} \geq 0$. As noted in [13, Proposition 3.3], any spread-extremal graph must be connected (as the join $G_{1} \vee G_{2}$ of nontrivial vertex-disjoint graphs is strictly greater than their union $\left.G_{1} \cup G_{2}\right)$. Then from the above, there is some graph $G(\mathbf{x}, \mathbf{z})$ that is a spread-extremal graph, with $\mathbf{x}, \mathbf{z}$ orthonormal and $\mathbf{x}$ positive [13, Lemma 3.5]. By the Perron-Frobenius Theorem, $\lambda_{1}$ is a simple eigenvalue. However, $\lambda_{n}$ may be repeated, and so $\mathbf{z}$ may not be unique up to normalization.

In addition, we enhance [13, Lemmas 3.4 and 3.5] using some helpful definitions and the language of threshold graphs. Whenever $G=G(\mathbf{x}, \mathbf{z})$ is understood, let $P=$ $P(\mathbf{x}, \mathbf{z}):=\left\{u \in V(G): \mathbf{z}_{u} \geq 0\right\}$ and $N=N(\mathbf{x}, \mathbf{z}):=V(G) \backslash P$. In addition, let $G[S]$,
$S \subset V$, be the subgraph induced by $S$, i.e., the graph with vertex set $S$ and containing all edges of $G$ between vertices in $S$.

For our purposes, we say that $G$ is a threshold graph if and only if there exists a function $\varphi: V(G) \rightarrow(-\infty, \infty]$ such that for all distinct $u, v \in V(G), u v \in E(G)$ if and only if $\varphi(u)+\varphi(v) \geq 0$. Here, $\varphi$ is a threshold function for $G$ (with 0 as its threshold). The following detailed lemma shows that any spread-extremal graph is the join of two threshold graphs with threshold functions that can be made explicit.

Lemma 2.1. Let $n>2$, and suppose $G$ is an $n$-vertex graph such that $s(G)=s(n)$. Denote by $\mathbf{x}$ and $\mathbf{z}$ extremal unit eigenvectors for $G$. Then
(i) For any two vertices $u, v$ of $G$, $u$ and $v$ are adjacent whenever $\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}>0$ and $u$ and $v$ are nonadjacent whenever $\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}<0$.
(ii) For any distinct $u, v \in V(G), \mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v} \neq 0$.
(iii) Let $P:=P(\mathbf{x}, \mathbf{z}), N:=N(\mathbf{x}, \mathbf{z})$ and let $G_{1}:=G[P]$ and $G_{2}:=G[N]$. Then $G=$ $G(\mathbf{x}, \mathbf{z})=G_{1} \vee G_{2}$.
(iv) For each $i \in\{1,2\}, G_{i}$ is a threshold graph with threshold function defined on all $u \in V\left(G_{i}\right)$ by

$$
\varphi(u):=\log \left|\frac{\mathbf{x}_{u}}{\mathbf{z}_{u}}\right| .
$$

Proof. Suppose $G$ is an $n$-vertex graph such that $s(G)=s(n)$, and write $A=\left(a_{u v}\right)_{u, v \in V(G)}$ for its adjacency matrix. Item (iil) is equivalent to Lemma 3.4 from [13]. For completeness, we include a proof. By Equation (1]) we have that

$$
s(G)=\max _{\mathbf{x}, \mathbf{z}} \mathbf{x}^{T} A \mathbf{x}-\mathbf{z}^{T} A \mathbf{z}=\sum_{u, v \in V(G)} a_{u v} \cdot\left(\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}\right),
$$

where the maximum is taken over all unit vectors of length $n$. If $\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}>0$ and $a_{u v}=0$, then $s(G+u v)>s(G)$, a contradiction. And if $\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}<0$ and $a_{u v}=1$, then $s(G-u v)>s(G)$, a contradiction. So Item (ii) holds.

For a proof of Item (iii), suppose $\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}=0$ and denote by $G^{\prime}$ the graph formed by adding or deleting the edge $u v$ from $G$. With $A^{\prime}$ denoting the adjacency matrix of $G^{\prime}$,

$$
s\left(G^{\prime}\right) \geq \mathbf{x}^{T} A^{\prime} \mathbf{x}-\mathbf{z}^{T} A^{\prime} \mathbf{z}=\mathbf{x}^{T} A \mathbf{x}-\mathbf{z}^{T} A \mathbf{z}=s(G) \geq s\left(G^{\prime}\right)
$$

so each inequality is an equality. It follows that $\mathbf{x}, \mathbf{z}$ are eigenvectors for $A^{\prime}$. Furthermore, without loss of generality, we may assume that $u v \in E(G)$. In particular, there exists some $\lambda^{\prime}$ such that

$$
\begin{aligned}
A \mathbf{x} & =\lambda \mathbf{x}, \\
\left(A-\mathbf{e}_{u} \mathbf{e}_{v}^{T}-\mathbf{e}_{v} \mathbf{e}_{u}^{T}\right) \mathbf{x} & =\lambda^{\prime} \mathbf{x} .
\end{aligned}
$$

So $\left(\mathbf{e}_{u} \mathbf{e}_{v}^{T}+\mathbf{e}_{v} \mathbf{e}_{u}^{T}\right) \mathbf{x}=\left(\lambda-\lambda^{\prime}\right) \mathbf{x}$. Let $w \in V(G) \backslash\{u, v\}$. By the above equation, $(\lambda-$ $\left.\lambda^{\prime}\right) \mathbf{x}_{w}=0$ and either $\lambda^{\prime}=\lambda$ or $\mathbf{x}_{w}=0$. To find a contradiction, it is sufficient to note that $G$ is a connected graph with Perron-Frobenius eigenvector $\mathbf{x}$. Indeed, let $P:=\left\{w \in V(G): \mathbf{z}_{w} \geq 0\right\}$ and let $N:=V(G) \backslash P$. Then for any $w \in P$ and any $w^{\prime} \in N$, $\mathbf{x}_{w} \mathbf{x}_{w^{\prime}}-\mathbf{z}_{w} \mathbf{z}_{w^{\prime}}>0$, and by Item (i), $w w^{\prime} \in E(G)$. So $G$ is connected. This completes the proof of Item (iii).

[^0]Now, we prove Item (iiii). To see that $G=G(\mathbf{x}, \mathbf{z})$, note by Items (ii) and (iii), for all distinct $u, v \in V(G), \mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}>0$ if and only if $u v \in E(G)$, and otherwise, $\mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v}<0$ and $u v \notin E(G)$. To see that $G=G_{1} \vee G_{2}$, note that for any $u \in P$ and any $v \in N, 0 \neq \mathbf{x}_{u} \mathbf{x}_{v}-\mathbf{z}_{u} \mathbf{z}_{v} \geq \mathbf{z}_{u} \cdot\left(-\mathbf{z}_{v}\right) \geq 0$.

Finally, we prove Item (iv). Suppose $u, v$ are distinct vertices such that either $u, v \in$ $P$ or $u, v \in N$. Allowing the possibility that $0 \in\left\{\mathbf{z}_{u}, \mathbf{z}_{v}\right\}$, the following equivalence holds:

$$
\begin{aligned}
\varphi(u)+\varphi(v) & \geq 0 & & \text { if and only if } \\
\log \left|\frac{\mathbf{x}_{u} \mathbf{x}_{v}}{\mathbf{z}_{u} \mathbf{z}_{v}}\right| & \geq 1 & & \text { if and only if } \\
\mathbf{x}_{u} \mathbf{x}_{v}-\left|\mathbf{z}_{u} \mathbf{z}_{v}\right| & \geq 0 . & &
\end{aligned}
$$

Since $\mathbf{z}_{u}, \mathbf{z}_{v}$ have the same sign, Item (iiv) holds. This completes the proof.
From [21], we recall the following useful characterization in terms of "nesting" neighborhoods: $G$ is a threshold graph if and only there exists a numbering $v_{1}, \cdots, v_{n}$ of $V(G)$ such that for all $1 \leq i<j \leq n$, if $v_{k} \in V(G) \backslash\left\{v_{i}, v_{j}\right\}, v_{j} v_{k} \in E(G)$ implies that $v_{i} v_{k} \in E(G)$. Given this ordering, if $k$ is the largest natural number such that $v_{k} v_{k+1} \in E(G)$, then the set $\left\{v_{1}, \cdots, v_{k}\right\}$ induces a clique and the set $\left\{v_{k+1}, \cdots, v_{n}\right\}$ induces an independent set.

Lemma 2.2 shows that both $P$ and $N$ have linear size.
Lemma 2.2. If $G$ is a spread-extremal graph, then both $P$ and $N$ have size $\geq n / 100$.
Proof. We will show that $P$ and $N$ both have size at least $\frac{n}{100}$. First, since $G$ is spreadextremal, it has spread more than $1.1 n$ and hence has smallest eigenvalue $\lambda_{n}<-\frac{n}{10}$. Without loss of generality, for the remainder of this proof we will assume that $|P| \leq|N|$, and that $v$ is a vertex satisfying $\left|\mathbf{z}_{v}\right|=\|\mathbf{z}\|_{\infty}$. By way of contradiction, assume that $|P|<\frac{n}{100}$.

If $v \in N$, then we have

$$
-\lambda_{n}=\lambda_{n} \frac{\mathbf{z}_{v}}{\left|\mathbf{z}_{v}\right|}=\sum_{u \sim v} \frac{\mathbf{z}_{u}}{\left|\mathbf{z}_{v}\right|} \leq \sum_{u \in P} \frac{\mathbf{z}_{u}}{\left|\mathbf{z}_{v}\right|} \leq|P|<\frac{n}{100},
$$

contradicting that $\lambda_{n}<\frac{-n}{10}$. Therefore, assume that $v \in P$. Then

$$
\begin{aligned}
\lambda_{n}^{2} & =\lambda_{n}^{2} \frac{\mathbf{z}_{v}}{\left|\mathbf{z}_{v}\right|}=\sum_{u \sim v} \sum_{w \sim u} \frac{\mathbf{z}_{w}}{\left|\mathbf{z}_{v}\right|} \leq \sum_{u \sim v} \sum_{\substack{w \sim u \\
w \in P}} \frac{\mathbf{z}_{w}}{\left|\mathbf{z}_{v}\right|} \\
& \leq|P||N|+2|E(G[P])| \leq|P||N|+|P|^{2} \leq \frac{99 n^{2}}{100^{2}}+\frac{n^{2}}{100^{2}} .
\end{aligned}
$$

This gives $\left|\lambda_{n}\right| \leq \frac{n}{10}$, a contradiction.
Lemma 2.3. If $\mathbf{x}$ and $\mathbf{z}$ are unit eigenvectors for $\lambda_{1}$ and $\lambda_{n}$, then $\|\mathbf{x}\|_{\infty}=O\left(n^{-1 / 2}\right)$ and $\|\mathbf{z}\|_{\infty}=O\left(n^{-1 / 2}\right)$.
Proof. During this proof we will assume that $\hat{u}$ and $\hat{v}$ are vertices satisfying $\|\mathbf{x}\|_{\infty}=\mathbf{x}_{\hat{u}}$ and $\|\mathbf{z}\|_{\infty}=\left|\mathbf{z}_{\hat{v}}\right|$ and without loss of generality that $\hat{v} \in N$. We will use the weak
estimates that $\lambda_{1}>\frac{n}{2}$ and $\lambda_{n}<-\frac{n}{10}$. Define sets

$$
\begin{aligned}
& A=\left\{w: \mathbf{x}_{w}>\frac{\mathbf{x}_{\hat{u}}}{4}\right\}, \\
& B=\left\{w: \mathbf{z}_{w}>-\frac{\mathbf{z}_{\hat{v}}}{20}\right\} .
\end{aligned}
$$

We have the estimates

$$
1=\mathbf{x}^{T} \mathbf{x} \geq \sum_{w \in A} \mathbf{x}_{w}^{2} \geq|A| \frac{\|\mathbf{x}\|_{\infty}^{2}}{16} \quad \text { and } \quad 1=\mathbf{z}^{T} \mathbf{z} \geq \sum_{w \in B} \mathbf{z}_{w}^{2} \geq|B| \frac{\|\mathbf{z}\|_{\infty}^{2}}{400}
$$

and so $\|\mathbf{x}\|_{\infty} \leq 4 / \sqrt{|A|}$ and $\|\mathbf{z}\|_{\infty} \leq 20 / \sqrt{|B|}$. To complete the proof, it suffices to show that $A$ and $B$ both have size $\Omega(n)$.

We now give a lower bound on the sizes of $A$ and $B$ using the eigenvalue-eigenvector equation and the weak bounds on $\lambda_{1}$ and $\lambda_{n}$.

$$
\frac{n}{2}\|\mathbf{x}\|_{\infty}=\frac{n}{2} \mathbf{x}_{\hat{u}}<\lambda_{1} \mathbf{x}_{\hat{u}}=\sum_{w \sim \hat{u}} \mathbf{x}_{w} \leq\|\mathbf{x}\|_{\infty}\left(|A|+\frac{1}{4}(n-|A|)\right),
$$

giving that $|A|>\frac{n}{3}$. Similarly,

$$
\frac{n}{10}\|\mathbf{z}\|_{\infty}=-\frac{n}{10} \mathbf{z}_{\hat{\vartheta}}<\lambda_{n} \mathbf{z}_{\hat{0}}=\sum_{w \sim \hat{\hat{0}}} \mathbf{z}_{w} \leq\|\mathbf{z}\|_{\infty}\left(|B|+\frac{1}{20}(n-|B|)\right),
$$

and so $|B|>\frac{n}{19}$. This implies that $\|\mathbf{x}\|_{\infty}<\frac{4}{\sqrt{3}} n^{-1 / 2}$ and $\|\mathbf{z}\|_{\infty} \leq \frac{20}{\sqrt{19}} n^{-1 / 2}$.

Lemma 2.4. Assume that $\mathbf{x}$ and $\mathbf{z}$ are unit vectors. Then there exists a constant $C$ such that for any pair of vertices $u$ and $v$, we have

$$
\left|\left(\lambda_{1} \mathbf{x}_{u}^{2}-\lambda_{n} \mathbf{z}_{u}^{2}\right)-\left(\lambda_{1} \mathbf{z}_{v}^{2}-\lambda_{n} \mathbf{z}_{v}^{2}\right)\right|<\frac{C}{n} .
$$

Proof. Let $u$ and $v$ be vertices, and create a graph $\tilde{G}$ by deleting $u$ and cloning $v$. That is, $V(\tilde{G})=\left\{v^{\prime}\right\} \cup V(G) \backslash\{u\}$ and

$$
E(\tilde{G})=E(G \backslash\{u\}) \cup\left\{v^{\prime} w: v w \in E(G)\right\} .
$$

Note that $v \nsim v^{\prime}$. Let $\tilde{A}$ be the adjacency matrix of $\tilde{G}$. Define two vectors $\tilde{\mathbf{x}}$ and $\tilde{\mathbf{z}}$ by

$$
\tilde{\mathbf{x}}_{w}= \begin{cases}\mathbf{x}_{w} & w \neq v^{\prime}, \\ \mathbf{x}_{v} & w=v^{\prime},\end{cases}
$$

and

$$
\tilde{\mathbf{z}}_{w}= \begin{cases}\mathbf{z}_{w} & w \neq v^{\prime}, \\ \mathbf{z}_{v} & w=v .\end{cases}
$$

Then $\tilde{\mathbf{x}}^{T} \tilde{\mathbf{x}}=1-\mathbf{x}_{u}^{2}+\mathbf{x}_{v}^{2}$ and $\tilde{\mathbf{z}}^{T} \tilde{\mathbf{z}}=1-\mathbf{z}_{u}^{2}+\mathbf{z}_{v}^{2}$. Similarly,

$$
\begin{aligned}
\tilde{\mathbf{x}}^{T} \tilde{A} \tilde{\mathbf{x}} & =\lambda_{1}-2 \mathbf{x}_{u} \sum_{u w \in E(G)} \mathbf{x}_{w}+2 \mathbf{x}_{v^{\prime}} \sum_{v w \in E(G)} \mathbf{x}_{w}-2 A_{u v} \mathbf{x}_{v} \mathbf{x}_{u} \\
& =\lambda_{1}-2 \lambda_{1} \mathbf{x}_{u}^{2}+2 \lambda_{1} \mathbf{x}_{v}^{2}-2 A_{u v} \mathbf{x}_{u} \mathbf{x}_{v},
\end{aligned}
$$

and

$$
\begin{aligned}
\tilde{\mathbf{z}}^{T} \tilde{A} \tilde{\mathbf{z}} & =\lambda_{n}-2 \mathbf{z}_{u} \sum_{u w \in E(G)} \mathbf{z}_{w}+2 \mathbf{z}_{v^{\prime}} \sum_{v w \in E(G)} \mathbf{z}_{w}-2 A_{u v} \mathbf{z}_{v} \mathbf{z}_{u} \\
& =\lambda_{n}-2 \lambda_{n} \mathbf{z}_{u}^{2}+2 \lambda_{n} \mathbf{z}_{v}^{2}-2 A_{u v} \mathbf{z}_{u} \mathbf{z}_{v}
\end{aligned}
$$

By Equation (11),

$$
\begin{aligned}
0 & \geq\left(\frac{\tilde{\mathbf{x}}^{T} \tilde{A} \tilde{\mathbf{x}}}{\tilde{\mathbf{x}}^{T} \tilde{\mathbf{x}}}-\frac{\tilde{\mathbf{z}}^{T} \tilde{A} \tilde{\mathbf{z}}}{\tilde{\mathbf{z}}^{T} \tilde{\mathbf{z}}}\right)-\left(\lambda_{1}-\lambda_{n}\right) \\
& =\left(\frac{\lambda_{1}-2 \lambda_{1} \mathbf{x}_{u}^{2}+2 \lambda_{1} \mathbf{x}_{v}^{2}-2 A_{u v} \mathbf{x}_{u} \mathbf{x}_{v}}{1-\mathbf{x}_{u}^{2}+\mathbf{x}_{v}^{2}}-\frac{\lambda_{n}-2 \lambda_{n} \mathbf{z}_{u}^{2}+2 \lambda_{n} \mathbf{z}_{v}^{2}-2 A_{u v} \mathbf{z}_{u} \mathbf{z}_{v}}{1-\mathbf{z}_{u}^{2}+\mathbf{z}_{v}^{2}}\right)-\left(\lambda_{1}-\lambda_{n}\right) \\
& =\frac{-\lambda_{1} \mathbf{x}_{u}^{2}+\lambda_{1} \mathbf{x}_{v}^{2}-2 A_{i j} \mathbf{x}_{u} \mathbf{x}_{v}}{1-\mathbf{x}_{u}^{2}+\mathbf{x}_{v}^{2}}-\frac{-\lambda_{n} \mathbf{z}_{u}^{2}+\lambda_{n} \mathbf{z}_{v}^{2}-2 A_{i j} \mathbf{z}_{u} \mathbf{z}_{v}}{1-\mathbf{z}_{u}^{2}+\mathbf{z}_{v}^{2}}
\end{aligned}
$$

Rearranging terms, we obtain
$\left(\lambda_{1} \mathbf{x}_{u}^{2}-\lambda_{1} \mathbf{x}_{v}^{2}\right)-\left(\lambda_{n} \mathbf{z}_{u}^{2}-\lambda_{n} \mathbf{z}_{v}^{2}\right) \geq \frac{\lambda_{n}\left(\mathbf{z}_{u}^{2}-\mathbf{z}_{v}^{2}\right)^{2}+2 A_{u, v} \mathbf{z}_{u} \mathbf{z}_{v}}{1-\mathbf{z}_{u}^{2}+\mathbf{z}_{v}^{2}}-\frac{\lambda_{1}\left(\mathbf{x}_{u}^{2}-\mathbf{x}_{v}^{2}\right)^{2}+2 A_{u, v} \mathbf{x}_{u} \mathbf{x}_{v}}{1-\mathbf{x}_{u}^{2}+\mathbf{x}_{v}^{2}}$.
By Lemma 2.3, we have that $\left|\mathbf{x}_{u}\right|,\left|\mathbf{x}_{v}\right|,\left|\mathbf{z}_{u}\right|$, and $\left|\mathbf{z}_{v}\right|$ are all $O\left(n^{-1 / 2}\right)$, and so the righthand side of the previous inequality is $O(1 / n)$. It follows that

$$
\left|\left(\lambda_{1} \mathbf{x}_{u}^{2}-\lambda_{1} \mathbf{x}_{v}^{2}\right)-\left(\lambda_{n} \mathbf{z}_{u}^{2}-\lambda_{n} \mathbf{z}_{v}^{2}\right)\right|<\frac{C}{n}
$$

for some absolute constant $C$. Rearranging terms gives the desired result.

## 3. THE SPREAD-EXTREMAL PROBLEM FOR GRAPHONS

Graphons (or graph functions) are analytical objects that may be used to study the limiting behavior of large, dense graphs. They were originally introduced in [6] and [19].
3.1. Introduction to graphons. Consider the set $\mathcal{W}$ of all bounded symmetric measurable functions $W:[0,1]^{2} \rightarrow[0,1]$ (by symmetric, we mean $W(x, y)=W(y, x)$ for all $\left.(x, y) \in[0,1]^{2}\right)$. A function $W \in \mathcal{W}$ is called a stepfunction if there is a partition of $[0,1]$ into subsets $S_{1}, S_{2}, \ldots, S_{m}$ such that $W$ is constant on every block $S_{i} \times S_{j}$. Every graph has a natural representation as a stepfunction in $\mathcal{W}$ taking values either 0 or 1 (such a graphon is referred to as a stepgraphon). In particular, given a graph $G$ on $n$ vertices indexed $\{1,2, \ldots, n\}$, we can define a measurable set $K_{G} \subseteq[0,1]^{2}$ as

$$
K_{G}=\bigcup_{u \sim v}\left[\frac{u-1}{n}, \frac{u}{n}\right] \times\left[\frac{v-1}{n}, \frac{v}{n}\right] .
$$

This represents the graph $G$ as a bounded symmetric measurable function $W_{G}$ that takes value 1 on $K_{G}$ and 0 everywhere else. For a measurable subset $U$, we will use $m(U)$ to denote its Lebesgue measure.

This representation of a graph as a measurable subset of $[0,1]^{2}$ lends itself to a visual presentation sometimes referred to as a "pixel picture"; see, for example, Figure for two representations of a bipartite graph as a measurable subset of $[0,1]^{2}$. Clearly, this indicates that such a representation is not unique; neither is the representation of a


FIGURE 2. Two presentations of a bipartite graph as a stepfunction
graph as a stepfunction. Using an equivalence relation on $\mathcal{W}$ derived from the socalled cut metric, we can identify graphons that are equivalent up to relabeling, and up to any differences on a set of measure zero (i.e. equivalent almost everywhere).

For all symmetric, bounded Lebesgue-measurable functions $W:[0,1]^{2} \rightarrow \mathbb{R}$, we let

$$
\|W\|_{\square}=\sup _{S, T \subseteq[0,1]}\left|\int_{S \times T} W(x, y) d x d y\right|
$$

Here, $\|\cdot\|_{\square}$ is referred to as the cut norm. Next, we can also define a semidistance $\delta_{\square}$ on $\mathcal{W}$ as follows. First, we define the weak isomorphism of graphons. Let $\mathcal{S}$ be the set of all measure-preserving functions on $[0,1]$. For every $\varphi \in \mathcal{S}$ and every $W \in \mathcal{W}$, define $W^{\varphi}:[0,1]^{2} \rightarrow[0,1]$ by

$$
W^{\varphi}(x, y):=W(\varphi(x), \varphi(y))
$$

for a.e. $(x, y) \in[0,1]^{2}$. Now for any $W_{1}, W_{2} \in \mathcal{W}$, let

$$
\delta_{\square}\left(W_{1}, W_{2}\right)=\inf _{\phi \in \mathcal{S}}\left\|W_{1}-W_{2} \circ \phi\right\|_{\square}
$$

Define the equivalence relation $\sim$ on $\mathcal{W}$ as follows: for all $W_{1}, W_{2} \in \mathcal{W}, W_{1} \sim W_{2}$ if and only if $\delta_{\square}\left(W_{1}, W_{2}\right)=0$. Furthermore, let $\hat{\mathcal{W}}:=\mathcal{W} / \sim$ be the quotient space of $\mathcal{W}$ under $\sim$. Note that $\delta_{\square}$ induces a metric on $\hat{\mathcal{W}}$. Crucially, by [20, Theorem 5.1], $\hat{\mathcal{W}}$ is a compact metric space.

Given $W \in \hat{\mathcal{W}}$, we define the Hilbert-Schmidt operator $A_{W}: \mathscr{L}^{2}[0,1] \rightarrow \mathscr{L}^{2}[0,1]$ by

$$
\left(A_{W} f\right)(x):=\int_{0}^{1} W(x, y) f(y) d y
$$

for all $f \in \mathscr{L}^{2}[0,1]$ and a.e. $x \in[0,1]$.
Since $W$ is symmetric and bounded, $A_{W}$ is a compact Hermitian operator. In particular, $A_{W}$ has a discrete, real spectrum $\Lambda(W)$ whose only possible accumulation point is 0 (cf. [5]), and so the maximum and minimum eigenvalues exist. Let $\mu(W)$ and $\nu(W)$ be the maximum and minimum eigenvalues of $A_{W}$, respectively, and define the spread of $W$ as

$$
\operatorname{spr}(W):=\mu(W)-\nu(W)
$$

By the Min-Max Theorem, we have that

$$
\mu(W)=\max _{\|f\|_{2}=1} \int_{0}^{1} \int_{0}^{1} W(x, y) f(x) f(y) d x d y
$$

and

$$
\nu(W)=\min _{\|f\|_{2}=1} \int_{0}^{1} \int_{0}^{1} W(x, y) f(x) f(y) d x d y
$$

Both $\mu$ and $\nu$ are continuous functions with respect to $\delta_{\square}$. In particular, we have the following.

Theorem 3.1 (cf. [6, Theorem 6.6] or [18, Theorem 11.54]). Let $\left\{W_{i}\right\}_{i}$ be a sequence of graphons converging to $W$ with respect to $\delta_{\square}$. Then as $n \rightarrow \infty$,

$$
\mu\left(W_{n}\right) \rightarrow \mu(W) \quad \text { and } \quad \nu\left(W_{n}\right) \rightarrow \nu(W) .
$$

If $W \sim W^{\prime}$ then $\mu(W)=\mu\left(W^{\prime}\right)$ and $\nu(W)=\nu\left(W^{\prime}\right)$. By compactness, we may consider the optimization problem on the factor space $\hat{\mathcal{W}}$

$$
\operatorname{spr}(\hat{\mathcal{W}})=\max _{W \in \hat{\mathcal{W}}} \operatorname{spr}(W)
$$

and there is a $W \in \hat{\mathcal{W}}$ that attains the maximum. Since every graph is represented by $W_{G} \in \hat{\mathcal{W}}$, this allows us to give an upper bound for $s(n)$ in terms of $\operatorname{spr}(\hat{\mathcal{W}})$. Indeed, by replacing the eigenvectors of $G$ with their corresponding stepfunctions, we can show Proposition 3.2.

