

A finite element method for quantum graphs

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We study the numerical solution of boundary and initial value problems for differential equations posed on graphs or networks. The graphs of interest are *quantum graphs*, i.e., metric graphs endowed with a differential operator acting on functions defined on the graph's edges with suitable side conditions. We describe and analyse the use of linear finite elements to discretize the spatial derivatives for a class of linear elliptic model problems. The solution of the discrete equations is discussed in detail in the context of a (nonoverlapping) domain decomposition approach. For model elliptic problems and a wide class of graphs, we show that a combination of Schur complement reduction and diagonally preconditioned conjugate gradients results in optimal complexity. For problems of parabolic type, we consider the use of exponential integrators based on Krylov subspace methods. Numerical results are given for both simple and complex graph topologies.

Keywords: quantum graphs; finite element method; complex graphs; sparse matrices; iterative methods; diffusion on graphs; exponential integrators.

1. Introduction

Differential equations posed on networks, or graphs, play an important role in the modeling and simulation of a wide variety of complex phenomena (see [Newman, 2010](#)). An important example is the study of diffusion phenomena on networks, which in the simplest cases reduces to the solution of initial value problems for linear systems of ordinary differential equations (ODEs). Although these models may be adequate for studying simple situations, such as those where the relations between the constituents of the system (corresponding to graph vertices) can be modeled by a simple binary relation (connected or not connected), more sophisticated models are necessary when dealing with more complex situations.

Metric and *quantum graphs* provide useful models for a variety of physical systems including conjugated molecules, quantum wires, photonic crystals, thin waveguides, carbon nanostructures, and so forth. We refer to [Berkolaiko & Kuchment \(2013\)](#) for details and references (see also [Lagnese *et al.* \(1994\)](#) for related problems not considered in [Berkolaiko & Kuchment \(2013\)](#)). A metric graph is a graph in which each edge is endowed with an implicit metric structure. Often (but not always), its edges can be identified with intervals on the real line. In technical terms, a metric graph is an example of one-dimensional topological manifold, or one-dimensional simplicial complex. A quantum graph is a metric graph equipped with a differential operator ('Hamiltonian') and suitable vertex conditions (see the next section for more precise definitions). This differential operator acts on functions defined on the edges and vertices of the underlying metric graph.

In physics and engineering applications, there is a strong interest in the spectral properties of the Hamiltonian, as well as on associated wave propagation, transport and diffusion phenomena. The literature on quantum graphs is vast; most papers deal with theoretical issues such as spectral theory, well posedness and so forth, or with physical applications. On the other hand, the literature devoted to computational issues is almost nonexistent, apart from a handful of references dealing with rather special situations (see, e.g., [Pesenson, 2005](#); [Nordenfelt, 2007](#); [Herty *et al.*, 2010](#); [Hild & Leugering, 2012](#); [Wybo *et al.*, 2015](#); [Zlotnik *et al.*, 2015](#)), and the numerical aspects are typically not the main focus.

Here, we take a first step in the systematic study of numerical methods for the solution of differential problems involving quantum graphs. We discuss a simple spatial discretization using linear finite elements and techniques for the fast solution of the discrete equations, using simple elliptic and parabolic model problems to illustrate our approach. We focus on this type of discretization because it allows us to highlight interesting relations connecting the discretized Hamiltonian with the graph Laplacian of the underlying combinatorial graph. In addition to equations posed on highly structured graphs, which are of interest in physics, we also consider the case of *complex* graphs with nontrivial topologies, in view of potential applications in fields such as physiology, biology and engineering. Not surprisingly, we observe significant differences with the numerical solution of PDEs posed on more typical spatial domains.

The remainder of the article is organized as follows. The necessary background information on metric and quantum graphs is provided in Section 2. In Section 3, we introduce a simple linear finite element method and analyse its convergence for a model quantum graph. The useful notion of *extended graph* is introduced in Section 4; in this section, we also give a detailed description of the matrices obtained from the finite element method. The conditioning of the mass matrix is studied in Section 5. Section 6 contains some remarks on the spectrum of discretized quantum graphs. Solution algorithms from the discretized equations, including preconditioned conjugate gradient (PCG) methods and schemes for integrating time-dependent problems, are discussed in Section 7. Numerical experiments for some simple elliptic and parabolic model problems are presented in Section 8. Finally, in Section 9, we present our conclusions and a list of topics for future work.

2. Definitions and notations

We give in sequence the definitions, the notations and the assumptions that we will use in the following. We refer to [Berkolaiko & Kuchment \(2013\)](#) for a comprehensive introduction to the theory of quantum graphs.