Proposition 3.2. Let $G$ be a graph on $n$ vertices. Then

$$
\lambda_{1}(G)=n \cdot \mu\left(W_{G}\right) \quad \text { and } \quad \lambda_{n}(G)=n \cdot \nu\left(W_{G}\right)
$$

Proposition 3.2 implies that $s(n) \leq n \cdot \operatorname{spr}(\hat{\mathcal{W}})$ for all $n$. Combined with Theorem 1.4, this gives Corollary 3.3(a similar argument implies Corollary 1.5).

Corollary 3.3. For all $n, s(n) \leq \frac{2 n}{\sqrt{3}}$.
This can be proved more directly using Theorem 1.1 and taking tensor powers.
3.2. Properties of spread-extremal graphons. Our main objective in the following two sections is to solve the maximum spread problem for graphons, in order to produce a tight estimate for $s(n)$. In this subsection, we prove some preliminary results that are largely a translation of what is known in the graph setting (see Section (2). First, we define what it means for a graphon to be connected, and show that spread-extremal graphons must be connected. We then prove a standard corollary of the Perron-Frobenius theorem. Finally, we prove graphon versions of Lemma 2.1 and Lemma 2.4.

Let $W_{1}$ and $W_{2}$ be graphons, and let $\alpha_{1}, \alpha_{2}$ be positive real numbers with $\alpha_{1}+\alpha_{2}=1$. We define the direct sum of $W_{1}$ and $W_{2}$ with weights $\alpha_{1}$ and $\alpha_{2}$, denoted $W=\alpha_{1} W_{1} \oplus$ $\alpha_{2} W_{2}$, as follows. Let $\varphi_{1}$ and $\varphi_{2}$ be the increasing affine maps that send $J_{1}:=\left[0, \alpha_{1}\right]$ and $J_{2}:=\left[\alpha_{1}, 1\right]$ to $[0,1]$, respectively. Then for all $(x, y) \in[0,1]^{2}$, let

$$
W(x, y):=\left\{\begin{array}{cl}
W_{i}\left(\varphi_{i}(x), \varphi_{i}(y)\right), & \text { if }(x, y) \in J_{i} \times J_{i} \text { for some } i \in\{1,2\} \\
0, & \text { otherwise }
\end{array} .\right.
$$

A graphon $W$ is connected if $W$ is not weakly isomorphic to a direct sum $\alpha_{1} W_{1} \oplus \alpha_{2} W_{2}$ where $\alpha_{1} \neq 0,1$. Equivalently, $W$ is connected if there does not exist a measurable subset $A \subseteq[0,1]$ of positive measure such that $W(x, y)=0$ for a.e. $(x, y) \in A \times A^{c}$.

Proposition 3.4. Suppose $W_{1}, W_{2}$ are graphons and $\alpha_{1}, \alpha_{2}$ are positive real numbers summing to 1 . Let $W:=\alpha_{1} W_{1} \oplus \alpha_{2} W_{2}$. Then as multisets,

$$
\Lambda(W)=\left\{\alpha_{1} u: u \in \Lambda\left(W_{1}\right)\right\} \cup\left\{\alpha_{2} v: v \in \Lambda\left(W_{2}\right)\right\} .
$$

Moreover, $\operatorname{spr}(W) \leq \alpha_{1} \operatorname{spr}\left(W_{1}\right)+\alpha_{2} \operatorname{spr}\left(W_{2}\right)$ with equality if and only if $W_{1}$ or $W_{2}$ is the all-zeroes graphon.

Proof. For convenience, let $\Lambda_{i}:=\left\{\alpha_{i} u: u \in \Lambda\left(W_{i}\right)\right\}$ for each $i \in\{1,2\}$ and $\Lambda:=$ $\Lambda(W)$. The first claim holds simply by considering the restriction of eigenfunctions to the intervals $\left[0, \alpha_{1}\right]$ and $\left[\alpha_{1}, 1\right]$.

For the second claim, we first write $\operatorname{spr}(W)=\alpha_{i} \mu\left(W_{i}\right)-\alpha_{j} \nu\left(W_{j}\right)$, where $i, j \in\{1,2\}$. Let $I_{i}:=\left[\min \left(\Lambda_{i}\right), \max \left(\Lambda_{i}\right)\right]$ for each $i \in\{1,2\}$ and $I:=[\min (\Lambda), \max (\Lambda)]$. Clearly $\alpha_{i} \operatorname{spr}\left(W_{i}\right)=\operatorname{diam}\left(I_{i}\right)$ for each $i \in\{1,2\}$ and $\operatorname{spr}(W)=\operatorname{diam}(I)$. Moreover, $I=I_{1} \cup I_{2}$. Since $0 \in I_{1} \cap I_{2}$, $\operatorname{diam}(I) \leq \operatorname{diam}\left(I_{1}\right)+\operatorname{diam}\left(I_{2}\right)$ with equality if and only if either $I_{1}$ or $I_{2}$ equals $\{0\}$. So the desired claim holds.

Furthermore, the following basic corollary of the Perron-Frobenius holds. For completeness, we prove it here.

Proposition 3.5. Let $W$ be a connected graphon, and write $f$ for an eigenfunction corresponding to $\mu(W)$. Then $f$ is nonzero with constant sign a.e.
Proof. Let $\mu=\mu(W)$. Since

$$
\mu=\max _{\|h\|_{2}=1} \int_{(x, y) \in[0,1]^{2}} W(x, y) h(x) h(y) d x d y
$$

it follows without loss of generality that $f \geq 0$ a.e. on $[0,1]$. Let $Z:=\{x \in[0,1]$ : $f(x)=0\}$. Then for a.e. $x \in Z$,

$$
0=\mu f(x)=\int_{y \in[0,1]} W(x, y) f(y) d y=\int_{y \in Z^{c}} W(x, y) f(y) d y
$$

Since $f>0$ on $Z^{c}$, it follows that $W(x, y)=0$ a.e. on $Z \times Z^{c}$. Clearly $m\left(Z^{c}\right) \neq 0$. If $m(Z)=0$ then the desired claim holds, so without loss of generality, $0<m(Z), m\left(Z^{c}\right)<$ 1. It follows that $W$ is disconnected, a contradiction to our assumption, which completes the proof.

We may now prove a graphon version of Lemma 2.1.
Lemma 3.6. Suppose $W$ is a graphon achieving maximum spread, and let $f, g$ be eigenfunctions for the maximum and minimum eigenvalues for $W$, respectively. Then the following claims hold:
(i) For a.e. $(x, y) \in[0,1]^{2}$,

$$
W(x, y)= \begin{cases}1, & \text { if } f(x) f(y)>g(x) g(y) \\ 0, & \text { otherwise }\end{cases}
$$

(ii) $f(x) f(y)-g(x) g(y) \neq 0$ for a.e. $(x, y) \in[0,1]^{2}$.

Proof. We proceed in the following order:

- Prove Item (i) holds for a.e. $(x, y) \in[0,1]^{2}$ such that $f(x) f(y) \neq g(x) g(y)$. We will call this Item (ii)*.
- Prove Item (iii).
- Deduce Item (i) also holds.

By Propositions 3.4 and 3.5, we may assume without loss of generality that $f>0$ a.e. on $[0,1]$. For convenience, we define the quantity $d(x, y):=f(x) f(y)-g(x) g(y)$. To prove Item (ii $)^{*}$, we first define a graphon $W^{\prime}$ by

$$
W^{\prime}(x, y)=\left\{\begin{aligned}
1, & \text { if } d(x, y)>0 \\
0, & \text { if } d(x, y)<0 \\
W(x, y), & \text { otherwise }
\end{aligned}\right.
$$

Then by inspection,

$$
\begin{aligned}
\operatorname{spr}\left(W^{\prime}\right) \geq & \int_{(x, y) \in[0,1]^{2}} W^{\prime}(x, y)(f(x) f(y)-g(x) g(y)) d x d y \\
= & \int_{(x, y) \in[0,1]^{2}} W(x, y)(f(x) f(y)-g(x) g(y)) d x d y \\
& +\int_{d(x, y)>0}(1-W(x, y)) d(x, y) d x d y-\int_{d(x, y)<0} W(x, y) d(x, y) d x d y \\
= & \operatorname{spr}(W)+\int_{d(x, y)>0}(1-W(x, y)) d(x, y) d x d y-\int_{d(x, y)<0} W(x, y) d(x, y) d x d y .
\end{aligned}
$$

Since $W$ maximizes spread, both integrals in the last line must be 0 , and hence Item (ii)* holds.

Now, we prove Item (iii). For convenience, we define $U$ to be the set of all pairs $(x, y) \in[0,1]^{2}$ so that $d(x, y)=0$. Now let $\tilde{W}$ be any graphon that differs from $W$ only on $U$. Then

$$
\begin{aligned}
\operatorname{spr}(\tilde{W}) \geq & \int_{(x, y) \in[0,1]^{2}} \tilde{W}(x, y)(f(x) f(y)-g(x) g(y)) d x d y \\
= & \int_{(x, y) \in[0,1]^{2}} W(x, y)(f(x) f(y)-g(x) g(y)) d x d y \\
& +\int_{(x, y) \in U}(\tilde{W}(x, y)-W(x, y))(f(x) f(y)-g(x) g(y)) d x d y \\
= & \operatorname{spr}(W) .
\end{aligned}
$$

Since $\operatorname{spr}(W) \geq \operatorname{spr}(\tilde{W}), f$ and $g$ are eigenfunctions for $\tilde{W}$, and we may write $\tilde{\mu}$ and $\tilde{\nu}$ for the corresponding eigenvalues. Now, we define

$$
\begin{aligned}
I_{\tilde{W}}(x) & :=(\tilde{\mu}-\mu) f(x) \\
& =\int_{y \in[0,1]}(\tilde{W}(x, y)-W(x, y)) f(y) d y \\
& =\int_{y \in[0,1],(x, y) \in U}(\tilde{W}(x, y)-W(x, y)) f(y) d y .
\end{aligned}
$$

Similarly, we define

$$
\begin{aligned}
J_{\tilde{W}}(x) & :=(\tilde{v}-v) g(x) \\
& =\int_{y \in[0,1]}(\tilde{W}(x, y)-W(x, y)) g(y) d y \\
& =\int_{y \in[0,1],(x, y) \in U}(\tilde{W}(x, y)-W(x, y)) g(y) d y .
\end{aligned}
$$

Since $f$ and $g$ are orthogonal,

$$
\int_{x \in[0,1]} I_{\tilde{W}}(x) J_{\tilde{W}}(x) d x=0
$$

By definition of $U$, we have that for a.e. $(x, y) \in U, 0=d(x, y)=f(x) f(y)-g(x) g(y)$. In particular, since $f(x), f(y)>0$ for a.e. $(x, y) \in[0,1]^{2}$, then a.e. $(x, y) \in U$ has $g(x) g(y)>0$. So by letting

$$
\begin{aligned}
& U_{+}:=\{(x, y) \in U: g(x), g(y)>0\}, \\
& U_{-}:=\{(x, y) \in U: g(x), g(y)<0\}, \text { and } \\
& U_{0}:=U \backslash\left(U_{+} \cup U_{-}\right),
\end{aligned}
$$

$U_{0}$ has measure 0.
First, let us fix $\tilde{W}$ to be the graphon defined by

$$
\tilde{W}(x, y)=\left\{\begin{array}{rl}
1, & \text { if }(x, y) \in U_{+} \\
W(x, y), & \text { otherwise }
\end{array} .\right.
$$

For this choice of $\tilde{W}$,

$$
\begin{aligned}
& I_{\tilde{W}}(x)=\int_{y \in[0,1],(x, y) \in U_{+}}(1-W(x, y)) f(y) d y, \text { and } \\
& J_{\tilde{W}}(x)=\int_{y \in[0,1],(x, y) \in U_{+}}(1-W(x, y)) g(y) d y
\end{aligned}
$$

Clearly $I_{\tilde{W}}$ and $J_{\tilde{W}}$ are nonnegative functions, so $I_{\tilde{W}}(x) J_{\tilde{W}}(x)=0$ for a.e. $x \in[0,1]$. Since $f(y)$ and $g(y)$ are positive for a.e. $(x, y) \in U, W(x, y)=1$ for a.e. on $U_{+}$.

If instead we fix $\tilde{W}(x, y)$ to be 0 for all $(x, y) \in U_{+}$, it follows by a similar argument that $W(x, y)=0$ for a.e. $(x, y) \in U_{+}$. So $U_{+}$has measure 0 . Repeating the same argument on $U_{-}$, we similarly conclude that $U_{-}$has measure 0 . This completes the proof of Item (iii).

Finally we note that Items (ii)* and (iii) together imply Item (ii).
From here, it is easy to see that any graphon maximizing the spread is a join of two threshold graphons. Next we prove the graphon version of Lemma 2.4.

Lemma 3.7. If $W$ is a graphon achieving the maximum spread with corresponding eigenfunctions $f, g$, then $\mu f^{2}-\nu g^{2}=\mu-v$ almost everywhere.
Proof. We will use the notation $(x, y) \in W$ to denote that $(x, y) \in[0,1]^{2}$ satisfies $W(x, y)=1$. Let $\varphi:[0,1] \rightarrow[0,1]$ be an arbitrary homeomorphism that is orientationpreserving in the sense that $\varphi(0)=0$ and $\varphi(1)=1$. Then $\varphi$ is a continuous strictly monotone increasing function which is differentiable almost everywhere. Now let $\tilde{f}:=$
$\varphi^{\prime} \cdot(f \circ \varphi), \tilde{g}:=\varphi^{\prime} \cdot(g \circ \varphi)$ and $\tilde{W}:=\left\{(x, y) \in[0,1]^{2}:(\varphi(x), \varphi(y)) \in W\right\}$. Using the substitutions $u=\varphi(x)$ and $v=\varphi(y)$,

$$
\begin{aligned}
\tilde{f} \tilde{W} \tilde{f} & =\int_{(x, y) \in[0,1]^{2}} \chi_{(\varphi(x), \varphi(y)) \in \tilde{W}} \varphi^{\prime}(x) \varphi^{\prime}(y) \cdot f(\varphi(x)) f(\varphi(y)) d x d y \\
& =\int_{(x, y) \in[0,1]^{2}} \chi_{(x, y) \in W} f(u) f(v) d u d v \\
& =\mu .
\end{aligned}
$$

Similarly, $\tilde{g} \tilde{W} \tilde{g}=\nu$.
Note however that the $L_{2}$ norms of $\tilde{f}, \tilde{g}$ may not be 1 . Indeed using the substitution $u=\varphi(x)$,

$$
\|\tilde{f}\|_{2}^{2}=\int_{x \in[0,1]} \varphi^{\prime}(x)^{2} f(\varphi(x))^{2} d x=\int_{u \in[0,1]} \varphi^{\prime}\left(\varphi^{-1}(u)\right) \cdot f(u)^{2} d u .
$$

We exploit this fact as follows. Suppose $I, J$ are disjoint subintervals of $[0,1]$ of the same positive length $m(I)=m(J)=\ell>0$, and for any $\varepsilon>0$ sufficiently small (in terms of $\ell$ ), let $\varphi$ be the (unique) piecewise linear function that stretches $I$ to length $(1+\varepsilon) m(I)$, shrinks $J$ to length $(1-\varepsilon) m(J)$, and shifts only the elements in between $I$ and $J$. Note that for a.e. $x \in[0,1]$,

$$
\varphi^{\prime}(x)=\left\{\begin{aligned}
1+\varepsilon, & \text { if } x \in I \\
1-\varepsilon, & \text { if } x \in J, \\
1, & \text { otherwise }
\end{aligned}\right.
$$

Again with the substitution $u=\varphi(x)$,

$$
\begin{aligned}
\|\tilde{f}\|_{2}^{2} & =\int_{x \in[0,1]} \varphi^{\prime}(x)^{2} \cdot f(\varphi(x))^{2} d x \\
& =\int_{[u \in[0,1]} \varphi^{\prime}\left(\varphi^{-1}(u)\right) f(u)^{2} d u \\
& =1+\varepsilon \cdot\left(\left\|\chi_{I} f\right\|_{2}^{2}-\left\|\chi_{J} f\right\|_{2}^{2}\right)
\end{aligned}
$$

The same equality holds for $\tilde{g}$ instead of $\tilde{f}$. After normalizing $\tilde{f}$ and $\tilde{g}$, by optimality of $W$, we get a difference of Rayleigh quotients as

$$
\begin{aligned}
0 & \leq(f W f-g W g)-\frac{\tilde{f} \tilde{W} \tilde{f}}{\|\tilde{f}\|_{2}^{2}}-\frac{\tilde{g} \tilde{W} \tilde{g}}{\|\tilde{g}\|_{2}^{2}} \\
& =\frac{\mu \varepsilon \cdot\left(\left\|\chi_{I} f\right\|_{2}^{2}-\left\|\chi_{J} f\right\|_{2}^{2}\right)}{1+\varepsilon \cdot\left(\left\|\chi_{I} f\right\|_{2}^{2}-\left\|\chi_{J} f\right\|_{2}^{2}\right)}-\frac{\nu \varepsilon \cdot\left(\left\|\chi_{I} g\right\|_{2}^{2}-\left\|\chi_{J} g\right\|_{2}^{2}\right)}{1+\varepsilon \cdot\left(\left\|\chi_{I} g\right\|_{2}^{2}-\left\|\chi_{J} g\right\|_{2}^{2}\right)} \\
& =(1+o(1)) \varepsilon \cdot\left(\int_{I}\left(\mu f(x)^{2}-\nu g(x)^{2}\right) d x-\int_{J}\left(\mu f(x)^{2}-\nu g(x)^{2}\right) d x\right)
\end{aligned}
$$

as $\varepsilon \rightarrow 0$. It follows that for all disjoint intervals $I, J \subseteq[0,1]$ of the same length, the corresponding integrals are the same. Taking finer and finer partitions of $[0,1]$, it follows that the integrand $\mu f(x)^{2}-\nu g(x)^{2}$ is constant almost everywhere. Since the average of this quantity over all $[0,1]$ is $\mu-\nu$, the desired claim holds.

## 4. FROM GRAPHONS TO STEPGRAPHONS

The main result of this section is as follows.
Theorem 4.1. Suppose $W$ maximizes $\operatorname{spr}(\hat{\mathcal{W}})$. Then $W$ is a stepfunction taking values 0 and 1 of the following form


Furthermore, the internal divisions separate according to the sign of the eigenfunction corresponding to the minimum eigenvalue of $W$.

We begin Section 4.1 by mirroring the argument in [30], which proved a conjecture of Nikiforov regarding the largest eigenvalue of a graph and its complement, $\mu+\bar{\mu}$. There, Terpai showed that performing two operations on graphons leads to a strict increase in $\mu+\bar{\mu}$. Furthermore based on previous work of Nikiforov from [24], the conjecture for graphs reduced directly to maximizing $\mu+\bar{\mu}$ for graphons. Using these operations, Terpai [30] reduced to a $4 \times 4$ stepgraphon and then completed the proof by hand.

In our case, we are not so lucky and are left with a $7 \times 7$ stepgraphon after performing similar but more technical operations, detailed in this section. In order to reduce to a $3 \times 3$ stepgraphon, we make use of interval arithmetic (see Section 5.2 and Appendices A and $B$ ). Furthermore, our proof requires an additional technical argument to translate the result for graphons (Theorem 5.1) to our main result for graphs (Theorem 1.1). In Section 4.2, we prove Theorem 4.1.
4.1. Averaging. For convenience, we introduce some terminology. For any graphon $W$ with $\lambda$-eigenfunction $h$, we say that $x \in[0,1]$ is typical (with respect to $W$ and $h$ ) if

$$
\lambda \cdot h(x)=\int_{y \in[0,1]} W(x, y) h(y) d y
$$

Note that a.e. $x \in[0,1]$ is typical. Additionally if $U \subseteq[0,1]$ is measurable with positive measure, then we say that $x_{0} \in U$ is average (on $U$, with respect to $W$ and $h$ ) if

$$
h\left(x_{0}\right)^{2}=\frac{1}{m(U)} \int_{y \in U} h(y)^{2} d y
$$

Given $W, h, U$, and $x_{0}$ as above, we define the $L_{2}[0,1]$ function $\mathrm{av}_{U, x_{0}} h$ by setting

$$
\left(\operatorname{av}_{U, x_{0}} h\right)(x):=\left\{\begin{array}{rl}
h\left(x_{0}\right), & \text { if } x \in U \\
h(x), & \text { otherwise }
\end{array} .\right.
$$

Clearly $\left\|\mathrm{av}_{U, x_{0}} h\right\|_{2}=\|h\|_{2}$. Additionally, we define the graphon $\mathrm{av}_{U, x_{0}} W$ by setting

$$
\operatorname{av}_{U, x_{0}} W(x, y):=\left\{\begin{array}{rl}
0, & (x, y) \in U \times U \\
W\left(x_{0}, y\right), & (x, y) \in U \times U^{c} \\
W\left(x, x_{0}\right), & (x, y) \in U^{c} \times U \\
W(x, y), & (x, y) \in U^{c} \times U^{c}
\end{array} .\right.
$$

In the graph setting, this is analogous to replacing $U$ with an independent set whose vertices are clones of $x_{0}$. Lemma 4.2 indicates how this cloning affects the eigenvalues.

Lemma 4.2. Suppose $W$ is a graphon with a $\lambda$-eigenfunction $h$, and suppose there exist disjoint measurable subsets $U_{1}, U_{2} \subseteq[0,1]$ of positive measures $\alpha$ and $\beta$, respectively. Let $U:=U_{1} \cup U_{2}$. Moreover, suppose $W=0$ a.e. on $(U \times U) \backslash\left(U_{1} \times U_{1}\right)$. Additionally, suppose $x_{0} \in U_{2}$ is typical and average on $U$, with respect to $W$ and $h$. Let $\tilde{h}:=\operatorname{av}_{U, x_{0}} h$ and $\tilde{W}:=\mathrm{av}_{U, x_{0}} W$. Then for a.e. $x \in[0,1]$,
(2)

$$
\left(A_{\tilde{W}} \tilde{h}\right)(x)=\lambda \tilde{h}(x)+\left\{\begin{array}{rl}
0, & \text { if } x \in U \\
m(U) \cdot W\left(x_{0}, x\right) h\left(x_{0}\right)-\int_{y \in U} W(x, y) h(y) d y, & \text { otherwise }
\end{array} .\right.
$$

Furthermore,

$$
\begin{equation*}
\left\langle A_{\tilde{W}} \tilde{h}, \tilde{h}\right\rangle=\lambda+\int_{(x, y) \in U_{1} \times U_{1}} W(x, y) h(x) h(y) d x d y \tag{3}
\end{equation*}
$$

Proof. We first prove Equation (2). Note that for a.e. $x \in U$,

$$
\begin{aligned}
\left(A_{\tilde{W}} \tilde{h}\right)(x) & =\int_{y \in[0,1]} \tilde{W}(x, y) \tilde{h}(y) d y \\
& =\int_{y \in U} \tilde{W}(x, y) \tilde{h}(y) d y+\int_{y \in[0,1] \backslash U} \tilde{W}(x, y) \tilde{h}(y) d y \\
& =\int_{y \in[0,1] \backslash U} W\left(x_{0}, y\right) h(y) d y \\
& =\int_{y \in[0,1]} W\left(x_{0}, y\right) h(y) d y-\int_{y \in U} W\left(x_{0}, y\right) h(y) d y \\
& =\lambda h\left(x_{0}\right) \\
& =\lambda \tilde{h}(x),
\end{aligned}
$$

as desired. Now note that for a.e. $x \in[0,1] \backslash U$,

$$
\begin{aligned}
\left(A_{\tilde{W}} \tilde{h}\right)(x) & =\int_{y \in[0,1]} \tilde{W}(x, y) \tilde{h}(y) d y \\
& =\int_{y \in U} \tilde{W}(x, y) \tilde{h}(y) d y+\int_{y \in[0,1] \backslash U} \tilde{W}(x, y) \tilde{h}(y) d y \\
& =\int_{y \in U} W\left(x_{0}, x\right) h\left(x_{0}\right) d y+\int_{y \in[0,1] \backslash U} W(x, y) h(y) d y \\
& =m(U) \cdot W\left(x_{0}, x\right) h\left(x_{0}\right)+\int_{y \in[0,1]} W(x, y) h(y) d y-\int_{y \in U} W(x, y) h(y) d y \\
& =\lambda h(x)+m(U) \cdot W\left(x_{0}, x\right) h\left(x_{0}\right)-\int_{y \in U} W(x, y) h(y) d y .
\end{aligned}
$$

So again, the claim holds and this completes the proof of Equation (2). Now we prove Equation (3). Indeed by Equation (2),

$$
\begin{aligned}
\left\langle\left(A_{\tilde{W}} \tilde{h}\right), \tilde{h}\right\rangle= & \int_{x \in[0,1]}\left(A_{\tilde{W}} \tilde{h}\right)(x) \tilde{h}(x) d x \\
= & \int_{x \in[0,1]} \lambda \tilde{h}(x)^{2} d x \\
& +\int_{x \in[0,1] \backslash U}\left(m(U) \cdot W\left(x_{0}, x\right) h\left(x_{0}\right)-\int_{y \in U} W(x, y) h(y) d y\right) \cdot h(x) d x \\
= & \lambda+m(U) \cdot h\left(x_{0}\right)\left(\int_{x \in[0,1]} W\left(x_{0}, x\right) h(x) d x-\int_{x \in U} W\left(x_{0}, x\right) h(x) d x\right) \\
& -\int_{y \in U}\left(\int_{x \in[0,1]} W(x, y) h(x) d x-\int_{x \in U} W(x, y) h(x) d x\right) \cdot h(y) d y \\
= & \lambda+m(U) \cdot h\left(x_{0}\right)\left(\lambda h\left(x_{0}\right)-\int_{y \in U} 0 d y\right) \\
& -\int_{y \in U}\left(\lambda h(y)^{2}-\int_{x \in U} W(x, y) h(x) h(y) d x\right) d y \\
= & \lambda+\lambda m(U) \cdot h\left(x_{0}\right)^{2}-\lambda \int_{y \in U} h(y)^{2} d y+\int_{(x, y) \in U \times U} W(x, y) h(x) h(y) d x d y \\
= & \lambda+\int_{(x, y) \in U_{1} \times U_{1}} W(x, y) h(x) h(y) d x d y,
\end{aligned}
$$

and this completes the proof of desired claims.
We have the following useful corollary.
Corollary 4.3. Suppose $\operatorname{spr}(W)=\operatorname{spr}(\hat{\mathcal{W}})$ with maximum and minimum eigenvalues $\mu, \nu$ corresponding respectively to eigenfunctions $f, g$. Moreover, suppose that there exist disjoint subsets $A, B \subseteq[0,1]$ and $x_{0} \in B$ so that the conditions of Lemma 4.2 are met for $W$ with $\lambda=\mu, h=f, U_{1}=A$, and $U_{2}=B$. Then,
(i) $W(x, y)=0$ for a.e. $(x, y) \in U^{2}$, and
(ii) $f$ is constant on $U$.

Proof. Without loss of generality, we assume that $\|f\|_{2}=\|g\|_{2}=1$. Write $\tilde{W}$ for the graphon and $\tilde{f}, \tilde{g}$ for the corresponding functions produced by Lemma 4.2. By Proposition 3.5, we may assume without loss of generality that $f>0$ a.e. on $[0,1]$. We first prove Item (ii). Note that

$$
\begin{align*}
\operatorname{spr}(\tilde{W}) & \geq \int_{(x, y) \in[0,1]^{2}} \tilde{W}(x, y)(\tilde{f}(x) \tilde{f}(y)-\tilde{g}(x) \tilde{g}(y)) d x d y \\
& =(\mu-\nu)+\int_{(x, y) \in A \times A} W(x, y)(f(x) f(y)-g(x) g(y)) d x d y \\
& =\operatorname{spr}(W)+\int_{(x, y) \in A \times A} W(x, y)(f(x) f(y)-g(x) g(y)) d x d y . \tag{4}
\end{align*}
$$

Since $\operatorname{spr}(W) \geq \operatorname{spr}(\tilde{W})$ and by Lemma 3.6(iii), $f(x) f(y)-g(x) g(y)>0$ for a.e. $(x, y) \in$ $A \times A$ such that $W(x, y) \neq 0$. Item (i) follows.

For Item (iii), we first note that $f$ is a $\mu$-eigenfunction for $\tilde{W}$. Indeed, if not, then the inequality in (4) holds strictly, a contradiction to the fact that $\operatorname{spr}(W) \geq \operatorname{spr}(\tilde{W})$. Again by Lemma 4.2,

$$
m(U) \cdot W\left(x_{0}, x\right) f\left(x_{0}\right)=\int_{y \in U} W(x, y) f(y) d y
$$

for a.e. $x \in[0,1] \backslash U$. Let $S_{1}:=\left\{x \in[0,1] \backslash U: W\left(x_{0}, x\right)=1\right\}$ and $S_{0}:=[0,1] \backslash\left(U \cup S_{1}\right)$. We claim that $m\left(S_{1}\right)=0$. Assume otherwise. By Lemma 4.2 and by Cauchy-Schwarz, for a.e. $x \in S_{1}$

$$
\begin{aligned}
m(U) \cdot f\left(x_{0}\right) & =m(U) \cdot W\left(x_{0}, x\right) f\left(x_{0}\right) \\
& =\int_{y \in U} W(x, y) f(y) d y \\
& \leq \int_{y \in U} f(y) d y \\
& \leq m(U) \cdot f\left(x_{0}\right)
\end{aligned}
$$

and by sandwiching, $W(x, y)=1$ and $f(y)=f\left(x_{0}\right)$ for a.e. $y \in U$. Since $m\left(S_{1}\right)>0$, it follows that $f(y)=f\left(x_{0}\right)=0$ for a.e. $y \in U$, as desired.

So we assume otherwise, that $m\left(S_{1}\right)=0$. Then for a.e. $x \in[0,1] \backslash U, W\left(x_{0}, x\right)=0$ and

$$
0=m(U) \cdot W\left(x_{0}, x\right) f\left(x_{0}\right)=\int_{y \in U} W(x, y) f(y) d y
$$

and since $f>0$ a.e. on $[0,1]$, it follows that $W(x, y)=0$ for a.e. $y \in U$. Altogether, $W(x, y)=0$ for a.e. $(x, y) \in([0,1] \backslash U) \times U$. So $W$ is a disconnected, a contradiction to Proposition 3.4. The desired claim holds.

### 4.2. Proof of Theorem 4.1.

Proof. For convenience, we write $\mu:=\mu(W)$ and $\nu:=\nu(W)$ and let $f, g$ denote the corresponding unit eigenfunctions. By Proposition 3.5, we may assume without loss of generality that $f>0$.

First, we show without loss of generality that $f, g$ are monotone on the sets $P:=$ $\{x \in[0,1]: g(x) \geq 0\}$ and $N:=[0,1] \backslash P$. Indeed, we define a total ordering $\leqslant$ on $[0,1]$ as follows. For all $x$ and $y$, we let $x \leqslant y$ if:
(i) $g(x) \geq 0$ and $g(y)<0$, or
(ii) Item (i) does not hold and $f(x)>f(y)$, or
(iii) Item (i) does not hold, $f(x)=f(y)$, and $x \leq y$.

By inspection, the function $\varphi:[0,1] \rightarrow[0,1]$ defined by

$$
\varphi(x):=m(\{y \in[0,1]: y \leqslant x\})
$$

is a weak isomorphism between $W$ and its entrywise composition with $\varphi$. By invariance of $\operatorname{spr}(\cdot)$ under weak isomorphism, we make the above replacement and write $f, g$ for the replacement eigenfunctions. That is, we are assuming that our graphon is relabeled so that $[0,1]$ respects $\leqslant$.

As above, let $P:=\{x \in[0,1]: g(x) \geq 0\}$ and $N:=[0,1] \backslash P$. By Lemma 3.7, $f$ and $-g$ are monotone nonincreasing on $P$. Additionally, $f$ and $g$ are monotone nonincreasing on $N$. Without loss of generality, we may assume that $W$ is of the form from Lemma 3.6. Now we let $S:=\{x \in[0,1]: f(x)<|g(x)|\}$ and $C:=[0,1] \backslash S$. By Lemma 3.6 we have that $W(x, y)=1$ for almost every $x, y \in C$ and $W(x, y)=0$ for almost every $x, y \in S \cap P$ or $x, y \in S \cap N$. We have used the notation $C$ and $S$ because the analogous sets in the graph setting form a clique or a stable set respectively. We first prove Claim A.

Claim A. Except on a set of measure $0, f$ takes on at most 2 values on $P \cap S$, and at most 2 values on $N \cap S$.

We first prove Claim $\Delta$ for $f$ on $P \cap S$. Let $D$ be the set of all discontinuities of $f$ on the interior of the interval $P \cap S$. Clearly $D$ consists only of jump-discontinuities. By the Darboux-Froda Theorem, $D$ is at most countable and moreover, $(P \cap S) \backslash D$ is a union of at most countably many disjoint intervals $\mathcal{J}$. Moreover, $f$ is continuous on the interior of each $I \in \mathcal{J}$.

We show now that $f$ is piecewise constant on the interiors of each $I \in \mathcal{J}$. Indeed, let $I \in \mathcal{J}$. Since $f$ is a $\mu$-eigenfunction function for $W$,

$$
\mu f(x)=\int_{y \in[0,1]} W(x, y) f(y) d y
$$

for a.e. $x \in[0,1]$. By continuity of $f$ on the interior of $I$, this equation holds everywhere on the interior of $I$. Additionally, since $f$ is continuous on the interior of $I$, by the Mean Value Theorem, there exists some $x_{0}$ in the interior of $I$ so that

$$
f\left(x_{0}\right)^{2}=\frac{1}{m(U)} \int_{x \in U} f(x)^{2} d x
$$

By Corollary 4.3, $f$ is constant on the interior of $U$, as desired.
If $|\mathcal{J}| \leq 2$, the desired claim holds, so we may assume otherwise. Then there exists distinct $I_{1}, I_{2}, I_{3} \in \mathcal{J}$. Moreover, $f$ equals a constant $f_{1}, f_{2}, f_{3}$ on the interiors of $I_{1}, I_{2}$, and $I_{3}$, respectively. Additionally, since $I_{1}, I_{2}$, and $I_{3}$ are separated from each other by at least one jump discontinuity, we may assume without loss of generality that $f_{1}<$ $f_{2}<f_{3}$. It follows that there exists a measurable subset $U \subseteq I_{1} \cup I_{2} \cup I_{3}$ of positive
measure so that

$$
f_{2}^{2}=\frac{1}{m(U)} \int_{x \in U} f(x)^{2} d x
$$

By Corollary 4.3, $f$ is constant on $U$, a contradiction. So Claim Aholds on $P \cap S$. For Claim $A$ on $N \cap S$, we may repeat this argument with $P$ and $N$ interchanged, and $g$ and $-g$ interchanged.

Now we show Claim B.
Claim B. For a.e. $(x, y) \in(P \times P) \cup(N \times N)$ such that $f(x) \geq f(y)$, we have that for a.e. $z \in[0,1], W(x, z)=0$ implies that $W(y, z)=0$.

We first prove Claim Bfor a.e. $(x, y) \in P \times P$. Suppose $W(y, z)=0$. By Lemma 3.6, in this case $z \in P$. Then for a.e. such $x, y$, by Lemma 3.7, $g(x) \leq g(y)$. By Lemma 3.6(ii), $W(x, z)=0$ implies that $f(x) f(z)<g(x) g(z)$. Since $f(x) \geq f(y)$ and $g(x) \leq g(y)$, $f(y) f(z)<g(y) g(z)$. Again by Lemma 3.6(il), $W(y, z)=0$ for a.e. such $x, y, z$, as desired. So the desired claim holds for a.e. $(x, y) \in P \times P$ such that $f(x) \geq f(y)$. We may repeat the argument for a.e. $(x, y) \in N \times N$ to arrive at the same conclusion.

Claim Collows directly from Lemma 3.7.
Claim C. For a.e. $x \in[0,1], x \in C$ if and only if $f(x) \geq 1$, if and only if $|g(x)| \leq 1$.
Finally, we show Claim D.
Claim D. Except on a set of measure $0, f$ takes on at most 3 values on $P \cap C$, and at most 3 values on $N \cap C$.

For a proof, we first write $P \cap S=S_{1} \cup S_{2}$ so that $S_{1}, S_{2}$ are disjoint, $f$ equals some constant $f_{1}$ a.e. on $S_{1}$, and $f$ equals some constant $f_{2}$ a.e. on $S_{2}$. By Lemma 3.7, $g$ equals some constant $g_{1}$ a.e. on $S_{1}$ and $g$ equals some constant $g_{2}$ a.e. on $S_{2}$. By definition of $P, g_{1}, g_{2} \geq 0$. Now suppose $x \in P \cap C$ so that

$$
\mu f(x)=\int_{y \in[0,1]} W(x, y) f(y) d y
$$

Then by Lemma 3.6(ii),

$$
\mu f(x)=\int_{y \in(P \cap C) \cup N} f(y) d y+\int_{y \in S_{1}} W(x, y) f(y) d y+\int_{y \in S_{2}} W(x, y) f(y) d y
$$

By Claim B this expression for $\mu f(x)$ may take on at most 3 values. So the desired claim holds on $P \cap C$. Repeating the same argument, Claim $\square$ also holds on $N \cap C$.

We are nearly done with the proof of the theorem, as we have now reduced $W$ to a $10 \times 10$ stepgraphon. To complete the proof, we show that we may reduce to at most $7 \times 7$. We now partition $P \cap C, P \cap S, N \cap C$, and $N \cap S$ so that $f$ and $g$ are constant a.e. on each part as:

- $P \cap C=U_{1} \cup U_{2} \cup U_{3}$,
- $P \cap S=U_{4} \cup U_{5}$,
- $N \cap C=U_{6} \cup U_{7} \cup U_{8}$, and
- $N \cap S=U_{9} \cup U_{10}$.

Then by Lemma 3.6(ii), there exists a matrix $\left(m_{i j}\right)_{i, j \in[10]}$ so that for all $(i, j) \in[10] \times[10]$, - $m_{i j} \in\{0,1\}$,

- $W(x, y)=m_{i j}$ for a.e. $(x, y) \in U_{i} \times U_{j}$,
- $m_{i j}=1$ if and only if $f_{i} f_{j}>g_{i} g_{j}$, and
- $m_{i j}=0$ if and only if $f_{i} f_{j}<g_{i} g_{j}$.

Additionally, we set $\alpha_{i}=m\left(U_{i}\right)$ and also denote by $f_{i}$ and $g_{i}$ the constant values of $f, g$ on each $U_{i}$, respectively, for each $i=1, \ldots, 10$. Furthermore, by Claim C and Lemma 3.6 we assume without loss of generality that $f_{1}>f_{2}>f_{3} \geq 1>f_{4}>f_{5}$ and that $f_{6}>f_{7}>f_{8} \geq 1>f_{9}>f_{10}$. Also by Lemma 3.7, $0 \leq g_{1}<g_{2}<g_{3} \leq 1<g_{4}<g_{5}$ and $0 \leq-g_{1}<-g_{2}<-g_{3} \leq 1<-g_{4}<-g_{5}$. Also, by Claim B , no two columns of $m$ are identical within the sets $\{1,2,3,4,5\}$ and within $\{6,7,8,9,10\}$. Shading $m_{i j}=1$ black and $m_{i j}=0$ white, we let


Therefore, $W$ is a stepgraphon with values determined by $M$ and the size of each block determined by the $\alpha_{i}$.

We claim that $0 \in\left\{\alpha_{3}, \alpha_{4}, \alpha_{5}\right\}$ and $0 \in\left\{\alpha_{8}, \alpha_{9}, \alpha_{10}\right\}$. For the first claim, assume to the contrary that all of $\alpha_{3}, \alpha_{4}, \alpha_{5}$ are positive and note that there exists some $x_{4} \in U_{4}$ such that

$$
\mu f_{4}=\mu f\left(x_{4}\right)=\int_{y \in[0,1]} W\left(x_{4}, y\right) f(y) d y .
$$

Moreover for some measurable subsets $U_{3}^{\prime} \subseteq U_{3}$ and $U_{5}^{\prime} \subseteq U_{5}$ of positive measure so that with $U:=U_{3}^{\prime} \cup U_{4} \cup U_{5}^{\prime}$,

$$
f\left(x_{4}\right)^{2}=\frac{1}{m(U)} \int_{y \in U} f(y)^{2} d y .
$$

Note that by Lemma 3.7, we may assume that $x_{4}$ is average on $U$ with respect to $g$ as well. The conditions of Corollary 4.3 are met for $W$ with $A=U_{3}^{\prime}, B=U_{4} \cup U_{5}^{\prime}, x_{0}=x_{4}$. Since $\int_{A \times A} W(x, y) f(x) f(y) d x d y>0$, this is a contradiction to Corollary 4.3, so the desired claim holds. The same argument may be used to prove that $0 \in\left\{\alpha_{8}, \alpha_{9}, \alpha_{10}\right\}$.

We now form the principal submatrix $M^{\prime}$ by removing the $i$-th row and column from $M$ if and only if $\alpha_{i}=0$. Since $\alpha_{i}=0, W$ is a stepgraphon with values determined by $M^{\prime}$. Let $M_{P}^{\prime}$ denote the principal submatrix of $M^{\prime}$ corresponding to the indices $i \in\{1, \ldots, 5\}$ so that $\alpha_{i}>0$. That is, $M_{P}^{\prime}$ corresponds to the upper left hand block of $M$. We use red to indicate rows and columns present in $M$ but not $M_{P}^{\prime}$. When forming the submatrix
$M_{P}^{\prime}$, we borrow the internal subdivisions which are present in the definition of $M$ above to denote where $f \geq 1$ and where $f<1$ (or between $S \cap P$ and $C \cap P$ ). Note that this is not the same as what the internal divisions denote in the statement of the theorem. Since $0 \in\left\{\alpha_{3}, \alpha_{4}, \alpha_{5}\right\}$, it follows that $M_{P}^{\prime}$ is a principal submatrix of

, or


In the second case, columns 2 and 3 are identical in $M^{\prime}$, and in the third case, columns 1 and 2 are identical in $M^{\prime}$. So without loss of generality, $M_{P}^{\prime}$ is a principal submatrix of one of


In each case, $M_{P}^{\prime}$ is a principal submatrix of


An identical argument shows that the principal submatrix of $M^{\prime}$ on the indices $i \in$ $\{6, \ldots, 10\}$ such that $\alpha_{i}>0$ is a principal submatrix of


Finally, we note that $0 \in\left\{\alpha_{1}, \alpha_{6}\right\}$. Indeed otherwise the corresponding columns are identical in $M^{\prime}$, a contradiction. So without loss of generality, row and column 6 were also removed from $M$ to form $M^{\prime}$. This completes the proof of the theorem.

## 5. SPREAD MAXIMUM GRAPHONS

In this section, we complete the proof of the graphon version of the spread conjecture of Gregory, Hershkowitz, and Kirkland from [13]. In particular, we prove Theorem 5.1. For convenience and completeness, we state this result in the following level of detail.

Theorem 5.1. If $W$ is a graphon that maximizes spread, then $W$ may be represented as follows. For all $(x, y) \in[0,1]^{2}$,

$$
W(x, y)= \begin{cases}0, & \text { if }(x, y) \in[2 / 3,1]^{2} \\ 1, & \text { otherwise }\end{cases}
$$

## Furthermore,

$$
\mu=\frac{1+\sqrt{3}}{3} \quad \text { and } \quad \nu=\frac{1-\sqrt{3}}{3}
$$

are the maximum and minimum eigenvalues of $W$, respectively, and if $f, g$ are unit eigenfunctions associated to $\mu, \nu$, respectively, then, up to a change in sign, they may be written as follows. For every $x \in[0,1]$,

$$
\begin{aligned}
& f(x)=\frac{1}{2 \sqrt{3+\sqrt{3}}} \cdot\left\{\begin{array}{cl}
3+\sqrt{3}, & \text { if } x \in[0,2 / 3] \\
2 \cdot \sqrt{3}, & \text { otherwise }
\end{array}\right. \text {, and } \\
& g(x)=\frac{1}{2 \sqrt{3-\sqrt{3}}} \cdot\left\{\begin{array}{cl}
3-\sqrt{3}, & \text { if } x \in[0,2 / 3] \\
-2 \cdot \sqrt{3}, & \text { otherwise }
\end{array}\right.
\end{aligned}
$$

To help outline our proof of Theorem 5.1, let the spread-extremal graphon have block sizes $\alpha_{1}, \ldots, \alpha_{7}$. Note that the spread of the graphon is the same as the spread of matrix $M^{*}$ in Figure 33, and so we will optimize the spread of $M^{*}$ over choices of $\alpha_{1}, \ldots, \alpha_{7}$. Let $G^{*}$ be the unweighted graph (with loops) corresponding to the matrix.

We proceed in the following steps.
(1) In Section A.1, we reduce the proof of Theorem 5.1 to 17 cases, each corresponding to a subset $S$ of $V\left(G^{*}\right)$. For each such $S$ we define an optimization problem $\mathrm{SPR}_{S}$, the solution to which gives us an upper bound on the spread of any graphon in the case corresponding to $S$.
(2) In Section 5.2, we appeal to interval arithmetic to translate these optimization problems into algorithms. Based on the output of the 17 programs we wrote, we eliminate 15 of the 17 cases. We address the multitude of formulas used throughout and relocate their statements and proofs to Appendix B.1.
(3) Finally in Section 5.3, we complete the proof of Theorem 5.1 by analyzing the 2 remaining cases. Here, we apply Viète's Formula for roots of cubic equations and make a direct argument.


Figure 3. The matrix $M^{*}$ with corresponding graph $G^{*}$, where $D_{\alpha}$ is the diagonal matrix with entries $\alpha_{1}, \ldots, \alpha_{7}$

For concreteness, we define $G^{*}$ on the vertex set $\{1, \ldots, 7\}$. Explicitly, the neighborhoods $N_{1}, \ldots, N_{7}$ of $1, \ldots, 7$ are defined as:

$$
\begin{array}{ll}
N_{1}:=\{1,2,3,4,5,6,7\} & N_{2}:=\{1,2,3,5,6,7\} \\
N_{3}:=\{1,2,5,6,7\} & N_{4}:=\{1,5,6,7\} \\
N_{5}:=\{1,2,3,4,5,6\} & N_{6}:=\{1,2,3,4,5\} \\
N_{7}:=\{1,2,3,4\} &
\end{array}
$$

More compactly, we may note that

$$
\begin{array}{lll}
N_{1}=\{1, \ldots, 7\} & N_{2}=N_{1} \backslash\{4\} & N_{3}=N_{2} \backslash\{3\}
\end{array} N_{4}=N_{3} \backslash\{2\} .
$$

5.1. Stepgraphon case analysis. Let $W$ be a graphon maximizing spread. By Theorem 4.1, we may assume that $W$ is a $7 \times 7$ stepgraphon corresponding to $G^{*}$. We will break into cases depending on which of the 7 weights $\alpha_{1}, \ldots \alpha_{7}$ are zero and which are positive. For some of these combinations the corresponding graphons are isomorphic, and in this section we will outline how one can show that we need only consider 17 cases rather than $2^{7}$.

We will present each case with the set of indices which have strictly positive weight. Additionally, we will use vertical bars to partition the set of integers according to its intersection with the sets $\{1\},\{2,3,4\}$ and $\{5,6,7\}$. Recall that vertices in block 1 are dominating vertices and vertices in blocks 5,6 , and 7 have negative entries in the eigenfunction corresponding to $\nu$. For example, we use $4 \mid 57$ to refer to the case that $\alpha_{4}, \alpha_{5}, \alpha_{7}$ are all positive and $\alpha_{1}=\alpha_{2}=\alpha_{3}=\alpha_{6}=0$; see Figure 4.

To give an upper bound on the spread of any graphon corresponding to case $4 \mid 57$ we solve a constrained optimization problem. Let $f_{4}, f_{5}, f_{7}$ and $g_{4}, g_{5}, g_{7}$ denote the eigenfunction entries for unit eigenfunctions $f$ and $g$ of the graphon. Then we maximize $\mu-\nu$ subject to


FIgURE 4. The family of graphons and the graph corresponding to case 4|57

$$
\begin{aligned}
\alpha_{4}+\alpha_{5}+\alpha_{7} & =1 \\
\alpha_{4} f_{4}^{2}+\alpha_{5} f_{5}^{2}+\alpha_{7} f_{7}^{2} & =1 \\
\alpha_{4} g_{4}^{2}+\alpha_{5} g_{5}^{2}+\alpha_{7} g_{7}^{2} & =1 \\
\mu f_{i}^{2}-\nu g_{i}^{2} & =\mu-v \quad \text { for all } i \in\{4,5,7\} \\
\mu f_{4}=\alpha_{5} f_{5}+\alpha_{7} f_{7}, \quad \mu f_{5} & =\alpha_{4} f_{4}+\alpha_{5} f_{5}, \quad \mu f_{7}=\alpha_{4} f_{4} \\
\nu g_{4}=\alpha_{5} g_{5}+\alpha_{7} g_{7}, \quad \nu g_{5} & =\alpha_{4} g_{4}+\alpha_{5} g_{5}, \quad \nu g_{7}=\alpha_{4} g_{4}
\end{aligned}
$$

The first three constraints say that the weights sum to 1 and that $f$ and $g$ are unit eigenfunctions. The fourth constraint is from Lemma 3.7. The final two lines of constraints say that $f$ and $g$ are eigenfunctions for $\mu$ and $\nu$ respectively. Since these equations must be satisfied for any spread-extremal graphon, the solution to this optimization problem gives an upper bound on any spread-extremal graphon corresponding to case 4|57. For each case we formulate a similar optimization problem in Appendix A.1.

First, if two distinct blocks of vertices have the same neighborhood, then without loss of generality we may assume that only one of them has positive weight. For example, see Figure 5: in case $123 \mid 567$, blocks 1 and 2 have the same neighborhood, and hence without loss of generality we may assume that only block 1 has positive weight. Furthermore, in this case the resulting graphon could be considered as case $13 \mid 567$ or equivalently as case $14 \mid 567$; the graphons corresponding to these cases are isomorphic. Therefore cases $123|567,13| 567$, and $14 \mid 567$ reduce to considering only case $14 \mid 567$.

Additionally, if there is no dominant vertex, then some pairs cases may correspond to isomorphic graphons and the optimization problems are equivalent up to flipping the sign of the eigenvector corresponding to $\nu$. For example, see Figure 6, in which cases $23 \mid 457$ and $24 \mid 567$ reduce to considering only a single one. However, because of how we choose to order the eigenfunction entries when setting up the constraints of the optimization problems, there are some examples of cases corresponding to isomorphic graphons that we solve as separate optimization problems. For example, the graphons corresponding to cases $1|24| 7$ and $1|4| 57$ are isomorphic, but we will consider them separate cases; see Figure 7 .

Repeated applications of these three principles show that there are only 17 distinct cases that we must consider. The details are straightforward to verify, see Lemma A.2.


Figure 5. Redundancy, then renaming: we can assume $\alpha_{2}=0$ in the family of graphons corresponding to $123 \mid 567$, which produces families of graphons corresponding to both cases $13 \mid 567$ and $14 \mid 567$

The distinct cases that we must consider are the following, summarized in Figure 8 .

$$
\mathcal{S}_{17}:=\left\{\begin{array}{r}
1|234| 567,1|24| 567,1|234| 57,1|4| 567,1|24| 57,1|234| 7,234 \mid 567 \\
24|567,4| 567,24|57,1| 567,1|4| 57,1|2| 47,1|57,4| 57,1|4| 7,1 \mid 7
\end{array}\right\} .
$$

5.2. Interval arithmetic. Interval arithmetic is a computational technique which bounds errors that accumulate during computation. For convenience, let $\mathbb{R}^{*}$ := $[-\infty,+\infty]$ be the extended real line. To enhance order floating point arithmetic, we replace extended real numbers with unions of intervals which are guaranteed to contain them. Moreover, we extend the basic arithmetic operations,,$+- \times, \div$, and $\sqrt{ }$ to operations on unions of intervals. This technique has real-world applications in the hard sciences, but has also been used in computer-assisted proofs. For two famous examples, we refer the interested reader to [15] for Hales' proof of the Kepler Conjecture on optimal sphere-packing in $\mathbb{R}^{2}$, and to [31] for Warwick's solution of Smale's 14th problem on the Lorenz attractor as a strange attractor.

As stated before, we consider extensions of the binary operations,,$+- \times$, and $\div$ as well as the unary operation $\sqrt{ }$ defined on $\mathbb{R}$ to operations on unions of intervals of extended real numbers. For example if $[a, b],[c, d] \subseteq \mathbb{R}$, then we may use the following extensions of,+- , and $\times$ :

$$
\begin{aligned}
{[a, b]+[c, d] } & =[a+c, b+d], \\
{[a, b]-[c, d] } & =[a-d, b-c], \text { and } \\
{[a, b] \times[c, d] } & =[\min \{a c, a d, b c, b d\}, \max \{a c, a d, b c, b d\}] .
\end{aligned}
$$

For $\div$, we must address the cases $0 \in[c, d]$ and $0 \notin[c, d]$. Here, we take the extension

$$
[a, b] \div[c, d]=\left[\min \left\{\frac{a}{c}, \frac{a}{d}, \frac{b}{c}, \frac{b}{d}\right\}, \max \left\{\frac{a}{c}, \frac{a}{d}, \frac{b}{c}, \frac{b}{d}\right\}\right],
$$



Figure 6. Changing the sign of $g$ : The optimization problems in these cases are equivalent
where

$$
1 \div[c, d]=\left\{\begin{aligned}
{\left[\min \left\{c^{-1}, d^{-1}\right\}, \max \left\{c^{-1}, d^{-1}\right\}\right], } & 0 \notin[c, d] \\
{\left[d^{-1},+\infty\right], } & c=0 \\
{\left[-\infty, c^{-1}\right], } & d=0 \\
{\left[-\infty, \frac{1}{c}\right] \cup\left[\frac{1}{d},+\infty\right], } & c<0<d
\end{aligned}\right.
$$

Additionally, we may let

$$
\sqrt{[a, b]}=\left\{\begin{array}{rl}
\emptyset, & b<0 \\
{[\sqrt{\max \{0, a\}}, \sqrt{b}],} & \text { otherwise }
\end{array} .\right.
$$

When endpoints of $[a, b]$ and $[c, d]$ include $-\infty$ or $+\infty$, the definitions above must be modified slightly in a natural way.

We use interval arithmetic to prove the strict upper bound $<2 / \sqrt{3}$ for the maximum graphon spread claimed in Theorem 5.1, for any solutions to 15 of the 17 constrained optimization problems $\mathrm{SPR}_{S}$ stated in Lemma A.2. The constraints in each $\mathrm{SPR}_{S}$ allow us to derive equations for the variables $\left(\alpha_{i}, f_{i}, g_{i}\right)_{i \in S}$ in terms of each other, and $\mu$ and $\nu$. For the reader's convenience, we relocate these formulas and their derivations to Appendix B.1. In the programs corresponding to each set $S \in \mathcal{S}_{17}$, we find two indices


Figure 7. The cases $1|24| 7$ and $1|4| 57$ correspond to the same family graphons but we consider the optimization problems separately, due to our prescribed ordering of the vertices
$i \in S \cap\{1,2,3,4\}$ and $j \in S \cap\{5,6,7\}$ such that for all $k \in S, \alpha_{k}, f_{k}$, and $g_{k}$ may be calculated, step-by-step, from $\alpha_{i}, \alpha_{j}, \mu$, and $\nu$. See Table 1 for each set $S \in \mathcal{S}_{17}$, organized by the chosen values of $i$ and $j$.

TABLE 1. The indices $i, j$ corresponding to the search space used to bound solutions to $\mathrm{SPR}_{S}$

|  | 1 | 2 | 3 | 4 |
| :--- | ---: | ---: | ---: | ---: |
| 5 | $1 \mid 57$ | $24 \mid 57$ | $1\|234\| 57$ | $4 \mid 57$ |
|  |  | $1\|24\| 57$ |  | $1\|4\| 57$ |
| 6 | $1 \mid 567$ | $24 \mid 567$ | $234 \mid 567$ | $4 \mid 567$ |
|  |  | $1\|24\| 567$ | $1\|234\| 567$ | $1\|4\| 57$ |
| 7 | $1 \mid 7$ | $1\|24\| 7$ | $1\|234\| 7$ | $1\|4\| 7$ |

In the program corresponding to a set $S \in \mathcal{S}_{17}$, we search a carefully chosen set $\Omega \subseteq[0,1]^{3} \times[-1,0]$ for values of $\left(\alpha_{i}, \alpha_{j}, \mu, \nu\right)$ which satisfy $\operatorname{SPR}_{S}$. We first divide $\Omega$ into a grid of "boxes". Starting at depth 0 , we test each box $B$ for feasibility by assuming that $\left(\alpha_{i}, \alpha_{j}, \mu, \nu\right) \in B$ and that $\mu-\nu \geq 2 / \sqrt{3}$. Next, we calculate $\alpha_{k}, f_{k}$, and $g_{k}$ for all $k \in S$ in interval arithmetic using the formulas from Appendix $B$. When the calculation detects that a constraint of $\mathrm{SPR}_{S}$ is not satisfied, e.g., by showing that some $\alpha_{k}, f_{k}$, or $g_{k}$ lies in an empty interval, or by constraining $\sum_{i \in S} \alpha_{i}$ to a union of intervals which


Figure 8. The set $S_{17}$, as a poset ordered by inclusion. Each element is a subset of $V\left(G^{*}\right)=\{1, \ldots, 7\}$, written without braces and commas. As noted in the proof of Lemma A.2, the sets $\{1\},\{2,3,4\}$, and $\{5,6,7\}$ have different behavior in the problems $\mathrm{SPR}_{S}$. For this reason, we use vertical bars to separate each $S \in \mathcal{S}_{17}$ according to the corresponding partition.
does not contain 1, then the box is deemed infeasible. Otherwise, the box is split into two boxes of equal dimensions, with the dimension of the cut alternating cyclically.

For each $S \in \mathcal{S}_{17}$, the program $\mathrm{SPR}_{S}$ has 3 norm constraints, $2|S|$ linear eigenvector constraints, $|S|$ elliptical constraints, $\binom{|S|}{2}$ inequality constraints, and $3|S|$ interval membership constraints. By using interval arithmetic, we have a computer-assisted proof of the following result.

Lemma 5.2. Suppose $S \in \mathcal{S}_{17} \backslash\{\{1,7\},\{4,5,7\}\}$. Then any solution to $\mathrm{SPR}_{S}$ attains $a$ value strictly less than $2 / \sqrt{3}$.

To better understand the role of interval arithmetic in our proof, consider Example 5.3.

Example 5.3. Suppose $\mu, \nu$, and $\left(\alpha_{i}, f_{i}, g_{i}\right)$ is a solution to $\operatorname{SPR}_{\{1, \ldots, 7\}}$. We show that $\left(\alpha_{3}, \mu, \nu\right) \notin[.7, .8] \times[.9,1] \times[-.2,-.1]$. By Proposition B.1, $\mathrm{g}_{3}^{2}=\frac{\nu\left(\alpha_{3}+2 \mu\right)}{\alpha_{3}(\mu+\nu)+2 \mu \nu}$. Using
interval arithmetic,

$$
\begin{aligned}
\nu\left(\alpha_{3}+2 \mu\right) & =[-.2,-.1] \times([.7, .8]+2 \times[.9,1]) \\
& =[-.2,-.1] \times[2.5,2.8]=[-.56,-.25], \text { and } \\
\alpha_{3}(\mu+\nu)+2 \mu \nu & =[.7, .8] \times([.9,1]+[-.2,-.1])+2 \times[.9,1] \times[-.2,-.1] \\
& =[.7, .8] \times[.7, .9]+[1.8,2] \times[-.2,-.1] \\
& =[.49, .72]+[-.4,-.18]=[.09, .54] .
\end{aligned}
$$

Thus

$$
g_{3}^{2}=\frac{\nu\left(\alpha_{3}+2 \mu\right)}{\alpha_{3}(\mu+\nu)+2 \mu \nu}=[-.56,-.25] \div[.09, .54]=[-6 . \overline{2},-.4 \overline{629}] .
$$

Since $g_{3}^{2} \geq 0$, we have a contradiction.
Example 5.3 illustrates a number of key elements. First, we note that through interval arithmetic, we are able to provably rule out the corresponding region. However, the resulting interval for the quantity $g_{3}^{2}$ is over fifty times bigger than any of the input intervals. This growth in the size of intervals is common, and so, in some regions, fairly small intervals for variables are needed to provably illustrate the absence of a solution. For this reason, using a computer to complete this procedure is ideal, as doing millions of calculations by hand would be untenable.

However, the use of a computer for interval arithmetic brings with it another issue. Computers have limited memory, and therefore cannot represent all numbers in $\mathbb{R}^{*}$. Instead, a computer can only store a finite subset of numbers, which we will denote by $F \subsetneq \mathbb{R}^{*}$. This set $F$ is not closed under the basic arithmetic operations, and so when some operation is performed and the resulting answer is not in $F$, some rounding procedure must be performed to choose an element of $F$ to approximate the exact answer. This issue is the cause of roundoff error in floating point arithmetic, and must be treated in order to use computer-based interval arithmetic as a proof.

PyInterval is one of many software packages designed to perform interval arithmetic in a manner which accounts for this crucial feature of floating point arithmetic. Given some $x \in \mathbb{R}^{*}$, let $f l_{-}(x)$ be the largest $y \in F$ satisfying $y \leq x$, and $f l_{+}(x)$ be the smallest $y \in F$ satisfying $y \geq x$. Then, in order to maintain a mathematically accurate system of interval arithmetic on a computer, once an operation is performed to form a union of intervals $\bigcup_{i=1}^{k}\left[a_{i}, b_{i}\right]$, the computer forms a union of intervals containing [ $f l_{-}\left(a_{i}\right), f l_{+}\left(b_{i}\right)$ ] for all $1 \leq i \leq k$. The programs which prove Lemma 5.2 can be found in [27].
5.3. Completing the proof of Theorem 5.1. Finally, we complete the second main result of this paper. We will need Lemma 5.4.

Lemma 5.4. If $\left(\alpha_{4}, \alpha_{5}, \alpha_{7}\right)$ is a solution to $\operatorname{SPR}_{\{4,5,7\}}$, then $\alpha_{7}=0$.
We delay the proof of Lemma 5.4 to Appendix $A$ because it is technical. We now proceed with the Proof of Theorem 5.1.

Proof of Theorem 5.1. Let $W$ be a graphon such that $\operatorname{spr}(W)=\max _{U \in \mathcal{W}} \operatorname{spr}(U)$. By Lemma 5.2 and Lemma 5.4, $W$ is a $2 \times 2$ stepgraphon. Let the weights of the parts be $\alpha_{1}$ and $1-\alpha_{1}$.

Thus, it suffices to demonstrate the uniqueness of the desired solution $\mu, \nu$, and $\left(\alpha_{i}, f_{i}, g_{i}\right)_{i \in\{1,7\}}$ to $\operatorname{SPR}_{\{1,7\}}$. Indeed, we first note that with

$$
N\left(\alpha_{1}\right):=\left[\begin{array}{cc}
\alpha_{1} & 1-\alpha_{1} \\
\alpha_{1} & 0
\end{array}\right]
$$

the quantities $\mu$ and $\nu$ are precisely the eigenvalues of the characteristic polynomial

$$
p(x)=x^{2}-\alpha_{1} x-\alpha_{1}\left(1-\alpha_{1}\right) .
$$

In particular,

$$
\mu=\frac{\alpha_{1}+\sqrt{\alpha_{1}\left(4-3 \alpha_{1}\right)}}{2}, \quad \nu=\frac{\alpha_{1}-\sqrt{\alpha_{1}\left(4-3 \alpha_{1}\right)}}{2},
$$

and

$$
\mu-v=\sqrt{\alpha_{1}\left(4-3 \alpha_{1}\right)} .
$$

Optimizing, it follows that $\left(\alpha_{1}, 1-\alpha_{1}\right)=(2 / 3,1 / 3)$. Calculating the eigenfunctions and normalizing them gives that $\mu, \nu$, and their respective eigenfunctions match those from the statement of Theorem 5.1.

## 6. FROM GRaphons to graphs

In this section, we show that Theorem 5.1 implies Conjecture for all $n$ sufficiently large; that is, the solution to the problem of maximizing the spread of a graphon implies the solution to the problem of maximizing the spread of a graph for sufficiently large $n$.

The outline for our argument is as follows. First, we define the spread-maximum graphon $W$ as in Theorem 5.1. Let $\left\{G_{n}\right\}$ be any sequence where each $G_{n}$ is a spreadmaximum graph on $n$ vertices and denote by $\left\{W_{n}\right\}$ the corresponding sequence of graphons. We show that, after applying measure-preserving transformations to each $W_{n}$, the extreme eigenvalues and eigenvectors of each $W_{n}$ converge suitably to those of $W$. It follows for $n$ sufficiently large that except for $o(n)$ vertices, $G_{n}$ is a join of a clique of $2 n / 3$ vertices and an independent set of $n / 3$ vertices (Lemma 6.2). Using results from Section 2, we precisely estimate the extreme eigenvector entries on this $o(n)$ set. Finally, Lemma 6.3 shows that the set of $o(n)$ exceptional vertices is actually empty, completing the proof.

Before proceeding with the proof, we state Corollary 6.1 of the Davis-Kahan theorem [10], stated for graphons.

Corollary 6.1. Suppose $W, W^{\prime}:[0,1]^{2} \rightarrow[0,1]$ are graphons. Let $\mu$ be an eigenvalue of $W$ with $f$ a corresponding unit eigenfunction. Let $\left\{h_{k}\right\}$ be an orthonormal eigenbasis for $W^{\prime}$ with corresponding eigenvalues $\left\{\mu_{k}^{\prime}\right\}$. Suppose that $\left|\mu_{k}^{\prime}-\mu\right|>\delta$ for all $k \neq 1$. Then

$$
\sqrt{1-\left\langle h_{1}, f\right\rangle^{2}} \leq \frac{\left\|A_{W^{\prime}-W} f\right\|_{2}}{\delta}
$$

Before proving Theorem 1.1, we prove the following approximate result. For all nonnegative integers $n_{1}, n_{2}, n_{3}$, let $G\left(n_{1}, n_{2}, n_{3}\right):=\left(K_{n_{1}} \dot{\cup} K_{n_{2}}^{c}\right) \vee K_{n_{3}}^{c}$.

Lemma 6.2. For all positive integers $n$, let $G_{n}$ denote a graph on $n$ vertices which maximizes spread. Then $G_{n}=G\left(n_{1}, n_{2}, n_{3}\right)$ for some nonnegative integers $n_{1}, n_{2}, n_{3}$ such that $n_{1}=(2 / 3+o(1)) n, n_{2}=o(n)$, and $n_{3}=(1 / 3+o(1)) n$.
Proof. Our argument outline is:
(1) show that the eigenvectors for the spread-extremal graphs resemble the eigenfunctions of the spread-extremal graphon in an $L_{2}$ sense
(2) show that with the exception of a small proportion of vertices, a spreadextremal graph is the join of a clique and an independent set
Let $\mathcal{P}:=[0,2 / 3]$ and $\mathcal{N}:=[0,1] \backslash \mathcal{P}$. By Theorem 5.1, the graphon $W$ which is the indicator function of the set $[0,1]^{2} \backslash \mathcal{N}^{2}$ maximizes spread. Denote by $\mu$ and $\nu$ its maximum and minimum eigenvalues, respectively. For every positive integer $n$, let $G_{n}$ denote a graph on $n$ vertices which maximizes spread, let $W_{n}$ be any stepgraphon corresponding to $G_{n}$, and let $\mu_{n}$ and $v_{n}$ denote the maximum and minimum eigenvalues of $W_{n}$, respectively. By Theorems 3.1 and 5.1, and compactness of $\hat{\mathcal{W}}$,

$$
\max \left\{\left|\mu-\mu_{n}\right|,\left|\nu-v_{n}\right|, \delta_{\square}\left(W, W_{n}\right)\right\} \rightarrow 0
$$

Moreover, we may apply measure-preserving transformations to each $W_{n}$ so that without loss of generality, $\left\|W-W_{n}\right\|_{\square} \rightarrow 0$. As in Theorem 5.1, let $f$ and $g$ be unit eigenfunctions that take values $f_{1}, f_{2}, g_{1}, g_{2}$. Furthermore, let $\varphi_{n}$ be a nonnegative unit $\mu_{n}$ eigenfunction for $W_{n}$ and let $\psi_{n}$ be a $\nu_{n}$-eigenfunction for $W_{n}$.

We show that without loss of generality, $\varphi_{n} \rightarrow f$ and $\psi_{n} \rightarrow g$ in the $L_{2}$ sense. Since $\mu$ is the only positive eigenvalue of $W$ and it has multiplicity 1 , taking $\delta:=\mu / 2$, Corollary 6.1 implies that

$$
\begin{aligned}
1-\left\langle f, \varphi_{n}\right\rangle^{2} & \leq \frac{4\left\|A_{W-W_{n}} f\right\|_{2}^{2}}{\mu^{2}} \\
& =\frac{4}{\mu^{2}} \cdot\left\langle A_{W-W_{n}} f, A_{W-W_{n}} f\right\rangle \\
& \leq \frac{4}{\mu^{2}} \cdot\left\|A_{W-W_{n}} f\right\|_{1} \cdot\left\|A_{W-W_{n}} f\right\|_{\infty} \\
& \leq \frac{4}{\mu^{2}} \cdot\left(\left\|A_{W-W_{n}}\right\|_{\infty \rightarrow 1}\|f\|_{\infty}\right) \cdot\|f\|_{\infty} \\
& \leq \frac{16\left\|W-W_{n}\right\|_{\square} \cdot\|f\|_{\infty}^{2}}{\mu^{2}},
\end{aligned}
$$

where the last inequality follows from Lemma 8.11 of [18]. Since $\|f\|_{\infty} \leq 1 / \mu$, this proves the first claim. The second claim follows by replacing $f$ with $g$, and $\mu$ with $|\nu|$.
Note. For the remainder of the proof, we will introduce quantities $\varepsilon_{i}>0$ in lieu of writing complicated expressions explicitly. When we introduce a new $\varepsilon_{i}$, we will remark that given $\varepsilon_{0}, \ldots, \varepsilon_{i-1}$ sufficiently small, $\varepsilon_{i}$ can be made sufficiently small enough to meet some other conditions.

Let $\varepsilon_{0}>0$ and for all $n \geq 1$, define

$$
\begin{aligned}
\mathcal{P}_{n} & :=\left\{x \in[0,1]:\left|\varphi_{n}(x)-f_{1}\right|<\varepsilon_{0} \text { and }\left|\psi_{n}(x)-g_{1}\right|<\varepsilon_{0}\right\}, \\
\mathcal{N}_{n} & :=\left\{x \in[0,1]:\left|\varphi_{n}(x)-f_{2}\right|<\varepsilon_{0} \text { and }\left|\psi_{n}(x)-g_{2}\right|<\varepsilon_{0}\right\}, \text { and } \\
\mathcal{E}_{n} & :=[0,1] \backslash\left(\mathcal{P}_{n} \cup \mathcal{N}_{n}\right) .
\end{aligned}
$$

Since

$$
\begin{aligned}
& \int_{\left|\varphi_{n}-f\right| \geq \varepsilon_{0}}\left|\varphi_{n}-f\right|^{2} d x \leq \int\left|\varphi_{n}-f\right|^{2} d x \rightarrow 0, \text { and } \\
& \int_{\left|\psi_{n}-g\right| \geq \varepsilon_{0}}\left|\psi_{n}-g\right|^{2} d x \leq \int\left|\psi_{n}-g\right|^{2} d x \rightarrow 0,
\end{aligned}
$$

it follows that

$$
\max \left\{m\left(\mathcal{P}_{n} \backslash \mathcal{P}\right), m\left(\mathcal{N}_{n} \backslash \mathcal{N}\right), m\left(\mathcal{E}_{n}\right)\right\} \rightarrow 0
$$

For all $u \in V\left(G_{n}\right)$, let $S_{u}$ be the subinterval of $[0,1]$ corresponding to $u$ in $W_{n}$, and denote by $\varphi_{u}$ and $\psi_{u}$ the constant values of $\varphi_{n}$ on $S_{u}$. For convenience, we define the following discrete analogues of $\mathcal{P}_{n}, \mathcal{N}_{n}, \mathcal{E}_{n}$ :

$$
\begin{aligned}
P_{n} & :=\left\{u \in V\left(G_{n}\right):\left|\varphi_{u}-f_{1}\right|<\varepsilon_{0} \text { and }\left|\psi_{u}-g_{1}\right|<\varepsilon_{0}\right\}, \\
N_{n} & :=\left\{u \in V\left(G_{n}\right):\left|\varphi_{u}-f_{2}\right|<\varepsilon_{0} \text { and }\left|\psi_{u}-g_{2}\right|<\varepsilon_{0}\right\}, \text { and } \\
E_{n} & :=V\left(G_{n}\right) \backslash\left(P_{n} \cup N_{n}\right) .
\end{aligned}
$$

Let $\varepsilon_{1}>0$. By Lemma 2.4 and using the fact that $\mu_{n} \rightarrow \mu$ and $\nu_{n} \rightarrow \nu$,

$$
\begin{equation*}
\left|\mu \varphi_{u}^{2}-\nu \psi_{u}^{2}-(\mu-\nu)\right|<\varepsilon_{1} \quad \text { for all } u \in V\left(G_{n}\right) \tag{5}
\end{equation*}
$$

for all $n$ sufficiently large. Let $\varepsilon_{0}^{\prime}>0$. We next need Claim I. which says that the eigenvector entries of the exceptional vertices behave as if they have neighborhood $N_{n}$.
Claim I. Suppose $\varepsilon_{0}$ is sufficiently small and $n$ is sufficiently large in terms of $\varepsilon_{0}^{\prime}$. Then for all $v \in E_{n}$,

$$
\begin{equation*}
\max \left\{\left|\varphi_{v}-\frac{f_{2}}{3 \mu}\right|,\left|\psi_{v}-\frac{g_{2}}{3 \nu}\right|\right\}<\varepsilon_{0}^{\prime} \tag{6}
\end{equation*}
$$

Indeed, suppose $v \in E_{n}$ and let

$$
U_{n}:=\left\{w \in V\left(G_{n}\right): v w \in E\left(G_{n}\right)\right\} \quad \text { and } \quad u_{n}:=\bigcup_{w \in U_{n}} S_{w}
$$

We take two cases, depending on the sign of $\psi_{v}$.
Case A. $\psi_{v} \geq 0$.
Recall that $f_{2}>0>g_{2}$. Furthermore, $\varphi_{v} \geq 0$ and by assumption, $\psi_{v} \geq 0$. It follows that for all $n$ sufficiently large, $f_{2} \varphi_{v}-g_{2} \psi_{v}>0$, so by Lemma 2.1, $N_{n} \subseteq U_{n}$. Since $\varphi_{n}$ is a $\mu_{n}$-eigenfunction for $W_{n}$,

$$
\begin{aligned}
\mu_{n} \varphi_{v} & =\int_{y \in[0,1]} W_{n}(x, y) \varphi_{n}(y) d y \\
& =\int_{y \in \mathcal{P}_{n} \cap u_{n}} \varphi_{n}(y) d y+\int_{y \in \mathcal{N}_{n}} \varphi_{n}(y) d y+\int_{y \in \varepsilon_{n} \cap u_{n}} \varphi_{n}(y) d y .
\end{aligned}
$$

Similarly,

$$
\begin{aligned}
v_{n} \psi_{v} & =\int_{y \in[0,1]} K_{n}(x, y) \psi_{n}(y) d y \\
& =\int_{y \in \mathcal{P}_{n} \cap u_{n}} \psi_{n}(y) d y+\int_{y \in \mathcal{N}_{n}} \psi_{n}(y) d y+\int_{y \in \varepsilon_{n} \cap u_{n}} \psi_{n}(y) d y .
\end{aligned}
$$

Let $\rho_{n}:=m\left(\mathcal{P}_{n} \cap \mathcal{U}_{n}\right)$. Note that for all $\varepsilon_{2}>0$, as long as $n$ is sufficiently large and $\varepsilon_{1}$ is sufficiently small, then

$$
\begin{equation*}
\max \left\{\left|\varphi_{v}-\frac{3 \rho_{n} f_{1}+f_{2}}{3 \mu}\right|,\left|\psi_{v}-\frac{3 \rho_{n} g_{1}+g_{2}}{3 v}\right|\right\}<\varepsilon_{2} \tag{7}
\end{equation*}
$$

Let $\varepsilon_{3}>0$. By Equations (5) and (7) and with $\varepsilon_{1}, \varepsilon_{2}$ sufficiently small,

$$
\left|\mu \cdot\left(\frac{3 \rho_{n} f_{1}+f_{2}}{3 \mu}\right)^{2}-v \cdot\left(\frac{3 \rho_{n} g_{1}+g_{2}}{3 v}\right)^{2}-(\mu-v)\right|<\varepsilon_{3} .
$$

Substituting the values of $f_{1}, f_{2}, g_{1}, g_{2}$ from Theorem 5.1 and simplifying, it follows that

$$
\left|\frac{\sqrt{3}}{2} \cdot \rho_{n}\left(3 \rho_{n}-2\right)\right|<\varepsilon_{3} .
$$

Let $\varepsilon_{4}>0$. It follows that if $n$ is sufficiently large and $\varepsilon_{3}$ is sufficiently small, then

$$
\begin{equation*}
\min \left\{\rho_{n},\left|2 / 3-\rho_{n}\right|\right\}<\varepsilon_{4} . \tag{8}
\end{equation*}
$$

Combining Equations (7) and (8), it follows that with $\varepsilon_{2}, \varepsilon_{4}$ sufficiently small, then

$$
\begin{aligned}
& \max \left\{\left|\varphi_{v}-\frac{f_{2}}{3 \mu}\right|,\left|\psi_{v}-\frac{g_{2}}{3 \mu}\right|\right\}<\varepsilon_{0}^{\prime} \text {, or } \\
& \max \left\{\left|\varphi_{v}-\frac{2 f_{1}+f_{2}}{3 \mu}\right|,\left|\psi_{v}-\frac{2 g_{1}+g_{2}}{3 \mu}\right|\right\}<\varepsilon_{0}^{\prime} .
\end{aligned}
$$

Note that

$$
f_{1}=\frac{2 f_{1}+f_{2}}{3 \mu} \quad \text { and } \quad g_{1}=\frac{2 g_{1}+g_{2}}{3 v}
$$

Since $v \in E_{n}$, the second inequality does not hold, which completes the proof of the desired claim.
Case B. $\psi_{v}<0$.
Recall that $f_{1}>g_{1}>0$. Furthermore, $\varphi_{v} \geq 0$ and by assumption, $\psi_{v}<0$. It follows that for all $n$ sufficiently large, $f_{1} \varphi_{v}-g_{1} \psi_{v}>0$, so by Lemma 2.1, $P_{n} \subseteq U_{n}$. Since $\varphi_{n}$ is a $\mu_{n}$-eigenfunction for $W_{n}$,

$$
\begin{aligned}
\mu_{n} \varphi_{v} & =\int_{y \in[0,1]} W_{n}(x, y) \varphi_{n}(y) d y \\
& =\int_{y \in \mathcal{N}_{n} \cap u_{n}} \varphi_{n}(y) d y+\int_{y \in \mathcal{P}_{n}} \varphi_{n}(y) d y+\int_{y \in \varepsilon_{n} \cap u_{n}} \varphi_{n}(y) d y .
\end{aligned}
$$

Similarly,

$$
\begin{aligned}
v_{n} \psi_{v} & =\int_{y \in[0,1]} W_{n}(x, y) \psi_{n}(y) d y \\
& =\int_{y \in \mathcal{N}_{n} \cap u_{n}} \psi_{n}(y) d y+\int_{y \in \mathcal{P}_{n}} \psi_{n}(y) d y+\int_{y \in \varepsilon_{n} \cap u_{n}} \psi_{n}(y) d y .
\end{aligned}
$$

Let $\rho_{n}:=m\left(\mathcal{N}_{n} \cap \mathcal{U}_{n}\right)$. Note that for all $\varepsilon_{5}>0$, as long as $n$ is sufficiently large and $\varepsilon_{1}$ is sufficiently small, then

$$
\begin{equation*}
\max \left\{\left|\varphi_{v}-\frac{2 f_{1}+3 \rho_{n} f_{2}}{3 \mu}\right|,\left|\psi_{v}-\frac{2 g_{1}+3 \rho_{n} g_{2}}{3 v}\right|\right\}<\varepsilon_{5} \tag{9}
\end{equation*}
$$

Let $\varepsilon_{6}>0$. By Equations (55) and (9) and with $\varepsilon_{1}, \varepsilon_{2}$ sufficiently small,

$$
\left|\mu \cdot\left(\frac{2 f_{1}+3 \rho_{n} f_{2}}{3 \mu}\right)^{2}-v \cdot\left(\frac{2 g_{1}+3 \rho_{n} g_{2}}{3 v}\right)^{2}-(\mu-v)\right|<\varepsilon_{6} .
$$

Substituting the values of $f_{1}, f_{2}, g_{1}, g_{2}$ from Theorem 5.1 and simplifying, it follows that

$$
\left|2 \sqrt{3} \cdot \rho_{n}\left(3 \rho_{n}-1\right)\right|<\varepsilon_{6} .
$$

Let $\varepsilon_{7}>0$. It follows that if $n$ is sufficiently large and $\varepsilon_{6}$ is sufficiently small, then

$$
\begin{equation*}
\min \left\{\rho_{n},\left|1 / 3-\rho_{n}\right|\right\}<\varepsilon_{7} . \tag{10}
\end{equation*}
$$

Combining Equations (7) and (10), it follows that with $\varepsilon_{2}, \varepsilon_{4}$ sufficiently small, then

$$
\begin{aligned}
\max \left\{\left|\varphi_{v}-\frac{2 f_{1}}{3 \mu}\right|,\left|\psi_{v}-\frac{2 g_{1}}{3 \mu}\right|\right\} & <\varepsilon_{0}^{\prime}, \text { or } \\
\max \left\{\left|\varphi_{v}-\frac{2 f_{1}+f_{2}}{3 \mu}\right|,\left|\psi_{v}-\frac{2 g_{1}+g_{2}}{3 \mu}\right|\right\} & <\varepsilon_{0}^{\prime} .
\end{aligned}
$$

Again, note that

$$
f_{1}=\frac{2 f_{1}+f_{2}}{3 \mu} \quad \text { and } \quad g_{1}=\frac{2 g_{1}+g_{2}}{3 \nu}
$$

Since $v \in E_{n}$, the second inequality does not hold.
Similarly, note that

$$
f_{2}=\frac{2 f_{1}}{3 \mu} \quad \text { and } \quad g_{2}=\frac{2 g_{1}}{3 \nu}
$$

Since $v \in E_{n}$, the first inequality does not hold, a contradiction. So the desired claim holds.

We now complete the proof of Lemma 6.3 by showing that for all $n$ sufficiently large, $G_{n}$ is the join of an independent set $N_{n}$ with a disjoint union of a clique $P_{n}$ and an independent set $E_{n}$.

As above, we let $\varepsilon_{0}, \varepsilon_{0}^{\prime}>0$ be arbitrary. By definition of $P_{n}$ and $N_{n}$ and by Equation (6) from Claim [1], then for all $n$ sufficiently large,

$$
\begin{aligned}
\max \left\{\left|\varphi_{v}-f_{1}\right|,\left|\psi_{v}-g_{1}\right|\right\}<\varepsilon_{0} & \text { for all } v \in P_{n} \\
\max \left\{\left|\varphi_{v}-\frac{f_{2}}{3 \mu}\right|,\left|\psi_{v}-\frac{g_{2}}{3 v}\right|\right\}<\varepsilon_{0}^{\prime} & \text { for all } v \in E_{n} \\
\max \left\{\left|\varphi_{v}-f_{2}\right|,\left|\psi_{v}-g_{2}\right|\right\}<\varepsilon_{0} & \text { for all } v \in N_{n}
\end{aligned}
$$

With rows and columns respectively corresponding to the vertex sets $P_{n}, E_{n}$, and $N_{n}$, we note the following inequalities: Indeed, note the following inequalities:

| $f_{1}^{2}>g_{1}^{2}$ | $f_{1} \cdot \frac{f_{2}}{3 \mu}<g_{1} \cdot \frac{g_{2}}{3 \nu}$ | $f_{1} f_{2}>g_{1} g_{2}$ |
| :---: | :---: | :---: |
|  | $\left(\frac{f_{2}}{3 \mu}\right)^{2}<\left(\frac{g_{2}}{3 \nu}\right)^{2}$ | $\frac{f_{2}}{3 \mu} \cdot f_{2}>\frac{g_{2}}{3 \nu}$ |
|  |  | $f_{2}^{2}<g_{2}^{2}$ |

Let $\varepsilon_{0}, \varepsilon_{0}^{\prime}$ be sufficiently small. Then for all $n$ sufficiently large and for all $u, v \in V\left(G_{n}\right)$, then $\varphi_{u} \varphi_{v}-\psi_{u} \psi_{v}<0$ if and only if $u, v \in E_{n}, u, v \in N_{n}$, or $(u, v) \in\left(P_{n} \times E_{n}\right) \cup\left(E_{n} \times P_{n}\right)$. By Lemma 2.1, since $m\left(P_{n}\right) \rightarrow 2 / 3$ and $m\left(N_{n}\right) \rightarrow 1 / 3$, the proof is complete.

We have now shown that the spread-extremal graph is of the form $\left(K_{n_{1}} \dot{\cup} K_{n_{2}}^{c}\right) \vee K_{n_{3}}^{c}$ where $n_{2}=o(n)$. Lemma 6.3 refines this to show that actually $n_{2}=0$.

Lemma 6.3. For all nonnegative integers $n_{1}, n_{2}, n_{3}, \operatorname{let} G\left(n_{1}, n_{2}, n_{3}\right):=\left(K_{n_{1}} \cup K_{n_{2}}^{c}\right) \vee K_{n_{3}}^{c}$. Then for all $n$ sufficiently large, the following holds. If $\operatorname{spr}\left(G\left(n_{1}, n_{2}, n_{3}\right)\right)$ is maximized subject to the constraint $n_{1}+n_{2}+n_{3}=n$ and $n_{2}=o(n)$, then $n_{2}=0$.

Proof outline. We aim to maximize the spread of $G\left(n_{1}, n_{2}, n_{3}\right)$ subject to $n_{2}=o(n)$. The spread of $G\left(n_{1}, n_{2}, n_{3}\right)$ is the same as the spread of the quotient matrix

$$
Q_{n}=\left[\begin{array}{ccc}
n_{1}-1 & 0 & n_{3} \\
0 & 0 & n_{3} \\
n_{1} & n_{2} & 0
\end{array}\right]
$$

We reparametrize with parameters $\varepsilon_{1}$ and $\varepsilon_{2}$ representing how far away $n_{1}$ and $n_{3}$ are proportionally from $\frac{2 n}{3}$ and $\frac{n}{3}$, respectively. Namely, $\varepsilon_{1}=\frac{2}{3}-\frac{n_{1}}{n}$ and $\varepsilon_{2}=\frac{1}{3}-\frac{n_{3}}{n}$. Then $\varepsilon_{1}+\varepsilon_{2}=\frac{n_{2}}{n}$. Hence maximizing the spread of $G\left(n_{1}, n_{2}, n_{3}\right)$ subject to $n_{2}=o(n)$ is equivalent to maximizing the spread of the matrix

$$
n\left[\begin{array}{ccc}
\frac{2}{3}-\varepsilon_{1}-\frac{1}{n} & 0 & \frac{1}{3}-\varepsilon_{2} \\
0 & 0 & \frac{1}{3}-\varepsilon_{2} \\
\frac{2}{3}-\varepsilon_{1} & \varepsilon_{1}+\varepsilon_{2} & 0
\end{array}\right]
$$

subject to the constraint that $\frac{2}{3}-\varepsilon_{1}$ and $\frac{1}{3}-\varepsilon_{2}$ are nonnegative integer multiples of $\frac{1}{n}$ and $\varepsilon_{1}+\varepsilon_{2}=o(1)$. In order to utilize calculus, we instead solve a continuous relaxation of the optimization problem.

As such, consider the following matrix.

$$
M_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right):=\left[\begin{array}{ccc}
\frac{2}{3}-\varepsilon_{1}-z & 0 & \frac{1}{3}-\varepsilon_{2} \\
0 & 0 & \frac{1}{3}-\varepsilon_{2} \\
\frac{2}{3}-\varepsilon_{1} & \varepsilon_{1}+\varepsilon_{2} & 0
\end{array}\right]
$$

Since $M_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right)$ is diagonalizable, we may let $S_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right)$ be the difference between the maximum and minimum eigenvalues of $M_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right)$. We consider the optimization problem $\mathcal{P}_{z, C}$ defined for all $z \in \mathbb{R}$ and all $C>0$ such that $|z|$ and $C$ are sufficiently small, by

$$
\left(\mathcal{P}_{z, C}\right):\left\{\begin{aligned}
\max & S_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right) \\
\text { s.t. } & \varepsilon_{1}, \varepsilon_{2} \in[-C, C]
\end{aligned}\right.
$$

We show that as long as $C$ and $|z|$ are sufficiently small, then the optimum of $\mathcal{P}_{z, C}$ is attained by

$$
\left(\varepsilon_{1}, \varepsilon_{2}\right)=\left((1+o(z)) \cdot \frac{7 z}{30},(1+o(z)) \cdot \frac{-z}{3}\right)
$$

Moreover we show that in the feasible region of $\mathcal{P}_{z, C}, S_{z, C}\left(\varepsilon_{1}, \varepsilon_{2}\right)$ is concave-down in $\left(\varepsilon_{1}, \varepsilon_{2}\right)$. We return to the original problem by imposing the constraint that $\frac{2}{3}-\varepsilon_{1}$ and
$\frac{1}{3}-\varepsilon_{2}$ are multiples of $\frac{1}{n}$. Together these two observations complete the proof of the lemma. Under these added constraints, the optimum is obtained when

$$
\left(\varepsilon_{1}, \varepsilon_{2}\right)=\left\{\begin{array}{rll}
(0,0), & n \equiv 0 & (\bmod 3) \\
(2 / 3,-2 / 3), & n \equiv 1 & (\bmod 3) \\
(1 / 3,-1 / 3), & n \equiv 2 & (\bmod 3)
\end{array} .\right.
$$

Since the details are straightforward but tedious calculus, we delay this part of the proof to Section A.3.

We may now complete the proof of Theorem 1.1.
Proof of Theorem 1.1. Suppose $G$ is a graph on $n$ vertices which maximizes spread. By Lemma 6.2, $G=\left(K_{n_{1}} \dot{\cup} K_{n_{2}}^{c}\right) \vee K_{n_{3}}^{c}$ for some nonnegative integers $n_{1}, n_{2}, n_{3}$ such that $n_{1}+n_{2}+n_{3}=n$ where

$$
\left(n_{1}, n_{2}, n_{3}\right)=\left(\left(\frac{2}{3}+o(1)\right), o(n),\left(\frac{1}{3}+o(1)\right) \cdot n\right)
$$

By Lemma 6.3, if $n$ is sufficiently large, then $n_{2}=0$. To complete the proof of the main result, it is sufficient to find the unique maximum of $\operatorname{spr}\left(K_{n_{1}} \vee K_{n_{2}}^{c}\right)$, subject to the constraint that $n_{1}+n_{2}=n$. This is determined in [13] to be the join of a clique on $\left\lfloor\frac{2 n}{3}\right\rfloor$ and an independent set on $\left\lceil\frac{n}{3}\right\rceil$ vertices. The interested reader can prove that $n_{1}$ is the nearest integer to $(2 n-1) / 3$ by considering the spread of the quotient matrix

$$
\left[\begin{array}{cc}
n_{1}-1 & n_{2} \\
n_{1} & 0
\end{array}\right]
$$

and optimizing the choice of $n_{1}$.

## 7. The Bipartite Spread Conjecture

In [13], the authors investigated the structure of graphs that maximize spread over all graphs with a fixed number of vertices $n$ and edges $m$. In particular, they proved the upper bound

$$
\begin{equation*}
s(G) \leq \lambda_{1}+\sqrt{2 m-\lambda_{1}^{2}} \leq 2 \sqrt{m} \tag{11}
\end{equation*}
$$

and noted that equality holds throughout if and only if $G$ is the union of isolated vertices and $K_{p, q}$, for some $p+q \leq n$ satisfying $m=p q$ [13, Theorem 1.5]. This led the authors to make Conjecture 2, namely, that if $G$ has $n$ vertices, $m \leq\left\lfloor n^{2} / 4\right\rfloor$ edges, and spread $s(n, m)$, then $G$ is bipartite [13, Conjecture 1.4]. Recall that $s(n, m)$ is the maximum spread over all graphs with $n$ vertices and $m$ edges, and $s_{b}(n, m), m \leq\left\lfloor n^{2} / 4\right\rfloor$, is the maximum spread over all bipartite graphs with $n$ vertices and $m$ edges. This conjecture can be equivalently stated as $s(n, m)=s_{b}(n, m)$ for all $n \in \mathbb{N}$ and $m \leq\left\lfloor n^{2} / 4\right\rfloor$. In this section, we prove an asymptotic form of this conjecture and provide an infinite family of counterexamples to the exact conjecture, which verifies that the error in our asymptotic result is of the correct order of magnitude (Theorem 1.2).

To compute the spread of certain graphs explicitly, we make use of the theory of equitable partitions. In particular, we note that if $\phi$ is an automorphism of $G$, then the quotient matrix of $A(G)$ with respect to $\phi$, denoted by $A_{\phi}$, satisfies $\Lambda\left(A_{\phi}\right) \subset \Lambda(A)$, and
therefore $s(G)$ is at least the spread of $A_{\phi}$ (for details, see [9, Section 2.3]). Additionally, we require two propositions, one regarding the largest spectral radius of subgraphs of $K_{p, q}$ of a given size, and another regarding the largest gap between sizes that correspond to a complete bipartite graph of order at most $n$.

Let $K_{p, q}^{m}, 0 \leq p q-m<\min \{p, q\}$, be the subgraph of $K_{p, q}$ resulting from removing $p q-m$ edges all incident to some vertex in the larger side of the bipartition (if $p=q$, the vertex can be from either set). In [17], the authors proved the following result.
Proposition 7.1. If $0 \leq p q-m<\min \{p, q\}$, then $K_{p, q}^{m}$ maximizes $\lambda_{1}$ over all subgraphs of $K_{p, q}$ of size $m$.

We also require estimates regarding the longest sequence of consecutive sizes $m<$ $\left\lfloor n^{2} / 4\right\rfloor$ for which there does not exist a complete bipartite graph on at most $n$ vertices and exactly $m$ edges. As pointed out by [4], the result follows quickly by induction. However, for completeness, we include a brief proof.
Proposition 7.2. The length of the longest sequence of consecutive sizes $m<\left\lfloor n^{2} / 4\right\rfloor$ for which there does not exist a complete bipartite graph on at most $n$ vertices and exactly $m$ edges is zero for $n \leq 4$ and at most $\sqrt{2 n-1}-1$ for $n \geq 5$.
Proof. We proceed by induction. By inspection, for every $n \leq 4$, $m \leq\left\lfloor n^{2} / 4\right\rfloor$, there exists a complete bipartite graph of size $m$ and order at most $n$, and so the length of the longest sequence is trivially zero for $n \leq 4$. When $n=m=5$, there is no complete bipartite graph of order at most five with exactly five edges. This is the only such instance for $n=5$, and so the length of the longest sequence for $n=5$ is one.

Now, suppose that the statement holds for graphs of order at most $n-1$, for some $n>5$. We aim to show the statement for graphs of order at most $n$. By our inductive hypothesis, it suffices to consider only sizes $m \geq\left\lfloor(n-1)^{2} / 4\right\rfloor$ and complete bipartite graphs on $n$ vertices. We have

$$
\left(\frac{n}{2}+k\right)\left(\frac{n}{2}-k\right) \geq \frac{(n-1)^{2}}{4} \quad \text { for } \quad|k| \leq \frac{\sqrt{2 n-1}}{2}
$$

When $1 \leq k \leq \sqrt{2 n-1} / 2$, the difference between the sizes of $K_{n / 2+k-1, n / 2-k+1}$ and $K_{n / 2+k, n / 2-k}$ is at most

$$
\left|E\left(K_{\frac{n}{2}+k-1, \frac{n}{2}-k+1}\right)\right|-\left|E\left(K_{\frac{n}{2}+k, \frac{n}{2}-k}\right)\right|=2 k-1 \leq \sqrt{2 n-1}-1 .
$$

Let $k^{*}$ be the largest value of $k$ satisfying $k \leq \sqrt{2 n-1} / 2$ and $n / 2+k \in \mathbb{N}$. Then

$$
\begin{aligned}
\left|E\left(K_{\frac{n}{2}+k^{*}, \frac{n}{2}-k^{*}}\right)\right| & <\left(\frac{n}{2}+\frac{\sqrt{2 n-1}}{2}-1\right)\left(\frac{n}{2}-\frac{\sqrt{2 n-1}}{2}+1\right) \\
& =\sqrt{2 n-1}+\frac{(n-1)^{2}}{4}-1
\end{aligned}
$$

and the difference between the sizes of $K_{n / 2+k^{*}, n / 2-k^{*}}$ and $K_{\left\lceil\frac{n-1}{2}\right\rceil,\left\lfloor\frac{n-1}{2}\right\rfloor}$ is at most

$$
\begin{aligned}
\left|E\left(K_{\frac{n}{2}+k^{*}, \frac{n}{2}-k^{*}}\right)\right|-\left|E\left(K_{\left\lceil\frac{n-1}{2} l,\left[\frac{n-1}{2}\right\rfloor\right.}\right)\right| & <\sqrt{2 n-1}+\frac{(n-1)^{2}}{4}-\left\lfloor\frac{(n-1)^{2}}{4}\right\rfloor-1 \\
& <\sqrt{2 n-1} .
\end{aligned}
$$

Combining these two estimates completes our inductive step, and the proof.

We are now prepared to prove an asymptotic version of [13, Conjecture 1.4], and provide an infinite class of counterexamples that illustrates that the asymptotic version under consideration is the tightest version of this conjecture possible.

Theorem 7.3 (Restatement of Theorems 1.2 and 1.3). For all $n, m \in \mathbb{N}$ satisfying $m \leq$ $\left\lfloor n^{2} / 4\right\rfloor$,

$$
\begin{equation*}
s(n, m)-s_{b}(n, m) \leq \frac{1+16 m^{-3 / 4}}{m^{3 / 4}} s(n, m) . \tag{12}
\end{equation*}
$$

In addition, for any $\varepsilon>0$, there exists some $n_{\varepsilon}$ such that

$$
\begin{equation*}
s(n, m)-s_{b}(n, m) \geq \frac{1-\varepsilon}{m^{3 / 4}} s(n, m) \tag{13}
\end{equation*}
$$

for all $n \geq n_{\varepsilon}$ and some $m \leq\left\lfloor n^{2} / 4\right\rfloor$ depending on $n$.
Proof. The main idea of the proof is as follows. To obtain Inequality (12) (the upper bound on $s(n, m)-s_{b}(n, m)$, we upper bound $s(n, m)$ by $2 \sqrt{m}$ using Inequality (11), and we lower bound $s_{b}(n, m)$ by the spread of some specific bipartite graph. To obtain Inequality (13) (the lower bound on $s(n, m)-s_{b}(n, m)$ ) for a specific $n$ and $m$, we explicitly compute $s_{b}(n, m)$ using Proposition 7.1, and lower bound $s(n, m)$ by the spread of some specific non-bipartite graph.

First, we analyze the spread of $K_{p, q}^{m}, 0<p q-m<q \leq p$, a quantity that will be used in the proofs of both Inequalities (12) and (13). Let us denote the vertices in the bipartition of $K_{p, q}^{m}$ by $u_{1}, \ldots, u_{p}$ and $v_{1}, \ldots, v_{q}$, and suppose without loss of generality that $u_{1}$ is not adjacent to $v_{1}, \ldots, v_{p q-m}$. Then

$$
\phi=\left(u_{1}\right)\left(u_{2}, \ldots, u_{p}\right)\left(v_{1}, \ldots, v_{p q-m}\right)\left(v_{p q-m+1}, \ldots, v_{q}\right)
$$

is an automorphism of $K_{p, q}^{m}$. The corresponding quotient matrix is given by

$$
A_{\phi}=\left(\begin{array}{cccc}
0 & 0 & 0 & m-(p-1) q \\
0 & 0 & p q-m & m-(p-1) q \\
0 & p-1 & 0 & 0 \\
1 & p-1 & 0 & 0
\end{array}\right)
$$

and has characteristic polynomial

$$
Q(p, q, m)=\operatorname{det}\left[A_{\phi}-\lambda I\right]=\lambda^{4}-m \lambda^{2}+(p-1)(m-(p-1) q)(p q-m)
$$

and, therefore, by the Perron-Frobenius Theorem,

$$
\begin{equation*}
s\left(K_{p, q}^{m}\right)=2\left(\frac{m+\sqrt{m^{2}-4(p-1)(m-(p-1) q)(p q-m)}}{2}\right)^{1 / 2} \tag{14}
\end{equation*}
$$

We are now prepared to prove Inequality (12) (the upper bound). For some fixed $n$ and $m \leq\left\lfloor n^{2} / 4\right\rfloor$, let $m=p q-r$, where $p, q, r \in \mathbb{N}, p+q \leq n$, and $r$ is as small as possible. If $r=0$, then by [13, Theorem 1.5] (described above), $s(n, m)=s_{b}(n, m)$, and we are done. Otherwise, we note that $0<r<\min \{p, q\}$, and so Equation (14) is applicable
(in fact, by Proposition 7.2, $r=O(\sqrt{n})$ ). Using the upper bound $s(n, m) \leq 2 \sqrt{m}$ and Equation (14), we have

$$
\begin{equation*}
\frac{s(n, p q-r)-s\left(K_{p, q}^{m}\right)}{s(n, p q-r)} \leq 1-\left(\frac{1}{2}+\frac{1}{2} \sqrt{1-\frac{4(p-1)(q-r) r}{(p q-r)^{2}}}\right)^{1 / 2} \tag{15}
\end{equation*}
$$

To upper bound $r$, we use Proposition 7.2 with $n^{\prime}=\lceil 2 \sqrt{m}\rceil \leq n$ and $m$. This implies that

$$
r \leq \sqrt{2\lceil 2 \sqrt{m}\rceil-1}-1<\sqrt{2(2 \sqrt{m}+1)-1}-1=\sqrt{4 \sqrt{m}+1}-1 \leq 2 m^{1 / 4}
$$

Recall that $\sqrt{1-x} \geq 1-x / 2-x^{2} / 2$ for all $x \in[0,1]$, and so

$$
\begin{aligned}
1-\left(\frac{1}{2}+\frac{1}{2} \sqrt{1-x}\right)^{1 / 2} & \leq 1-\left(\frac{1}{2}+\frac{1}{2}\left(1-\frac{1}{2} x-\frac{1}{2} x^{2}\right)\right)^{1 / 2}=1-\left(1-\frac{1}{4}\left(x+x^{2}\right)\right)^{1 / 2} \\
& \leq 1-\left(1-\frac{1}{8}\left(x+x^{2}\right)-\frac{1}{32}\left(x+x^{2}\right)^{2}\right) \\
& \leq \frac{1}{8} x+\frac{1}{4} x^{2}
\end{aligned}
$$

for $x \in[0,1]$. To simplify Inequality (15), we observe that

$$
\frac{4(p-1)(q-r) r}{(p q-r)^{2}} \leq \frac{4 r}{m} \leq \frac{8}{m^{3 / 4}}
$$

Therefore,

$$
\frac{s(n, p q-r)-s\left(K_{p, q}^{m}\right)}{s(n, p q-r)} \leq \frac{1}{m^{3 / 4}}+\frac{16}{m^{3 / 2}}
$$

This completes the proof of Inequality (12) (the upper bound).
Finally, we proceed with the proof of Inequality (13) (the lower bound). Let us fix some $0<\varepsilon<1$, and consider some sufficiently large $n$. Let $m=(n / 2+k)(n / 2-k)+1$, where $k$ is the smallest number satisfying $n / 2+k \in \mathbb{N}$ and $\hat{\varepsilon}:=1-2 k^{2} / n<\varepsilon / 2$ (here we require $n=\Omega\left(1 / \varepsilon^{2}\right)$ ). Denote the vertices in the bipartition of $K_{n / 2+k, n / 2-k}$ by $u_{1}, \ldots, u_{n / 2+k}$ and $v_{1}, \ldots, v_{n / 2-k}$, and consider the graph $K_{n / 2+k, n / 2-k}^{+}:=K_{n / 2+k, n / 2-k} \cup$ $\left\{\left(v_{1}, v_{2}\right)\right\}$ resulting from adding one edge to $K_{n / 2+k, n / 2-k}$ between two vertices in the smaller side of the bipartition. Then

$$
\phi=\left(u_{1}, \ldots, u_{n / 2+k}\right)\left(v_{1}, v_{2}\right)\left(v_{3}, \ldots, v_{n / 2-k}\right)
$$

is an automorphism of $K_{n / 2+k, n / 2-k}^{+}$, and

$$
A_{\phi}=\left(\begin{array}{ccc}
0 & 2 & n / 2-k-2 \\
n / 2+k & 1 & 0 \\
n / 2+k & 0 & 0
\end{array}\right)
$$

has characteristic polynomial

$$
\begin{aligned}
\operatorname{det}\left[A_{\phi}-\lambda I\right] & =-\lambda^{3}+\lambda^{2}+\left(n^{2} / 4-k^{2}\right) \lambda-(n / 2+k)(n / 2-k-2) \\
& =-\lambda^{3}+\lambda^{2}+\left(\frac{n^{2}}{4}-\frac{(1-\hat{\varepsilon}) n}{2}\right) \lambda-\left(\frac{n^{2}}{4}-\frac{(3-\hat{\varepsilon}) n}{2}-\sqrt{2(1-\hat{\varepsilon}) n}\right) .
\end{aligned}
$$

By matching higher order terms, we obtain

$$
\lambda_{\max }\left(A_{\phi}\right)=\frac{n}{2}-\frac{1-\hat{\varepsilon}}{2}+\frac{\left(8-(1-\hat{\varepsilon})^{2}\right)}{4 n}+o(1 / n)
$$

$$
\lambda_{\min }\left(A_{\phi}\right)=-\frac{n}{2}+\frac{1-\hat{\varepsilon}}{2}+\frac{\left(8+(1-\hat{\varepsilon})^{2}\right)}{4 n}+o(1 / n)
$$

and

$$
s\left(K_{n / 2+k, n / 2-k}^{+}\right) \geq n-(1-\hat{\varepsilon})-\frac{(1-\hat{\varepsilon})^{2}}{2 n}+o(1 / n)
$$

Next, we aim to compute $s_{b}(n, m), m=(n / 2+k)(n / 2-k)+1$. By Proposition 7.1, $s_{b}(n, m)$ is equal to the maximum of $s\left(K_{n / 2+\ell, n / 2-\ell}^{m}\right)$ over all $\ell \in[0, k-1], k-\ell \in \mathbb{N}$. As previously noted, the quantity $s\left(K_{n / 2+\ell, n / 2-\ell}^{m}\right)$ is given exactly by Equation (14), and so the optimal choice of $\ell$ minimizes

$$
\begin{aligned}
f(\ell) & :=(n / 2+\ell-1)\left(k^{2}-\ell^{2}-1\right)\left(n / 2-\ell-\left(k^{2}-\ell^{2}-1\right)\right) \\
& =(n / 2+\ell)\left((1-\hat{\varepsilon}) n / 2-\ell^{2}\right)\left(\hat{\varepsilon} n / 2+\ell^{2}-\ell\right)+O\left(n^{2}\right) .
\end{aligned}
$$

We have

$$
f(k-1)=(n / 2+k-2)(2 k-2)(n / 2-3 k+3)
$$

and if $\ell \leq \frac{4}{5} k$, then $f(\ell)=\Omega\left(n^{3}\right)$. Therefore the minimizing $\ell$ is in $\left[\frac{4}{5} k, k\right]$. The derivative of $f(\ell)$ is given by

$$
\begin{aligned}
f^{\prime}(\ell)= & \left(k^{2}-\ell^{2}-1\right)\left(n / 2-\ell-k^{2}+\ell^{2}+1\right) \\
& -2 \ell(n / 2+\ell-1)\left(n / 2-\ell-k^{2}+\ell^{2}+1\right) \\
& +(2 \ell-1)(n / 2+\ell-1)\left(k^{2}-\ell^{2}-1\right) .
\end{aligned}
$$

For $\ell \in\left[\frac{4}{5} k, k\right]$,

$$
\begin{aligned}
f^{\prime}(\ell) & \leq \frac{n\left(k^{2}-\ell^{2}\right)}{2}-\ell n\left(n / 2-\ell-k^{2}+\ell^{2}\right)+2 \ell(n / 2+\ell)\left(k^{2}-\ell^{2}\right) \\
& \leq \frac{9 k^{2} n}{50}-\frac{4}{5} k n\left(n / 2-k-\frac{9}{25} k^{2}\right)+\frac{18}{25}(n / 2+k) k^{3} \\
& =\frac{81 k^{3} n}{125}-\frac{2 k n^{2}}{5}+O\left(n^{2}\right) \\
& =k n^{2}\left(\frac{81(1-\hat{\varepsilon})}{250}-\frac{2}{5}\right)+O\left(n^{2}\right)<0
\end{aligned}
$$

for sufficiently large $n$. This implies that the optimal choice is $\ell=k-1$, and $s_{b}(n, m)=$ $s\left(K_{n / 2+k-1, n / 2-k+1}^{m}\right)$. The characteristic polynomial $Q\left(n / 2+k-1, n / 2-k+1, n^{2} / 4-\right.$ $k^{2}+1$ ) equals

$$
\lambda^{4}-\left(n^{2} / 4-k^{2}+1\right) \lambda^{2}+2(n / 2+k-2)(n / 2-3 k+3)(k-1) .
$$

By matching higher order terms, the extreme root of $Q$ is given by

$$
\lambda=\frac{n}{2}-\frac{1-\hat{\varepsilon}}{2}-\sqrt{\frac{2(1-\hat{\varepsilon})}{n}}+\frac{27-14 \hat{\varepsilon}-\hat{\varepsilon}^{2}}{4 n}+o(1 / n),
$$

and so

$$
s_{b}(n, m)=n-(1-\hat{\varepsilon})-2 \sqrt{\frac{2(1-\hat{\varepsilon})}{n}}+\frac{27-14 \hat{\varepsilon}-\hat{\varepsilon}^{2}}{2 n}+o(1 / n),
$$



FIGURE 9. The graphs $K_{7,3}^{+}, K_{6,4}^{22}$, and $K_{5,5}^{22}$ from Example 7.4. The additional edge in $K_{7,3}^{+}$is between the two pink vertices, and the missing edges in $K_{6,4}^{22}$ and $K_{7,3}^{22}$ are between the black vertex and the grey vertices.
and

$$
\begin{aligned}
\frac{s(n, m)-s_{b}(n, m)}{s(n, m)} & \geq \frac{2^{3 / 2}(1-\hat{\varepsilon})^{1 / 2}}{n^{3 / 2}}-\frac{14-8 \hat{\varepsilon}}{n^{2}}+o\left(1 / n^{2}\right) \\
& =\frac{(1-\hat{\varepsilon})^{1 / 2}}{m^{3 / 4}}+\frac{(1-\hat{\varepsilon})^{1 / 2}}{(n / 2)^{3 / 2}}\left[1-\frac{(n / 2)^{3 / 2}}{m^{3 / 4}}\right]-\frac{14-8 \hat{\varepsilon}}{n^{2}}+o\left(1 / n^{2}\right) \\
& \geq \frac{1-\varepsilon / 2}{m^{3 / 4}}+o\left(1 / m^{3 / 4}\right)
\end{aligned}
$$

This completes the proof.
A concrete example where the Bipartite Spread Conjecture fails is as follows.
Example 7.4. Let $n=10$ and $m=22$, and consider the graphs $K_{7,3}^{+}, K_{6,4}^{22}$, and $K_{5,5}^{22}$ (see Figure 9 for a visual representation). We have

$$
\begin{aligned}
& s\left(K_{7,3}^{+}\right)=8 \cos \left(\frac{1}{3} \tan ^{-1}(3 \sqrt{29127})\right)+\frac{8}{\sqrt{3}} \sin \left(\frac{1}{3} \tan ^{-1}(3 \sqrt{29127})\right) \approx 9.238 \\
& s\left(K_{6,4}^{22}\right)=2(11+\sqrt{101})^{1 / 2} \approx 9.176, \text { and } \\
& s\left(K_{5,5}^{22}\right)=2(11+\sqrt{97})^{1 / 2} \approx 9.132
\end{aligned}
$$

Every bipartite graph with 10 vertices and 22 edges is either a subgraph of $K_{5,5}$ or $K_{6,4}$, and so, by Proposition 7.1, $s_{b}(10,22)=\max \left\{s\left(K_{6,4}^{22}\right), s\left(K_{5,5}^{22}\right)\right\}=s\left(K_{6,4}^{22}\right)$. The graph $K_{7,3}^{+}$ is a lower bound for $s(10,22)$, and so

$$
s(10,22)-s_{b}(10,22) \geq s\left(K_{7,3}^{+}\right)-s\left(K_{6,4}^{22}\right) \geq .06
$$

It is unclear whether this is the smallest counterexample (w.r.t. $n$ ) to the Bipartite Spread Conjecture, though we note that in [13] the authors verified the conjecture by computer for all graphs of order $n \leq 8$. Verifying that Example 7.4 is the smallest counterexample would require computations for $n=9$, a task that we leave to the motivated reader.

## 8. Concluding remarks

In this work we provided a proof of the spread conjecture for sufficiently large $n$, a proof of an asymptotic version of the bipartite spread conjecture, and an infinite class of counterexamples illustrating that our asymptotic version of this conjecture is the strongest result possible. There are a number of interesting future avenues of research, some of which we briefly describe below. These avenues consist primarily of considering the spread of more general classes of graphs (for instance, directed graphs, graphs with loops) or considering more general objective functions.

Our proof of the spread conjecture for sufficiently large $n$ immediately implies a nearly-tight estimate for the adjacency matrix of undirected graphs with loops, also commonly referred to as symmetric $0-1$ matrices. Given a directed graph $G=(V, \mathcal{A})$, the corresponding adjacency matrix $A$ has entry $A_{i, j}=1$ if the $\operatorname{arc}(i, j) \in \mathcal{A}$, and is zero otherwise. In this case, $A$ is not necessarily symmetric, and may have complex eigenvalues. One interesting question is what digraph of order $n$ maximizes the spread of its adjacency matrix, where spread is defined as the diameter of the spectrum. Is this more general problem also maximized by the same set of graphs as in the undirected case? This problem for either loop-less directed graphs or directed graphs with loops is an interesting question, and the latter is equivalent to asking the above question for the set of all $0-1$ matrices.

Another approach is to restrict ourselves to undirected graphs or undirected graphs with loops, and further consider the competing interests of simultaneously producing a graph with both $\lambda_{1}$ and $-\lambda_{n}$ large by understanding the trade-off between these two goals. To this end, we propose considering the class of objective functions

$$
f(G ; \beta)=\beta \lambda_{1}(G)-(1-\beta) \lambda_{n}(G), \quad \beta \in[0,1] .
$$

When $\beta=0$, this function is maximized by the complete bipartite graph $K_{[n / 2], \mid n / 2]}$ and when $\beta=1$, this function is maximized by the complete graph $K_{n}$. This paper treats the specific case of $\beta=1 / 2$, but none of the mathematical techniques used in this work rely on this restriction. In fact, the structural graph-theoretic results of Section 2, suitably modified for arbitrary $\beta$, still hold (see the thesis [32, Section 3.3.1] for this general case). Understanding the behavior of the optimum between these three wellstudied choices of $\beta=0,1 / 2,1$ is an interesting future avenue of research.

More generally, any linear combination of graph eigenvalues could be optimized over any family of graphs. Many sporadic examples of this problem have been studied. Nikiforov [23] proposed a general framework for it and proved some conditions under which the problem is well behaved. We conclude with some specific instances of the problem that we think deserve consideration.

Given a graph $F$, maximizing $\lambda_{1}$ over the family of $n$-vertex $F$-free graphs can be thought of as a spectral version of Turán's problem. Many papers have been written about this problem, which was proposed in generality in [25]. We remark that these results can often strengthen classical results in extremal graph theory. Maximizing
$\lambda_{1}+\lambda_{n}$ over the family of triangle-free graphs has been considered in [8] and is related to an old conjecture of Erdős on how many edges must be removed from a triangle-free graph to make it bipartite [12]. In general it would be interesting to maximize $\lambda_{1}+\lambda_{n}$ over the family of $K_{r}$-free graphs. When a graph is regular, the difference between $\lambda_{1}$ and $\lambda_{2}$ (the spectral gap) is related to the graph's expansion properties. Aldous and Fill [3] asked for the minimum of $\lambda_{1}-\lambda_{2}$ over the family of $n$-vertex connected regular graphs. Partial results were given by [1, 2, 7, 14]. A nonregular version of the problem was proposed by Stanić [28], who asked for the minimum of $\lambda_{1}-\lambda_{2}$ over connected $n$ vertex graphs. Finally, maximizing $\lambda_{3}$ or $\lambda_{4}$ over the family of $n$-vertex graphs seems to be a surprisingly difficult question and even the asymptotics are not known (see [26]).

## Appendix A. Technical proofs

A.1. Reduction to 17 cases. Now, we introduce the following specialized notation. For any nonempty set $S \subseteq V\left(G^{*}\right)$ and any labeled partition $\left(I_{i}\right)_{i \in S}$ of [0,1], we define the stepgraphon $W_{\mathcal{J}}$ as follows. For all $i, j \in S, W_{\mathcal{J}}$ equals 1 on $I_{i} \times I_{j}$ if and only if $i j$ is an edge (or loop) of $G^{*}$, and 0 otherwise. If $\alpha=\left(\alpha_{i}\right)_{i \in S}$ where $\alpha_{i}=m\left(I_{i}\right)$ for all $i \in S$, we may write $W_{\alpha}$ to denote the graphon $W_{\mathcal{J}}$ up to weak isomorphism.

To make the observations from Section 5.1 more explicit, we note that Theorem 4.1 implies that a spread-optimal graphon has the form $W=W_{\mathcal{J}}$ where $\mathcal{J}=\left(I_{i}\right)_{i \in S}$ is a labeled partition of $[0,1], S \subseteq[7]$, and each $I_{i}$ is measurable with positive measure. Since $W$ is a stepgraphon, its extreme eigenfunctions may be taken to be constant on $I_{i}$, for all $i \in S$. With $f, g$ denoting the extreme eigenfunctions for $W$, we may let $f_{i}$ and $g_{i}$ be the constant value of $f$ and $g$, respectively, on step $S_{i}$, for all $i \in S$. Appealing again to Theorem 4.1, we may assume without loss of generality that $f_{i} \geq 0$ for all $i \in S$, and for all $i \in S, g_{i} \geq 0$ implies that $i \in\{1,2,3,4\}$. By Lemma 3.7, for each $i \in S$, $\mu f_{i}^{2}-\nu g_{i}^{2}=\mu-\nu$. Combining these facts, we note that $f_{i}$ and $g_{i}$ belong to specific intervals as in Figure 10.


FIGURE 10. Intervals containing the quantities $f_{i}$ and $g_{i}$. Note that $f_{i}$ and $g_{i}$ are only defined for all $i \in S$.

For convenience, we define the following sets $F_{i}$ and $G_{i}$, for all $i \in S$. First, let $\mathcal{U}:=[0,1]$ and $\mathcal{V}:=[1,+\infty]$. With some abuse of notation, we denote $-\mathcal{U}=[-1,0]$ and $-\mathcal{V}=[-\infty,-1]$.

For each $i \in V\left(G^{*}\right)$, we define the intervals $F_{i}$ and $G_{i}$ by

$$
\left(F_{i}, G_{i}\right):=\left\{\begin{aligned}
(\mathcal{V}, \mathcal{U}), & i \in\{1,2\}, \\
(\mathcal{U}, \mathcal{V}), & i \in\{3,4\}, \\
(\mathcal{V},-\mathcal{U}), & i=5, \\
(\mathcal{U},-\mathcal{V}), & i \in\{6,7\} .
\end{aligned}\right.
$$

Given that the set $S$ and the quantities $\left(\alpha_{i}, f_{i}, g_{i}\right)_{i \in S}$ are clear from context, we label the following equation:

$$
\begin{equation*}
\sum_{i \in S} \alpha_{i}=\sum_{i \in S} \alpha_{i} f_{i}^{2}=\sum_{i \in S} \alpha_{i} g_{i}^{2}=1 . \tag{16}
\end{equation*}
$$

Furthermore when $i \in S$ is understood from context, we define the equations

$$
\begin{align*}
\mu f_{i}^{2}-v g_{i}^{2} & =\mu-v,  \tag{17}\\
\sum_{j \in N_{i} \cap S} \alpha_{j} f_{j} & =\mu f_{i}  \tag{18}\\
\sum_{j \in N_{i} \cap S} \alpha_{j} g_{j} & =v g_{i} . \tag{19}
\end{align*}
$$

Additionally, we consider the following inequalities. For all $S \subseteq V\left(G^{*}\right)$ and all distinct $i, j \in S$,

$$
f_{i} f_{j}-g_{i} g_{j}\left\{\begin{array}{l}
\geq 0, \quad i j \in E\left(G^{*}\right)  \tag{20}\\
\leq 0, \quad i j \notin E\left(G^{*}\right)
\end{array}\right.
$$

Finally, for all nonempty $S \subseteq V\left(G^{*}\right)$, we define the constrained-optimization problem $\mathrm{SPR}_{S}$ by:
$\left(\mathrm{SPR}_{S}\right):\left\{\begin{array}{rll}\max & \mu-v \\ \text { s.t. } & \text { Equation (16) } & \\ & \text { Equations (17), (18), and (19) } & \text { for all } i \in S \\ & \text { Inequality (20) } & \text { for all distinct } i, j \in S \\ & \left(\alpha_{i}, f_{i}, g_{i}\right) \in[0,1] \times F_{i} \times G_{i} \quad & \text { for all } i \in S \\ & \mu, v \in \mathbb{R}\end{array}\right.$.
For completeness, we state and prove the following observation.
Proposition A.1. Let $W \in \mathcal{W}$ such that $\operatorname{spr}(W)=\max _{U \in \mathcal{W}} \operatorname{spr}(U)$ and write $\mu, \nu$ for the maximum and minimum eigenvalues of $W$, with corresponding unit eigenfunctions $f, g$. Then for some nonempty set $S \subseteq V\left(G^{*}\right)$, the following holds. There exists a triple $\left(I_{i}, f_{i}, g_{i}\right)_{i \in S}$, where $\left(I_{i}\right)_{i \in S}$ is a labeled partition of $[0,1]$ with parts of positive measure and $f_{i}, g_{i} \in \mathbb{R}$ for all $i \in S$, such that:
(i) $W=W_{\jmath}$.
(ii) Allowing the replacement of $f$ by $-f$ and of $g$ by -g , for all $i \in S$, $f$ and $g$ equal $f_{i}$ and $g_{i}$ a.e. on $I_{i}$.
(iii) With $\alpha_{i}:=m\left(I_{i}\right)$ for all $i \in S, \operatorname{SPR}_{S}$ is solved by $\mu, \nu$, and $\left(\alpha_{i}, f_{i}, g_{i}\right)_{i \in S}$.

Proof. First we prove Item (ii). By Theorem 4.1 and the definition of $G^{*}$, there exist a nonempty set $S \subseteq V\left(G^{*}\right)$ and a labeled partition $\mathcal{J}=\left(I_{i}\right)_{i \in S}$ such that $W=W_{\mathcal{J}}$. By merging any parts of measure 0 into some part of positive measure, we may assume without loss of generality that $m\left(I_{i}\right)>0$ for all $i \in S$. So Item (ii) holds.

For Item (iii), the eigenfunctions corresponding to the maximum and minimum eigenvalues of a stepgraphon must be constant on each block by convexity and the Courant-Fischer Min-Max Theorem.

Finally, we prove Item (iiii), we first prove that for all $i \in V\left(G^{*}\right),\left(f_{i}, g_{i}\right) \in F_{i} \times G_{i}$. By Lemma 3.7 ,

$$
\mu f_{i}^{2}-v g_{i}^{2}=\mu-v
$$

for all $i \in S$. In particular, either $f_{i}^{2} \leq 1 \leq g_{i}^{2}$ or $g_{i}^{2} \leq 1 \leq f_{i}^{2}$. By Lemma 3.6, for all $i, j \in S, f_{i} f_{j}-g_{i} g_{j} \neq 0$ and $i j \in E(G)$ if and only if $f_{i} f_{j}-g_{i} g_{j}>0$. Note that the loops of $G^{*}$ are 1,2 , and 5 . It follows that for all $i \in S, f_{i}^{2}>1>g_{i}^{2}$ if and only if $i \in\{1,2,5\}$, and $g_{i}^{2}>1>f_{i}^{2}$, otherwise. Since $f$ is positive on [0,1], this completes the proof that $f_{i} \in F_{i}$ for all $i \in S$. Similarly since $g$ is positive on $\bigcup_{i \in\{1,2,3,4\} \cap S} I_{i}$ and negative on $\bigcup_{i \in\{5,6,7\}} I_{i}$, by inspection $g_{i} \in G_{i}$ for all $i \in S$. Similarly, Inequalities (20) follow directly from Lemma 3.6.

Continuing, we note the following. Since $W$ is a stepgraphon, if $\lambda \neq 0$ is an eigenvalue of $W$, there exists a $\lambda$-eigenfunction $h$ for $W$ such that for all $i \in S, h=h_{i}$ on $I_{i}$ for some $h_{i} \in \mathbb{R}$. Moreover for all $i \in S$, since $m\left(I_{i}\right)>0$,

$$
\lambda h_{i}=\sum_{i \in S} \alpha_{i} h_{i}
$$

In particular, any solution to $\mathrm{SPR}_{S}$ is at most $\mu-\nu$. Since $f, g$ are eigenfunctions corresponding to $W$ and the eigenvalues $\mu, \nu$, respectively, Equations (18) and (19) hold. Finally since $\left(I_{i}\right)_{i \in S}$ is a partition of $[0,1]$ and since $\|f\|_{2}^{2}=\|g\|_{2}^{2}=1$, Equation (16) holds. So $\mu, \nu$, and $\left(\alpha_{i}, f_{i}, g_{i}\right)_{i \in S}$ lie in the domain of SPR $_{S}$. This completes the proof of item (iiii), and the desired claim.

We enhance Proposition A. 1 as follows.
Lemma A.2. Proposition A.1 holds with the added assumption that $S \in \mathcal{S}_{17}$.
Proof. We begin our proof with Claim A.
Claim A. Suppose $i \in S$ and $j \in V\left(G^{*}\right)$ are distinct such that $N_{i} \cap S=N_{j} \cap S$. Then Proposition A.1 holds with the set $S^{\prime}:=(S \backslash\{i\}) \cup\{j\}$ replacing $S$.

First, we define the following quantities. For all $k \in S^{\prime} \backslash\{j\}$, let $\left(f_{k}^{\prime}, g_{k}^{\prime}, I_{k}^{\prime}\right):=$ $\left(f_{k}, g_{k}, I_{k}\right)$, and also let $\left(f_{j}^{\prime}, g_{j}^{\prime}\right):=\left(f_{i}, g_{i}\right)$. If $j \in S$, let $I_{j}^{\prime}:=I_{i} \cup I_{j}$, and otherwise, let $I_{j}^{\prime}:=I_{i}$. Additionally let $\mathcal{J}^{\prime}:=\left(I_{k}^{\prime}\right)_{k \in S^{\prime}}$ and for each $k \in S^{\prime}$, let $\alpha_{k}^{\prime}:=m\left(I_{k}^{\prime}\right)$. By the criteria from Proposition A.1, the domain criterion $\left(\alpha_{k}^{\prime}, f_{k}^{\prime}, g_{k}^{\prime}\right) \in[0,1] \times F_{k} \times G_{k}$ as well as Equation (17) holds for all $k \in S^{\prime}$. Since we are reusing $\mu, \nu$, the constraint $\mu, \nu \in \mathbb{R}$ also holds.

It suffices to show that Equation (16) holds, and that Equations (18) and (19) hold for all $k \in S^{\prime}$. To do this, we first note that for all $k \in S^{\prime}, f=f_{k}^{\prime}$ and $g=g_{k}^{\prime}$ on $I_{k}^{\prime}$. By definition, $f=f_{k}$ and $g=g_{k}$ on $I_{k}^{\prime}=I_{k}$ for all $k \in S^{\prime} \backslash\{j\}$ as needed by Claim A. Now suppose $j \notin S$. Then $f=f_{i}=f_{j}^{\prime}$ and $g=g_{i}=g_{j}^{\prime}$ and $I_{j}^{\prime}=I_{i}$ on the set $I_{i}=I_{j}^{\prime}$, matching Claim A. Finally, suppose $j \in S$. Note by definition that $f=f_{i}=f_{j}^{\prime}$ and $g=g_{i}=g_{j}^{\prime}$ on $I_{i}$. Since and $I_{j}^{\prime}=I_{i} \cup I_{j}$, it suffices to prove that $f=f_{j}^{\prime}$ and $g=g_{j}^{\prime}$ on $I_{j}$. We first show that $f_{j}=f_{i}$ and $g_{j}=g_{i}$. Indeed,

$$
\mu f_{j}=\sum_{k \in N_{j} \cap S} \alpha_{k} f_{k}=\sum_{k \in N_{i} \cap S} \alpha_{k} f_{k}=\mu f_{i}
$$

and since $\mu \neq 0, f_{j}=f_{i}$. Similarly, $g_{j}=g_{i}$. So $f=f_{j}=f_{i}=f_{j}^{\prime}$ and $g=g_{j}=g_{i}=g_{j}^{\prime}$ on the set $I_{j}^{\prime}=I_{i} \cup I_{j}$.

Finally, we claim that $W_{\mathcal{J}^{\prime}}=W$. Indeed, this follows directly from Lemma 3.6 and the fact that $W=W_{J}$. Since $\mathcal{J}^{\prime}$ is a partition of $[0,1]$ and since $f, g$ are unit eigenfunctions for $W$ Equation (16) holds, and Equations (18) and (19) hold for all $k \in S^{\prime}$. This completes the proof of Claim $A$.

Next, we prove Claim B.
Claim B. If $S$ satisfies the criteria of Proposition A.1, then without loss of generality the following holds.
(a) If there exists some $i \in S$ such that $N_{i}=S$, then $i=1$.
(b) $S \cap\{1,2,3,4\} \neq \emptyset$.
(c) $S \cap\{2,3,4\}$ is one of $\emptyset,\{4\},\{2,4\}$, and $\{2,3,4\}$.
(d) $S \cap\{5,6,7\}$ is one of $\{7\},\{5,7\}$, and $\{5,6,7\}$.

Since $N_{1} \cap S=S=N_{i}$, item (a) follows from Claim A applied to the pair $(i, 1)$. Since $f, g$ are orthogonal and $f$ is positive on $[0,1], g$ is positive on a set of positive measure, so item (b) holds.

To prove item (C), we have 4 cases. If $S \cap\{2,3,4\}=\{2\}$, then $N_{2} \cap S=N_{1} \cap S$ and we may apply Claim A to the pair (2,1). If $S \cap\{2,3,4\}=\{3\}$ or $\{3,4\}$, then $N_{3} \cap S=N_{4} \cap S$ and we may apply Claim A to the pair $(3,4)$. If $S \cap\{2,3,4\}=\{2,3\}$, then $N_{2} \cap S=N_{1} \cap S$ and we may apply Claim A to the pair ( 2,1 ). So item (d) holds. For item (d), we reduce $S \cap\{5,6,7\}$ to one of $\emptyset,\{7\},\{5,7\}$, and $\{5,6,7\}$ in the same fashion. To eliminate the case where $S \cap\{5,6,7\}=\emptyset$, we simply note that since $f$ and $g$ are orthogonal and $f$ is positive on $[0,1], g$ is negative on a set of positive measure. This completes the proof of Claim B.

TABLE 2. The 21 sets which arise from repeated applications of Claim B

|  | $\emptyset$ | 4 | $\{2,4\}$ | $\{2,3,4\}$ |
| :---: | :---: | ---: | ---: | ---: |
| $\{7\}$ | $1 \mid 7$ | $4 \mid 7$ | $24 \mid 7$ | $234 \mid 7$ |
|  |  | $1\|4\| 7$ | $1\|24\| 7$ | $1\|234\| 7$ |
| 5,7 | $1 \mid 57$ | $4 \mid 57$ | $24 \mid 57$ | $234 \mid 57$ |
|  |  | $1\|4\| 57$ | $1\|24\| 57$ | $1\|234\| 57$ |
| $5,6,7$ | $1 \mid 567$ | $4 \mid 567$ | $24 \mid 567$ | $234 \mid 567$ |
|  |  | $1\|4\| 567$ | $1\|24\| 567$ | $1\|234\| 567$ |

After repeatedly applying Claim $B$, we may replace $S$ with one of the cases found in Table Let $\mathcal{S}_{21}$ denote the sets in Table 2 . By definition,

$$
\mathcal{S}_{21}=\mathcal{S}_{17} \bigcup\{\{4,7\},\{2,4,7\},\{2,3,4,7\},\{2,3,4,5,7\}\}
$$

Finally, we eliminate the 4 cases in $\mathcal{S}_{21} \backslash \mathcal{S}_{17}$. If $S=\{4,7\}$, then $W$ is a bipartite graphon, hence $\operatorname{spr}(W) \leq 1$, a contradiction since $\max _{U \in \mathcal{W}} \operatorname{spr}(W)>1$.

For the three remaining cases, let $\tau$ be the permutation on $\{2, \ldots, 7\}$ defined as follows. For all $i \in\{2,3,4\}, \tau(i):=i+3$ and $\tau(i+3):=i$. If $S$ is among $\{2,4,7\},\{2,3,4,7\}$, $\{2,3,4,5,7\}$, we apply $\tau$ to $S$ in the following sense. Replace $g$ with $-g$ and replace $\left(\alpha_{i}, I_{i}, f_{i}, g_{i}\right)_{i \in S}$ with $\left(\alpha_{\tau(i)}, I_{\tau(i)}, f_{\tau(i)},-g_{\tau(i)}\right)_{i \in \tau(S)}$. By careful inspection, it follows that
$\tau(S)$ satisfies the criteria from Proposition A.1. Since $\tau(\{2,4,7\})=\{4,5,7\}, \tau(\{2,3,4,7\})$ $=\{4,5,6,7\}$, and $\tau(\{2,3,4,5,7\})=\{2,4,5,6,7\}$, this completes the proof.
A.2. Proof of Lemma 5.4. Let $\left(\alpha_{4}, \alpha_{5}, \alpha_{7}\right)$ be a solution to $\operatorname{SPR}_{457}$.

First, let $T:=\left\{\left(\varepsilon_{1}, \varepsilon_{2}\right) \in(-1 / 3,2 / 3) \times(-2 / 3,1 / 3): \varepsilon_{1}+\varepsilon_{2} \in(0,1)\right\}$, and for all $\varepsilon=\left(\varepsilon_{1}, \varepsilon_{2}\right) \in T$, let

$$
M(\varepsilon):=\left[\begin{array}{ccc}
2 / 3-\varepsilon_{1} & 0 & 1 / 3-\varepsilon_{2} \\
0 & 0 & 1 / 3-\varepsilon_{2} \\
2 / 3-\varepsilon_{1} & \varepsilon_{1}+\varepsilon_{2} & 0
\end{array}\right]
$$

As a motivation, suppose $\mu, \nu$, and $\left(\alpha_{4}, \alpha_{5}, \alpha_{7}\right)$ are part of a solution to $\operatorname{SPR}_{\{4,5,7\}}$. Then with $\varepsilon:=\left(\varepsilon_{1}, \varepsilon_{2}\right)=\left(2 / 3-\alpha_{5}, 1 / 3-\alpha_{4}\right), \varepsilon \in T$ and $\mu, \nu$ are the maximum and minimum eigenvalues of $M(\varepsilon)$, respectively. By the end of the proof, we show that any solution of $\operatorname{SPR}_{\{4,5,7\}}$ has $\alpha_{7}=0$.

To proceed, we prove the following claims.
Claim A. For all $\varepsilon \in T, M(\varepsilon)$ has two distinct positive eigenvalues and one negative eigenvalue.

Since $M(\varepsilon)$ is diagonalizable, it has 3 real eigenvalues which we may order as $\mu \geq$ $\delta \geq \nu$. Since $\mu \delta \nu=\operatorname{det}(M(\varepsilon))=-\alpha_{4} \alpha_{5} \alpha_{7} \neq 0<0, M(\varepsilon)$ has an odd number of negative eigenvalues. Since $0<\alpha_{5}=\mu+\delta+\nu$, it follows that $\mu \geq \delta>0>\nu$. Finally, note by the Perron-Frobenius Theorem that $\mu>\delta$. This completes the proof of Claim A.

Next, we define the following quantities, treated as functions of $\varepsilon$ for all $\varepsilon \in T$. For convenience, we suppress the argument " $\varepsilon$ " in most places. Let $k(x)=a x^{3}+b x^{2}+c x+d$ be the characteristic polynomial of $M(\varepsilon)$. By inspection,

$$
\begin{aligned}
& a=1 \quad b=\varepsilon_{1}-\frac{2}{3} \\
& c=\frac{\left(3 \varepsilon_{2}+2\right)\left(3 \varepsilon_{2}-1\right)}{9}
\end{aligned} \quad d=\frac{\left(\varepsilon_{1}+\varepsilon_{2}\right)\left(3 \varepsilon_{1}-2\right)\left(3 \varepsilon_{2}-1\right)}{9} .
$$

Continuing, let

$$
\begin{array}{lr}
p:=\frac{3 a c-b^{2}}{3 a^{2}} & q:=\frac{2 b^{3}-9 a b c+27 a^{2} d}{27 a^{3}}, \\
A:=2 \sqrt{\frac{-p}{3}} & B:=\frac{-b}{3 a}, \\
\phi:=\arccos \left(\frac{3 q}{A p}\right) . &
\end{array}
$$

Let $S(\varepsilon)$ be the difference between the maximum and minimum eigenvalues of $M(\varepsilon)$. We show Claim B.

Claim B. For all $\varepsilon \in T$,

$$
S(\varepsilon)=\sqrt{3} \cdot A(\varepsilon) \cdot \cos \left(\frac{2 \phi(\varepsilon)-\pi}{6}\right)
$$

Moreover, $S$ is analytic on $T$.

Indeed, by Viéte's Formula, using the fact that $k(x, y)$ has exactly 3 distinct real roots, the quantities $a(\varepsilon), \ldots, \phi(x, y)$ are analytic on $T$. Moreover, the eigenvalues of $M(\varepsilon)$ are $x_{0}, x_{1}, x_{2}$ where, for all $k \in\{0,1,2\}$,

$$
x_{k}(\varepsilon)=A(\varepsilon) \cdot \cos \left(\frac{\phi+2 \pi \cdot k}{3}\right)+B(\varepsilon)
$$

Moreover, $x_{0}(\varepsilon), x_{1}(\varepsilon), x_{2}(\varepsilon)$ are analytic on $T$. For all $k, \ell \in\{1,2,3\}$, let

$$
D(k, \ell, x):=\cos \left(x+\frac{2 \pi k}{3}\right)-\cos \left(x+\frac{2 \pi \ell}{3}\right) .
$$

For all $(k, \ell) \in\{(0,1),(0,2),(2,1)\}$, note the trigonometric identities

$$
D(k, \ell, x)=\sqrt{3} \cdot\left\{\begin{aligned}
\cos \left(x-\frac{\pi}{6}\right), & (k, \ell)=(0,1) \\
\cos \left(x+\frac{\pi}{6}\right), & (k, \ell)=(0,2) \\
\sin (x), & (k, \ell)=(2,1)
\end{aligned}\right.
$$

By inspection, for all $x \in(0, \pi / 3)$,

$$
D(0,1)>\max \{D(0,2), D(2,1)\} \geq \min \{D(0,2), D(2,1)\} \geq 0 .
$$

Since $A>0$ and $\phi \in(0, \pi / 3)$, the claimed equality holds. Since $x_{0}(\varepsilon), x_{1}(\varepsilon)$ are analytic, $S(\varepsilon)$ is analytic on $T$. This completes the proof of Claim $B$.

Next, we compute the derivatives of $S(\varepsilon)$ on $T$. For convenience, denote by $A_{i}, \phi_{i}$, and $S_{i}$ for the partial derivatives of $A$ and $\phi$ by $\varepsilon_{i}$, respectively, for $i \in\{1,2\}$. Furthermore, let

$$
\psi(\varepsilon):=\frac{2 \phi(\varepsilon)-\pi}{6} .
$$

Claim follows directly from Claim $B$.
Claim C. For all $i \in T$, then on the set $T$, we have

$$
3 S_{i}=3 A_{i} \cdot \cos (\psi)-\cdot A \phi_{i} \sin (\psi) .
$$

Moreover, each expression is analytic on $T$.
Finally, we solve $\operatorname{SPR}_{\{4,5,7\}}$.
Claim D. If $\left(\alpha_{4}, \alpha_{5}, \alpha_{7}\right)$ is a solution to $\operatorname{SPR}_{\{4,5,7\}}$, then $0 \in\left\{\alpha_{4}, \alpha_{5}, \alpha_{7}\right\}$.
With $\left(\alpha_{4}, \alpha_{5}, \alpha_{7}\right):=\left(1 / 3-\varepsilon_{2}, 2 / 3-\varepsilon_{1}, \varepsilon_{1}+\varepsilon_{2}\right)$ and using the fact that $S$ is analytic on $T$, it is sufficient to eliminate all common zeroes of $S_{1}$ and $S_{2}$ on $T$. With the help of a computer algebra system and the formulas for $S_{1}$ and $S_{2}$ from Claim C, we replace the system $S_{1}=0$ and $S_{2}=0$ with a polynomial system of equations $P=0$ and $Q=0$ whose real solution set contains all previous solutions. Here,

$$
P(\varepsilon)=9 \varepsilon_{1}^{3}+18 \varepsilon_{1}^{2} \varepsilon_{2}+54 \varepsilon_{1} \varepsilon_{2}^{2}+18 \varepsilon_{2}^{3}-15 \varepsilon_{1}^{2}-33 \varepsilon_{1} \varepsilon_{2}-27 \varepsilon_{2}^{2}+5 \varepsilon_{1}+\varepsilon_{2}
$$

and $Q=43046721 \varepsilon_{1}^{18} \varepsilon_{2}+\cdots+\left(-532480 \varepsilon_{2}\right)$ is a polynomial of degree 19 , with coefficients between -184862311457373 and 192054273812559. For brevity, we do not express $Q$ explicitly.

To complete the proof of Claim $\mathbb{D}$, it suffices to show that no common real solution to $P=Q=0$ which lies in $T$ also satisfies $S_{1}=S_{2}=0$. Again using a computer algebra system, we first find all common zeroes of $P$ and $Q$ on $\mathbb{R}^{2}$. Included are the rational
solutions $(2 / 3,-2 / 3),(-1 / 3,1 / 3),(0,0),(2 / 3,1 / 3)$, and $(2 / 3,-1 / 6)$ which do not lie in T. Furthermore, the solution ( $1.2047 \ldots, 0.0707 \ldots$ ) may also be eliminated. For the remaining 4 zeroes, $S_{1}, S_{2} \neq 0$. A notebook showing these calculations can be found in [27].

Claim E. If $\mu, \nu$, and $\alpha=\left(\alpha_{4}, \alpha_{5}, \alpha_{7}\right)$ is part of a solution to $\operatorname{SPR}_{\{4,5,7\}}$ such that $\mu-\nu \geq 1$, then $\alpha_{7}=0$.

By definition of $\operatorname{SPR}_{\{4,5,7\}}, \mu$ and $\nu$ are eigenvalues of the matrix

$$
N(\alpha):=\left[\begin{array}{ccc}
\alpha_{5} & 0 & \alpha_{4} \\
0 & 0 & \alpha_{4} \\
\alpha_{5} & \alpha_{7} & 0
\end{array}\right]
$$

Furthermore, $N(\alpha)$ has characteristic polynomial

$$
p(x)=x^{3}-\alpha_{5} x^{2}-\alpha_{4} \cdot\left(\alpha_{5}+\alpha_{7}\right)+\alpha_{4} \alpha_{5} \alpha_{7} .
$$

Recall that $\alpha_{4}+\alpha_{5}+\alpha_{7}=1$. By Claim $\mathbb{D}, 0 \in\{4,5,7\}$, and it follows that $p \in\left\{p_{4}, p_{5}, p_{7}\right\}$ where

$$
\begin{aligned}
& p_{4}(x):=x^{2} \cdot\left(x-\alpha_{5}\right), \\
& p_{5}(x):=x \cdot\left(x^{2}-\alpha_{4}\left(1-\alpha_{4}\right)\right), \text { and } \\
& p_{7}(x):=x \cdot\left(x^{2}-\left(1-\alpha_{4}\right) x-\alpha_{4}\left(1-\alpha_{4}\right)\right) .
\end{aligned}
$$

If $p=p_{4}$, then $\mu-v=\alpha_{5} \leq 1$, and if $p=p_{5}$, then $\mu-v=2 \sqrt{\alpha_{4}\left(1-\alpha_{4}\right)} \leq 1$. So $p=p_{7}$, which completes the proof of Claim E.

This completes the proof of Lemma 5.4.
A.3. Proof of Lemma 6.3. First, we find $S_{z}\left(\varepsilon_{1}, \varepsilon_{3}\right)$ using Viète's Formula. In doing so, we define functions $k_{z}\left(\varepsilon_{1}, \varepsilon_{2} ; x\right), \ldots, \delta_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right)$. To ease the burden on the reader, we suppress the subscript $z$ and the arguments $\varepsilon_{1}, \varepsilon_{2}$ when convenient and unambiguous. Let $k(x)=a x^{3}+b x^{2}+c x+d$ be the characteristic polynomial of $M_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right)$. By inspection,

$$
\begin{array}{lr}
a=1 & b=\varepsilon_{1}+z-\frac{2}{3} \\
c=\frac{\left(3 \varepsilon_{2}+2\right)\left(3 \varepsilon_{2}-1\right)}{9} & d=\frac{\left(\varepsilon_{1}+\varepsilon_{2}\right)\left(3 \varepsilon_{1}+3 z-2\right)\left(3 \varepsilon_{2}-1\right)}{9} .
\end{array}
$$

Continuing, let

$$
\begin{array}{lr}
p:=\frac{3 a c-b^{2}}{3 a^{2}} & q:=\frac{2 b^{3}-9 a b c+27 a^{2} d}{27 a^{3}}, \\
A:=2 \sqrt{\frac{-p}{3}} & B:=\frac{-b}{3 a}, \\
\phi:=\arccos \left(\frac{3 q}{A p}\right) . &
\end{array}
$$

By Viète's Formula, the roots of $k_{z}\left(\varepsilon_{1}, \varepsilon_{2} ; x\right)$ are the suggestively defined quantities:

$$
\begin{aligned}
\mu:=A \cos \left(\frac{\phi}{3}\right)+B & \nu:=A \cos \left(\frac{\phi+2 \pi}{3}\right)+B, \\
\delta:=A \cos \left(\frac{\phi+4 \pi}{3}\right)+B . &
\end{aligned}
$$

First, We prove Claim A.
Claim A. If $\left(\varepsilon_{1}, \varepsilon_{2}, z\right)$ is sufficiently close to $(0,0,0)$, then

$$
\begin{equation*}
S_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right)=A_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right) \sqrt{3} \cdot \cos \left(\frac{2 \phi_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right)-\pi}{6}\right) . \tag{21}
\end{equation*}
$$

Indeed, suppose $z>0$ and $z \rightarrow 0$. Then for all $\left(\varepsilon_{1}, \varepsilon_{2}\right) \in(-3 z, 3 z), \varepsilon_{1}, \varepsilon_{2} \rightarrow 0$. With the help of a computer algebra system, we substitute in $z=0$ and $\varepsilon_{1}, \varepsilon_{2}=0$ to find the limits:

$$
\begin{aligned}
(a, b, c, d) & \rightarrow\left(1, \frac{-2}{3}, \frac{-2}{9}, 0\right) \\
(p, q) & \rightarrow\left(\frac{-10}{27}, \frac{-52}{729}\right) \\
(A, B, \phi) & \rightarrow\left(\frac{2 \sqrt{10}}{9}, \frac{2}{9}, \arccos \left(\frac{13 \sqrt{10}}{50}\right)\right) .
\end{aligned}
$$

Using a computer algebra system, these substitutions imply that

$$
(\mu, \nu, \delta) \rightarrow(0.9107 \ldots,-0.2440 \ldots, 0)
$$

So for all $z$ sufficiently small, $S=\mu-\nu$. After some trigonometric simplification,

$$
\mu-v=A \cdot\left(\cos \left(\frac{\phi}{3}\right)-\cos \left(\frac{\phi+2 \phi}{3}\right)\right)=A \sqrt{3} \cdot \cos \left(\frac{2 \phi-\pi}{6}\right)
$$

and Equation (21). This completes the proof of Claim A.
Now we prove Claim B.
Claim B. There exists a constant $C_{0}^{\prime}>0$ such that the following holds. If $|z|$ is sufficiently small, then $S_{z}$ is concave-down on $\left[-C_{0}, C_{0}\right]^{2}$ and strictly decreasing on $\left[-C_{0}, C_{0}\right]^{2} \backslash\left[-C_{0} z, C_{0} z\right]^{2}$.

First, we define

$$
D_{z}\left(\varepsilon_{1}, \varepsilon_{2}\right):=\left.\left(\frac{\partial^{2} S_{z}}{\partial \varepsilon_{1}^{2}} \cdot \frac{\partial^{2} S_{z}}{\partial \varepsilon_{2}^{2}}-\left(\frac{\partial^{2} S_{z}}{\partial \varepsilon_{1} \partial \varepsilon_{2}}\right)^{2}\right)\right|_{\left(\varepsilon_{1}, \varepsilon_{2}, z\right)}
$$

As a function of $\left(\varepsilon_{1}, \varepsilon_{2}\right), D_{z}$ is the determinant of the Hessian matrix of $S_{z}$. Using a computer algebra system, we note that

$$
\begin{aligned}
D_{0}(0,0) & =22.5 \ldots, \quad \text { and } \\
\left.\left(\frac{\partial^{2} S}{\partial \varepsilon_{1}^{2}}, \frac{\partial^{2} S}{\partial \varepsilon_{1} \partial \varepsilon_{2}}, \frac{\partial^{2} S}{\partial \varepsilon_{2}^{2}}\right)\right|_{(0,0,0)} & =(-8.66 \ldots,-8.66 \ldots,-11.26 \ldots) .
\end{aligned}
$$

Since $S$ is analytic to $(0,0,0)$, there exist constants $C_{1}, C_{2}>0$ such that the following holds. For all $z \in\left[-C_{1}, C_{1}\right], S_{z}$ is concave-down on $\left[-C_{1}, C_{1}\right]^{2}$. This completes the
proof of the first claim. Moreover for all $z \in\left[-C_{1}, C_{1}\right]$ and for all $\left(\varepsilon_{1}, \varepsilon_{2}\right) \in\left[-C_{1}, C_{1}\right]^{2}$,

$$
\left.\max \left\{\frac{\partial^{2} S_{z}}{\partial \varepsilon_{1}^{2}}, \frac{\partial^{2} S_{z}}{\partial \varepsilon_{1} \partial \varepsilon_{2}}, \frac{\partial^{2} S_{z}}{\partial \varepsilon_{2}^{2}}\right\}\right|_{\left(\varepsilon_{1}, \varepsilon_{2}, z\right)} \leq-C_{2} .
$$

To complete the proof of the second claim, note also that since $S$ is analytic at $(0,0,0)$, there exist constants $C_{3}, C_{4}>0$ such that for all $z \in\left[-C_{3}, C_{3}\right]$ and all $\left(\varepsilon_{1}, \varepsilon_{2}\right) \in$ $\left[-C_{3}, C_{3}\right]^{2}$,

$$
\frac{\partial^{2} S}{\partial z \partial \varepsilon_{i}} \leq C_{4}
$$

Since $(0,0)$ is a local maximum of $S_{0}$,

$$
\begin{aligned}
\left.\frac{\partial S}{\partial \varepsilon_{i}}\right|_{\left(\varepsilon_{1}, \varepsilon_{2}, z\right)} & =\left.\frac{\partial S}{\partial \varepsilon_{i}}\right|_{(0,0,0)}+\left.\int_{w=0}^{z} \frac{\partial^{2} S}{\partial z \partial \varepsilon_{i}}\right|_{(0,0, w)} d w+\left.\int_{\mathbf{u}=(0,0)}^{\left(\varepsilon_{1}, \varepsilon_{2}\right)} \frac{\partial^{2} S}{\partial \mathbf{u} \partial \varepsilon_{i}}\right|_{(\mathbf{u}, z)} d \mathbf{u} \\
& \leq C_{4} \cdot z-C_{2} \cdot\left\|\left(\varepsilon_{1}, \varepsilon_{2}\right)\right\|_{2} .
\end{aligned}
$$

Since $C_{2}, C_{4}>0$, this completes the proof of Claim $B$.
Next, we prove Claim C.
Claim C. If $z$ is sufficiently small, then $\mathcal{P}_{z, C_{0}}$ is solved by a unique point $\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}\right)=$ $\left(\varepsilon_{1}^{*}(z), \varepsilon_{2}^{*}(z)\right)$. Moreover as $z \rightarrow 0$,

$$
\begin{equation*}
\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}\right)=\left((1+o(z)) \frac{7 z}{30},(1+o(z)) \frac{-z}{3}\right) . \tag{22}
\end{equation*}
$$

Indeed, the existence of a unique maximum $\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}\right)$ on $\left[-C_{0}, C_{0}\right]^{2}$ follows from the fact that $S_{z}$ is strictly concave-down and bounded on $\left[-C_{0}, C_{0}\right]^{2}$ for all $z$ sufficiently small. Since $S_{z}$ is strictly decreasing on $\left[-C_{0}, C_{0}\right]^{2} \backslash\left(-C_{0} z, C_{0} z\right)^{2}$, it follows that $\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}\right)$ $\in\left(-C_{0} z, C_{0} z\right)$. For the second claim, note that since $S$ is analytic at $(0,0,0)$,

$$
0=\left.\frac{\partial S}{\partial \varepsilon_{i}}\right|_{\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}, z\right)}=\sqrt{3} \cdot\left(\frac{\partial A}{\partial \varepsilon_{i}} \cdot \cos \left(\frac{2 \phi-\pi}{6}\right)-\frac{A}{3} \cdot \frac{\partial \phi}{\partial \varepsilon_{i}} \cdot \sin \left(\frac{2 \phi-\pi}{6}\right)\right)
$$

for both $i=1$ and $i=2$. Let

$$
\tau_{i}:=\frac{3 \cdot \frac{\partial A}{\partial \varepsilon_{i}}}{A \cdot \frac{\partial \phi}{\partial \varepsilon_{i}}}
$$

for both $i=1$ and $i=2$. Then by Equation (21),

$$
\left.\arctan \left(\tau_{i}\right)\right|_{\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}, z\right)}=\left.\frac{2 \phi-\pi}{6}\right|_{\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}, z\right)}
$$

for both $i=1$ and $i=2$. We first consider linear approximation of the above quantities under the limit $\left(\varepsilon_{1}, \varepsilon_{2}, z\right) \rightarrow(0,0,0)$. Here, we write $f\left(\varepsilon_{1}, \varepsilon_{2}, z\right) \sim g\left(\varepsilon_{1}, \varepsilon_{2}, z\right)$ to mean that

$$
f\left(\varepsilon_{1}, \varepsilon_{2}, z\right)=\left(1+o\left(\max \left\{\left|\varepsilon_{1}\right|,\left|\varepsilon_{2}\right|,|z|\right\}\right)\right) \cdot g\left(\varepsilon_{1}, \varepsilon_{2}, z\right) .
$$

With the help of a computer algebra system, we note that

$$
\begin{aligned}
\arctan \left(\tau_{1}\right) & \sim \frac{-78 \varepsilon_{1}-96 \varepsilon_{2}-3 z-40 \arctan \left(\frac{1}{3}\right)}{40}, \\
\arctan \left(\tau_{2}\right) & \sim \frac{-64 \varepsilon_{1}-103 \varepsilon_{2}-14 z-20 \arctan \left(\frac{1}{3}\right)}{20}, \\
\frac{2 \phi-\pi}{6} & \sim \frac{108 \varepsilon_{1}+81 \varepsilon_{2}+18 z+20 \arccos \left(\frac{13 \sqrt{10}}{50}\right)-10 \pi}{60} .
\end{aligned}
$$

By inspection, the constant terms match due to the identity

$$
-\arctan \left(\frac{1}{3}\right)=\frac{1}{3} \arccos \left(\frac{13 \sqrt{10}}{50}\right)-\frac{\pi}{6}
$$

Since $\max \left\{\left|\varepsilon_{1}^{*}\right|,\left|\varepsilon_{2}^{*}\right|\right\} \leq C_{0} z$, replacing $\left(\varepsilon_{1}, \varepsilon_{2}\right)$ with $\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}\right)$ implies that

$$
\begin{aligned}
\frac{-78 \varepsilon_{1}^{*}-96 \varepsilon_{2}^{*}-3 z}{2} & =(1+o(z)) \cdot\left(36 \varepsilon_{1}^{*}+27 \varepsilon_{2}^{*}+6 z\right), \quad \text { and } \\
-64 \varepsilon_{1}^{*}-103 \varepsilon_{2}^{*}-14 z & =(1+o(z)) \cdot\left(36 \varepsilon_{1}^{*}+27 \varepsilon_{2}^{*}+6 z\right)
\end{aligned}
$$

as $z \rightarrow 0$. After applying Gaussian Elimination to this 3 -variable system of 2 equations, it follows that

$$
\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}\right)=\left((1+o(z)) \cdot \frac{7 z}{30},(1+o(z)) \cdot \frac{-z}{3}\right) .
$$

This completes the proof of Claim C.
For the next step, we prove Claim D. First, let $Q_{n}$ denote the program formed from $\mathcal{P}_{n^{-1}, C_{0}}$ subject to the added constraint that $n \cdot\left(\frac{2}{3}-\varepsilon_{1}\right), n \cdot\left(\frac{1}{3}-\varepsilon_{2}\right) \in \mathbb{Z}$.
Claim D. For all $n$ sufficiently large, $Q_{n}$ is solved by a unique point $\left(n_{1}^{*}, n_{3}^{*}\right)$ which satisfies $n_{1}^{*}+n_{3}^{*}=n$.

Note by Lemma 6.2 that for all $n$ sufficiently large,

$$
\max \left\{\left|\frac{n_{1}}{n}-\frac{2}{3}\right|,\left|\frac{n_{3}}{n}-\frac{1}{3}\right|\right\} \leq C_{0} .
$$

Moreover, by Claim C, $\mathcal{P}_{n^{-1}}$ is solved uniquely by

$$
\left(\varepsilon_{1}^{*}, \varepsilon_{2}^{*}\right)=\left((1+o(z)) \cdot \frac{7}{30 n},(1+o(z)) \cdot \frac{-1}{3 n}\right) .
$$

Since

$$
\frac{2 n}{3}-n \cdot \varepsilon_{1}^{*}=\frac{2 n}{3}-(1+o(1)) \cdot \frac{7}{30}
$$

and $7 / 30<1 / 3$, it follows for $n$ sufficiently large that $2 n / 3-n \cdot \varepsilon_{1}^{*} \in I_{1}$ where

$$
I_{1}:=\left\{\begin{array}{ll}
\left(\frac{2 n}{3}-1, \frac{2 n}{3}\right), & 3 \mid n \\
\left.\left(\left\lfloor\frac{2 n}{3}\right\rfloor, \left\lvert\, \frac{2 n}{3}\right.\right\rceil\right), & 3+n
\end{array} .\right.
$$

Similarly since

$$
n \cdot\left(\varepsilon_{1}^{*}+\varepsilon_{2}^{*}\right)=(1+o(1)) \cdot\left(\frac{7}{30}-\frac{1}{3}\right)=(1+o(1)) \cdot \frac{-1}{10}
$$

and $1 / 10<1 / 3$, it follows that $n \cdot\left(\varepsilon_{1}^{*}+\varepsilon_{2}^{*}\right) \in(-1,0)$. Altogether,

$$
\left(\frac{2 n}{3}-n \cdot \varepsilon_{1}, n \cdot\left(\varepsilon_{1}^{*}+\varepsilon_{2}^{*}\right)\right) \in I_{1} \times(-1,0)
$$

Note that to solve $Q_{n}$, it is sufficient to maximize $S_{n^{-1}}$ on the set

$$
\left[-C_{0}, C_{0}\right]^{2} \bigcap\left\{\left(n_{1} / n, n_{3} / n\right)\right\}_{u, v \in \mathbb{N}} .
$$

Since $S_{n^{-1}}$ is concave-down on $I_{1} \times(-1,0),\left(n_{1}^{*}, n-n_{1}^{*}-n_{3}^{*}\right)$ is a corner of the square $I_{1} \times(-1,0)$. So $n_{1}^{*}+n_{2}^{*}=n$, which implies Claim D. This completes the proof of the main result.

## Appendix B. A Computer-Assisted proof of Lemma 5.2

In this appendix, we derive a number of formulas that a stepgraphon corresponding to some set $S \subseteq\{1,2,3,4,5,6,7\}$ in Lemma A.2 satisfies, and detail how these formulas are used to provide a computer-assisted proof of Lemma 5.2.
B.1. Formulas. In this subsection, we derive the formulas used in our computerassisted proof, from the equations described in Section A.1. First, we define a number of functions which will ease the notational burden in the results that follow. Let

$$
\begin{aligned}
F_{1}(x):= & (\mu+\nu) x+2 \mu \nu, \\
F_{2}(x):= & 2(\mu \nu+(\mu+\nu) x)^{2}+(\mu+\nu) x^{3}, \\
F_{3}(x):= & 4 \mu^{2} \nu^{2} \cdot(\mu \nu+(\mu+\nu) x)^{2} \\
& -2(\mu+\nu) x^{3} \cdot((\mu+\nu) x+\mu \nu)((\mu+\nu) x+3 \mu \nu) \\
& -(\mu+\nu) x^{5} \cdot(2 \mu \nu+(\mu+\nu) x), \\
F_{4}(x):= & 4 \mu^{2} \nu^{2} x \cdot((3(\mu+\nu) x+\mu \nu) \cdot(2(\mu+\nu) x+\mu \nu)-\mu \nu(\mu+\nu) x) \\
& +4(\mu+\nu) x^{4} \cdot\left(((\mu+\nu) x+\mu \nu)^{2}+(\mu+\nu)^{2} \cdot((\mu+\nu) x+4 \mu \nu)\right) \\
& +(\mu+\nu)^{2} x^{7} .
\end{aligned}
$$

Letting $S:=\left\{i \in\{1, \ldots, 7\}: \alpha_{i}>0\right\}$, we prove the following six formulas.
Proposition B.1. Let $i \in\{1,2,5\} \cap S$ and $j \in\{3,4,6,7\} \cap S$ be such that $N_{i} \cap S=$ $\left(N_{j} \cap S\right) \dot{\cup}\{j\}$. Then

$$
f_{j}^{2}=\frac{\left(\alpha_{j}+2 \nu\right) \mu}{F_{1}\left(\alpha_{j}\right)}, \quad g_{j}^{2}=\frac{\left(\alpha_{j}+2 \mu\right) \nu}{F_{1}\left(\alpha_{j}\right)}
$$

and

$$
f_{i}=\left(1+\frac{\alpha_{j}}{\mu}\right) f_{j}, \quad g_{i}=\left(1+\frac{\alpha_{j}}{v}\right) g_{j}
$$

Moreover, $F_{1}\left(\alpha_{j}\right)$ and $\alpha_{j}+2 v$ are negative.
Proof. By Lemma 3.7,

$$
\begin{aligned}
& \mu f_{i}^{2}-v g_{i}^{2}=\mu-v, \\
& \mu f_{j}^{2}-v g_{j}^{2}=\mu-v .
\end{aligned}
$$

By taking the difference of the eigenvector equations for $f_{i}$ and $f_{j}$ (and also $g_{i}$ and $g_{j}$ ), we obtain

$$
\begin{aligned}
\alpha_{j} f_{j} & =\mu\left(f_{i}-f_{j}\right), \\
\alpha_{j} g_{j} & =v\left(g_{i}-g_{j}\right),
\end{aligned}
$$

or, equivalently,

$$
\begin{aligned}
& f_{i}=\left(1+\frac{\alpha_{j}}{\mu}\right) f_{j} \\
& g_{i}=\left(1+\frac{\alpha_{j}}{v}\right) g_{j}
\end{aligned}
$$

This leads to the system of equations

$$
\left[\begin{array}{cc}
\mu & -v \\
\mu \cdot\left(1+\frac{\alpha_{j}}{\mu}\right)^{2} & -v \cdot\left(1+\frac{\alpha_{j}}{v}\right)^{2}
\end{array}\right] \cdot\left[\begin{array}{c}
f_{j}^{2} \\
g_{j}^{2}
\end{array}\right]=\left[\begin{array}{l}
\mu-v \\
\mu-v
\end{array}\right] .
$$

If the corresponding matrix is invertible, then after substituting the claimed formulas for $f_{j}^{2}, g_{j}^{2}$ and simplifying, it follows that they are the unique solutions. To verify that $F_{1}\left(\alpha_{j}\right)$ and $\alpha_{j}+2 v$ are negative, it is sufficient to inspect the formulas for $f_{j}$ and $g_{j}$, noting that $\nu$ is negative and both $\mu$ and $\alpha_{j}+2 \mu$ are positive.

Suppose the matrix is not invertible. By assumption $\mu, \nu \neq 0$, and so

$$
\left(1+\frac{\alpha_{j}}{\mu}\right)^{2}=\left(1+\frac{\alpha_{j}}{v}\right)^{2}
$$

But, since $i \in\{1,2,5\}$ and $j \in\{3,4,6,7\}$,

$$
1>f_{j}^{2} g_{i}^{2}=f_{i}^{2} g_{i}^{2} \cdot\left(1+\frac{\alpha_{j}}{\mu}\right)^{2}=f_{i}^{2} g_{i}^{2} \cdot\left(1+\frac{\alpha_{j}}{\nu}\right)^{2}=f_{i}^{2} g_{j}^{2}>1
$$

a contradiction.
Proposition B.2. Let $i \in\{1,2,5\} \cap S$ and $j \in\{3,4,6,7\} \cap S$ be such that $N_{i} \cap S=$ ( $N_{j} \cap S$ ) $\dot{\cup}\{i\}$. Then

$$
f_{i}^{2}=\frac{\left(\alpha_{i}-2 \nu\right) \mu}{-F_{1}\left(-\alpha_{i}\right)}, \quad g_{i}^{2}=\frac{\left(\alpha_{i}-2 \mu\right) \nu}{-F_{1}\left(-\alpha_{i}\right)}
$$

and

$$
f_{j}=\left(1-\frac{\alpha_{i}}{\mu}\right) f_{i}, \quad g_{j}=\left(1-\frac{\alpha_{i}}{\nu}\right) g_{i}
$$

Moreover, $-F_{1}\left(-\alpha_{i}\right)$ is positive and $\alpha_{i}-2 \mu$ is negative.
Proof. The proof of Proposition B.1, slightly modified, gives the desired result.
Proposition B.3. Suppose $i, j, k \in S$ where $(i, j, k)$ is either $(2,3,4)$ or $(5,6,7)$. Then

$$
f_{k}=\frac{\mu f_{j}-\alpha_{i} f_{i}}{\mu}, \quad g_{k}=\frac{\nu g_{j}-\alpha_{i} g_{i}}{\nu}
$$

and

$$
\alpha_{i}=\frac{2 \mu^{2} v^{2} \alpha_{j}}{F_{2}\left(\alpha_{j}\right)}
$$

Proof. Using the eigenfunction equations for $f_{j}, f_{k}$ and for $g_{j}, g_{k}$, it follows that

$$
f_{k}=\frac{\mu f_{j}-\alpha_{i} f_{i}}{\mu}, \quad g_{k}=\frac{\nu g_{j}-\alpha_{i} g_{i}}{\nu}
$$

Combined with Lemma 3.7, it follows that

$$
\begin{aligned}
0 & =\mu f_{k}^{2}-\nu g_{k}^{2}-(\mu-\nu) \\
& =\mu\left(\frac{\mu f_{j}-\alpha_{i} f_{i}}{\mu}\right)^{2}-v\left(\frac{\nu g_{j}-\alpha_{i} g_{i}}{\nu}\right)^{2}-(\mu-\nu) .
\end{aligned}
$$

After expanding, we note that the right-hand side can be expressed purely in terms of $\mu, \nu, \alpha_{i}, f_{i}^{2}, f_{i} f_{j}, f_{j}^{2}, g_{i}^{2}, g_{i} g_{j}$, and $g_{j}^{2}$. Note that Proposition B. 1 gives explicit formulas for $f_{i}^{2}, f_{i} f_{j}$, and $f_{j}^{2}$, as well as $g_{i}^{2}, g_{i} g_{j}$, and $g_{j}^{2}$, purely in terms of $\mu, \nu$, and $\alpha_{j}$. With the help of a computer algebra system, we make these substitutions and factor the righthand side as:

$$
0=(\mu-\nu) \cdot \alpha_{i} \cdot \frac{2 \mu^{2} \nu^{2} \cdot \alpha_{j}-F_{2}\left(\alpha_{j}\right) \cdot \alpha_{i}}{\mu^{2} \nu^{2} \cdot F_{1}\left(\alpha_{i}\right)}
$$

Since $\alpha_{i},(\mu-\nu) \neq 0$, the desired claim holds.
Proposition B.4. Suppose $1, i, j, k \in S$ where $(i, j, k)$ is either $(2,3,4)$ or $(5,6,7)$. Then

$$
f_{1}=\frac{\mu f_{i}+\alpha_{k} f_{k}}{\mu}, \quad g_{1}=\frac{\nu g_{i}+\alpha_{k} g_{k}}{\nu}
$$

and

$$
\alpha_{k}=\frac{\alpha_{j} \cdot F_{2}\left(\alpha_{j}\right)^{2}}{F_{3}\left(\alpha_{j}\right)}
$$

Proof. Using the eigenfunction equations for $f_{1}, f_{i}, f_{j}, f_{k}$ and for $g_{1}, g_{i}, g_{j}, g_{k}$, it follows that

$$
f_{1}=\frac{\mu f_{i}+\alpha_{k} f_{k}}{\mu}, \quad g_{1}=\frac{\nu g_{i}+\alpha_{k} g_{k}}{\nu}
$$

and

$$
f_{k}=\frac{\mu f_{j}-\alpha_{i} f_{i}}{\mu}, \quad g_{k}=\frac{\nu g_{j}-\alpha_{i} g_{i}}{\nu}
$$

Altogether,

$$
f_{1}=\frac{\mu^{2} f_{i}+\alpha_{k}\left(\mu f_{j}-\alpha_{i} f_{i}\right)}{\mu^{2}}, \quad g_{1} \quad=\frac{\nu^{2} g_{i}+\alpha_{k}\left(\nu g_{j}-\alpha_{i} g_{i}\right)}{\nu^{2}}
$$

Combined with Lemma 3.6, it follows that

$$
\begin{aligned}
0 & =\mu f_{1}^{2}-\nu g_{1}^{2}-(\mu-\nu) \\
& =\mu\left(\frac{\mu^{2} f_{i}+\alpha_{k}\left(\mu f_{j}-\alpha_{i} f_{i}\right)}{\mu^{2}}\right)^{2}-\nu\left(\frac{\nu^{2} g_{i}+\alpha_{k}\left(\nu g_{j}-\alpha_{i} g_{i}\right)}{\nu^{2}}\right)^{2}-(\mu-\nu)
\end{aligned}
$$

After expanding, we note that the right-hand side can be expressed purely in terms of $\mu, \nu, f_{i}^{2}, f_{i} f_{j}, f_{j}^{2}, g_{i}^{2}, g_{i} g_{j}$, and $\alpha_{i}$. Note that Proposition B.1 gives explicit formulas for
$f_{i}^{2}, f_{i} f_{j}, f_{j}^{2}, g_{i}^{2}, g_{i} g_{j}$, and $g_{j}^{2}$ purely in terms of $\mu, \nu$, and $\alpha_{j}$. With the help of a computer algebra system, we make these substitutions and factor the right-hand side as:

$$
0=2 \alpha_{k} \cdot(\mu-\nu) \cdot \frac{\alpha_{j} \cdot F_{2}\left(\alpha_{j}\right)^{2}-\alpha_{k} \cdot F_{3}\left(\alpha_{j}\right)}{F_{1}\left(\alpha_{j}\right) \cdot F_{2}\left(\alpha_{j}\right)^{2}}
$$

So the desired claim holds.

Proposition B.5. Suppose $1, i, k \in S$ and $j \notin S$ where $(i, j, k)$ is either $(2,3,4)$ or $(5,6,7)$. Then,

$$
f_{1}=\frac{\mu f_{i}+\alpha_{k} f_{k}}{\mu}, \quad g_{1}=\frac{\nu g_{i}+\alpha_{k} g_{k}}{\nu}
$$

and

$$
\alpha_{k}=\frac{2 \alpha_{i} \mu^{2} v^{2}}{F_{2}\left(-\alpha_{i}\right)}
$$

Proof. Using the eigenfunction equations for $f_{1}, f_{i}, f_{j}, f_{k}$ and for $g_{1}, g_{i}, g_{j}, g_{k}$, it follows that

$$
f_{1}=\frac{\mu f_{i}+\alpha_{k} f_{k}}{\mu}, \quad g_{1}=\frac{\nu g_{i}+\alpha_{k} g_{k}}{\nu}
$$

and

$$
f_{k}=\frac{\mu f_{i}-\alpha_{i} f_{i}}{\mu}, \quad g_{k}=\frac{\nu g_{i}-\alpha_{i} g_{i}}{\nu}
$$

Altogether,

$$
f_{1}=\frac{\mu^{2} f_{i}+\alpha_{k}\left(\mu f_{i}-\alpha_{i} f_{i}\right)}{\mu^{2}}, \quad g_{1}=\frac{\nu^{2} g_{i}+\alpha_{k}\left(\nu g_{i}-\alpha_{i} g_{i}\right)}{\nu^{2}}
$$

Combined with Lemma 3.6, it follows that

$$
\begin{aligned}
0 & =\mu f_{1}^{2}-\nu g_{1}^{2}-(\mu-v) \\
& =\mu\left(\frac{\mu^{2} f_{i}+\alpha_{k}\left(\mu f_{i}-\alpha_{i} f_{i}\right)}{\mu^{2}}\right)^{2}-\nu\left(\frac{\nu^{2} g_{i}+\alpha_{k}\left(\nu g_{i}-\alpha_{i} f_{i}\right)}{\nu^{2}}\right)^{2}-(\mu-\nu)
\end{aligned}
$$

After expanding, we note that the right-hand side can be expressed purely in terms of $\mu, \nu, f_{i}^{2}, f_{i} f_{j}, f_{j}^{2}, g_{i}^{2}, g_{i} g_{j}$, and $\alpha_{i}$. Note that Proposition B. 1 gives explicit formulas for $f_{i}^{2}, f_{i} f_{j}, f_{j}^{2}, g_{i}^{2}, g_{i} g_{j}$, and $g_{j}^{2}$ purely in terms of $\mu, \nu$, and $\alpha_{j}$. With the help of a computer algebra system, we make these substitutions and factor the right-hand side as:

$$
0=2 \alpha_{k} \cdot(\mu-\nu) \cdot \frac{\alpha_{j} \cdot F_{2}\left(\alpha_{j}\right)^{2}-\alpha_{k} \cdot F_{3}\left(\alpha_{j}\right)}{F_{1}\left(\alpha_{j}\right) \cdot F_{2}\left(\alpha_{j}\right)^{2}}
$$

So the desired claim holds.
Proposition B.6. Suppose $1 \notin S$ and $i, j, k, \ell \in S$ where $(i, j, k, \ell)$ is either $(2,3,4,7)$ or (5, 6, 7, 4). Then

$$
\alpha_{k}=\frac{F_{4}(x)}{F_{3}(x)}
$$

Proof. Using the eigenfunction equations for $f_{\ell}, f_{i}, f_{j}, f_{k}$ and for $g_{\ell}, g_{i}, g_{j}, g_{k}$, it follows that

$$
f_{\ell}=\frac{\alpha_{i} f_{i}+\alpha_{j} f_{j}+\alpha_{k} f_{k}}{\mu}, \quad g_{1}=\frac{\alpha_{i} g_{i}+\alpha_{j} g_{j}+\alpha_{k} g_{k}}{\nu}
$$

and

$$
f_{k}=\frac{\mu f_{j}-\alpha_{i} f_{i}}{\mu}, \quad g_{k}=\frac{\nu g_{j}-\alpha_{i} g_{i}}{\nu}
$$

Altogether,

$$
f_{\ell}=\frac{\mu \alpha_{i} f_{i}+\alpha_{j} f_{j}+\alpha_{k}\left(\mu f_{j}-\alpha_{i} f_{i}\right)}{\mu^{2}}, \quad g_{\ell} \quad=\frac{\nu \alpha_{i} g_{i}+\alpha_{j} g_{j}+\alpha_{k}\left(\nu g_{j}-\alpha_{i} g_{i}\right)}{\nu^{2}}
$$

Combined with Lemma 3.7, it follows that

$$
\begin{aligned}
0= & \mu f_{\ell}^{2}-\nu g_{\ell}^{2}-(\mu-\nu) \\
= & \mu\left(\frac{\mu \alpha_{i} f_{i}+\alpha_{j} f_{j}+\alpha_{k}\left(\mu f_{j}-\alpha_{i} f_{i}\right)}{\mu^{2}}\right)^{2}-\nu\left(\frac{\nu \alpha_{i} g_{i}+\alpha_{j} g_{j}+\alpha_{k}\left(\nu g_{j}-\alpha_{i} g_{i}\right)}{\nu^{2}}\right)^{2} \\
& -(\mu-\nu) .
\end{aligned}
$$

After expanding, we note that the right-hand side can be expressed purely in terms of $\mu, \nu, f_{i}^{2}, f_{i} f_{j}, f_{j}^{2}, g_{i}^{2}, g_{i} g_{j}$, and $\alpha_{i}$. Note that Proposition B. 1 gives explicit formulas for $f_{i}^{2}, f_{i} f_{j}, f_{j}^{2}, g_{i}^{2}, g_{i} g_{j}, g_{j}^{2}, \alpha_{i}, \alpha_{j}, \alpha_{k}$ purely in terms of $\mu, \nu$, and $\alpha_{j}$. With the help of a computer algebra system, we make these substitutions and factor the right-hand side as:

$$
0=2(\mu-\nu) \cdot \alpha_{k} \cdot \frac{F_{4}\left(\alpha_{j}\right)-\alpha_{k} \cdot F_{3}\left(\alpha_{j}\right)}{F_{1}\left(\alpha_{j}\right) \cdot F_{2}\left(\alpha_{j}\right)^{2}} .
$$

Proposition B.7. Suppose 2,4,5,7 $\in S$ and let $\alpha_{\neq 4,7}:=\sum_{\substack{i \in S, i \neq 4,7}} \alpha_{i}$. Then

$$
\begin{aligned}
& \alpha_{4}=\frac{\left(1-\alpha_{\neq 4,7}\right) f_{7}-\mu\left(f_{2}-f_{7}\right)}{f_{4}+f_{7}} \\
& \alpha_{7}=\frac{\left(1-\alpha_{\neq 4,7}\right) f_{4}-\mu\left(f_{5}-f_{2}\right)}{f_{4}+f_{7}}
\end{aligned}
$$

and

$$
\begin{aligned}
& \alpha_{4}=\frac{\left(\left(1-\alpha_{\neq 4,7}\right) g_{7}-v\left(g_{2}-g_{7}\right)\right.}{g_{4}+g_{7}}, \\
& \alpha_{7}=\frac{\left(1-\alpha_{\neq 4,7}\right) g_{4}-v\left(g_{5}-g_{2}\right)}{g_{4}+g_{7}} .
\end{aligned}
$$

Proof. Taking the difference of the eigenvector equations for $f_{2}$ and $f_{5}$, and for $g_{2}$ and $g_{5}$, we have

$$
\alpha_{7} f_{7}-\alpha_{4} f_{4}=\mu\left(f_{2}-f_{5}\right), \quad \alpha_{7} g_{7}-\alpha_{4} g_{4}=\nu\left(g_{2}-g_{5}\right)
$$

Combining these equalities with the equation $\alpha_{4}+\alpha_{7}=1-\alpha_{\neq 4,7}$ completes the proof.
B.2. Algorithm. In this subsection, we briefly detail how the computer-assisted proof of Lemma 5.2 works. This proof is via interval arithmetic, and, at a high level, consists largely of iteratively decomposing the domain of feasible choices of $\left(\alpha_{3}, \alpha_{6}, \mu, \nu\right)$ for a given $S$ into smaller subregions (boxes) until all subregions violate some required equality or inequality. We provide two similar, but slightly different computer assisted proofs of this fact, and both of these can be found at the spread_numeric GitHub repository [27]. The first, found in folder interval1, is a shorter and simpler version, containing slightly fewer formulas, albeit at the cost of overall computation and run time. The second, found in the folder interval2, contains slightly more formulas and makes a greater attempt to optimize computation and run time. Below, we further detail the exact output and run time of both versions (exact output can be found in [27]), but for now, we focus on the main aspects of both proofs, and consider both together, saving a more detailed discussion of the differences for later.

These algorithms are implemented in Python using the PyInterval package. The algorithms consist of two parts: a main file containing useful formulas and subroutines and 17 different files used to rule out each of the 17 cases for $S$. The main file, casework_helper, contains functions with the formulas of Subsection B.1 (suitably modified to limit error growth), and functions used to check that certain equalities and inequalities are satisfied. In particular, casework_helper contains formulas for

- $\alpha_{2}$, assuming $\{2,3,4\} \subset S$ (using Proposition B.3)
- $\alpha_{4}$, assuming $\{1,2,3,4\} \subset S$ (using Proposition B.4)
- $\alpha_{4}$, assuming $\{2,3,4,7\} \subset S, 1 \notin S$ (using Proposition B.6)
- $\alpha_{4}$, assuming $\{1,2,4\} \subset S, 3 \notin S$ (using Proposition B.5)
- $f_{3}$ and $g_{3}$, assuming $\{2,3\} \subset S$ (using Proposition B.1)
- $f_{2}$ and $g_{2}$, assuming $\{2,3\} \subset S$ (using Proposition B.1)
- $f_{4}$ and $g_{4}$, assuming $\{2,3,4\} \subset S$ (using Proposition B.3)
- $f_{1}$ and $g_{1}$, assuming $\{1,2,4\} \subset S$ (using Propositions B. 4 and B.5)
- $f_{2}$ and $g_{2}$, assuming $\{2,4\} \subset S, 3 \notin S$ (using Proposition B.2)
- $f_{4}$ and $g_{4}$, assuming $\{2,4\} \subset S, 3 \notin S$ (using Proposition B.2)
as a function of $\alpha_{3}, \mu$, and $\nu$ (and $\alpha_{2}$ and $\alpha_{4}$, which can be computed as functions of $\alpha_{3}$, $\mu$, and $\nu$ ). Some of the formulas are slightly modified compared to their counterparts in this Appendix, for the purpose of minimizing accumulated error. Each formula is performed using interval arithmetic, while restricting the resulting interval solution to the correct range. In addition, we recall that we have the inequalities
- $\alpha_{i} \in[0,1]$, for $i \in S$
- $\left|g_{2}\right|,\left|f_{3}\right| \leq 1,\left|f_{2}\right|,\left|g_{3}\right| \geq 1$, for $\{2,3\} \subset S$
- $\left|f_{4}\right| \leq 1,\left|g_{4}\right| \geq 1$, for $4 \in S$
- $\left|f_{1}\right| \geq 1,\left|g_{1}\right| \leq 1$, for $\{1,2,4\} \in S$
- $\left|f_{4}\right|,\left|g_{2}\right| \leq 1,\left|f_{2}\right|,\left|g_{4}\right| \geq 1$, for $\{2,4\} \in S, 3 \notin S$
- $\alpha_{3}+2 \nu \leq 0$, for $\{2,3\} \in S$ (using Proposition B.1)
- $\alpha_{2}-2 \mu \leq 0$, for $\{2,4\} \in S, 3 \notin S$ (using Proposition B.2).

These inequalities are also used at various points in the algorithms. This completes a brief overview of the casework_helper file. Next, we consider the different files used to test feasibility for a specific choice of $S \subset\{1, \ldots, 7\}$, each denoted by case\{elements of $S\}$, i.e., for $S=\{1,4,5,7\}$, the associated file is case1457. For each specific case, there are
a number of different properties which can be checked, including eigenvector equations, bounds on edge density, norm equations for the eigenvectors, and the ellipse equations. Each of these properties has an associated function which returns FALSE if the property cannot be satisfied, given the intervals for each variable, and returns TRUE otherwise. The implementation of each of these properties is rather intuitive, and we refer the reader to the programs themselves (which contain comments) for exact details [27]. Each feasibility file consists of two parts. The first part is a function is_feasible (mu,nu,a3,a6) that, given bounding intervals for $\mu, \nu, \alpha_{3}, \alpha_{6}$, computes intervals for all other variables (using interval arithmetic) and checks feasibility using the functions in the casework_helper file. If any checked equation or inequality in the file is proven to be unsatisfiable (i.e., see Example 5.3), then this function outputs 'FALSE', otherwise the function outputs 'TRUE' by default. The second part is a divide and conquer algorithm that breaks the hypercube

$$
\left(\mu, \nu, \alpha_{3}, \alpha_{6}\right) \in[.65,1] \times[-.5,-.15] \times[0,1] \times[0,1]
$$

into sub-boxes of size $1 / 20$ by $1 / 20$ by $1 / 10$ by $1 / 10$, checks feasibility in each box using is_feasible, and subdivides any box that does not rule out feasibility (i.e., subdivides any box that returns 'TRUE'). This subdivision breaks a single box into two boxes of equal size, by subdividing along one of the four variables. The variable used for this subdivision is chosen iteratively, in the order $\alpha_{3}, \alpha_{6}, \mu, \nu, \alpha_{3}, \ldots$. The entire divide and conquer algorithm terminates after all sub-boxes, and therefore, the entire domain

$$
\left(\mu, v, \alpha_{3}, \alpha_{6}\right) \in[.65,1] \times[-.5,-.15] \times[0,1] \times[0,1]
$$

has been shown to be infeasible, at which point the algorithm prints 'infeasible'. Alternatively, if the number of subdivisions reaches some threshold, then the algorithm terminates and outputs 'feasible'.

Next, we briefly detail the output of the algorithms casework_helper/intervals1 and casework_helper/intervals2. Both algorithms ruled out 15 of the 17 choices for $S$ using a maximum depth of 26 , and failed to rule out cases $S=\{4,5,7\}$ and $S=\{1,7\}$ up to depth 51 . For the remaining 15 cases, intervals 1 considered a total of 5.5 million boxes, was run serially on a personal computer, and terminated in slightly over twelve hours. For these same 15 cases, intervals2 considered a total of 1.3 million boxes, was run in parallel using the Penn State math department's 'mathcalc' computer, and terminated in under 140 minutes. The exact output for both versions of the spread_numeric algorithm can be found in [27].

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[^0]:    ${ }^{1}$ Here, we take the usual convention that for all $x \in(-\infty, \infty], \infty+x=x+\infty=\infty$.