DEFINITION 2.1 A *combinatorial graph* $\Gamma = (\mathcal{V}, \mathcal{E})$ is a collection of a finite number of vertices and of edges connecting pairs of distinct vertices. We will denote by $\mathcal{V} = \{v_1, \dots, v_N\}$ the set of vertices and by $\mathcal{E} = \{e_j = (v_i, v_k)\}_{j=1, \dots, M}$ the set of edges. Thus, an edge can be identified with a pair of vertices. The graph is *undirected* if no orientation is assigned to the edges; in this case, we do not distinguish between (v_i, v_j) and (v_j, v_i) . Otherwise, the graph is *directed*. For an undirected graph, we define for each vertex v_i its *degree* d_{v_i} as the number of edges e_k such that $e_k = (v_i, v_j)$. Since only a single edge (at most) is allowed between any two vertices, d_{v_i} is the number of vertices adjacent to v_i (i.e., the number of ‘immediate neighbors’ of v_i in Γ). We restrict our attention to graphs with no self-loops: $(v_i, v_i) \notin \mathcal{E}$, for all i . A graph is *connected* if from each vertex v_i in \mathcal{V} there exists a path $(v_i, v_k), (v_k, v_j), \dots$ made by edges in \mathcal{E} connecting it to any of the other vertices.

In this work, we only consider sparse graphs, which we define as those graphs with $M = O(N)$. We assume that the initial graph Γ is undirected, but we assign an arbitrarily chosen direction to each

edge to define an *incidence matrix* of Γ . This is a matrix $\mathbf{E} \in \mathbb{R}^{N \times M}$ where each column corresponds to an edge and has only two nonzero entries corresponding to the two vertices identifying the edge. We will arbitrarily fix the first nonzero entry in the column to the value 1 and the second nonzero entry to the value -1 (this is equivalent to assigning an orientation to the edges). We emphasize that the choice of the signs is irrelevant for the purposes of this exposition. The advantage of using \mathbf{E} will become more clear after the introduction of the concept of quantum graph. It is important to remark that \mathbf{E} has an interesting interpretation as a finite-dimensional operator mimicking the divergence operator on differentiable functions (see [Arioli & Manzini, 2003, 2006](#) for a similar discussion related to mixed finite element problems). Finally, we recall that \mathbf{E} has rank $N - 1$ in the case of a connected graph.

The matrix \mathbf{E}^T is also interpretable as the finite-dimensional equivalent of the gradient operator acting on differentiable functions. It is also important to observe that the classical (combinatorial) *graph Laplacian* coincides with the matrix $\mathbf{L}_\Gamma = \mathbf{E}\mathbf{E}^T$. Note that \mathbf{L}_Γ does not depend on the choice of orientation used for the edges, since $\mathbf{E}\mathbf{Q}(\mathbf{E}\mathbf{Q})^T = \mathbf{E}\mathbf{E}^T$ for any $M \times M$ diagonal matrix \mathbf{Q} with entries ± 1 on the main diagonal. The graph Laplacian is a symmetric positive semidefinite M -matrix. The multiplicity of 0 as an eigenvalue of \mathbf{L}_Γ equals the number of connected components of Γ ; if Γ is connected, the null space of \mathbf{L}_Γ is one-dimensional and is spanned by the vector of all ones.

Let the matrix \mathbf{D} be the diagonal of the matrix \mathbf{L}_Γ . The diagonal entries of \mathbf{D} are just the degrees of the corresponding vertices. The matrix $\mathbf{A}\mathbf{d} = \mathbf{D} - \mathbf{L}_\Gamma$ is the *adjacency matrix* of the graph, i.e., the matrix where the (i, j) entry is either 1 or 0 according to whether $(v_i, v_j) \in \mathcal{E}$.

DEFINITION 2.2 A connected graph Γ is said to be a *metric graph* if,

1. to each edge e is assigned a length ℓ_e such that $0 < \ell_e < \infty$, and
2. each edge is assigned a coordinate $x^e \in [0, \ell_e]$, which increases in a specified (but otherwise arbitrary) direction along the edge.

In general, the edges could be simple differentiable curves (i.e., no loops). However, to simplify the notation, we assume that each edge is a straight line joining the two vertices defining the edge. The direction used to assign coordinates to points on a given edge will be the same one used to define the incidence matrix. In our definition of a metric graph, we assume that all lengths are finite. We refer to [Berkolaiko & Kuchment \(2013\)](#) for discussions of infinite metric graphs with some edges having infinite length. We define the *volume* of a metric graph Γ as $\mathbf{vol}(\Gamma) = \sum_{e \in \mathcal{E}} \ell_e$. In the following, we will discuss only finite graphs with all edge lengths finite, hence with $\mathbf{vol}(\Gamma) < \infty$. Note that we do not assume that metric graphs are embedded in \mathbb{R}^n for some n . With the structure described above, the metric graph Γ becomes a one-dimensional domain, where for each edge e we have a variable x^e representing locally the global coordinate variable x .

As noted earlier, a sequence of contiguous vertices defines a path in Γ formed by $\{e_j\}_{j=1}^{\hat{M}}$ and the associated *path length* is simply $\sum \ell_{e_j}$. We define the *distance* $d(v_i, v_j)$ between two vertices v_i and v_j as the length of a shortest path in Γ between them. This notion of distance can be extended in a natural way to define the distance between any two points (possibly lying on different edges) in the one-dimensional simplicial complex. Endowed with this distance, Γ is readily seen to be a metric space.

Next, we proceed to introduce function spaces on a metric graph Γ and linear differential operators defined on these spaces.

DEFINITION 2.3 The space $L^2(\Gamma) = \bigoplus_e L^2(e)$ is the space of all square-integrable measurable functions u on the edges of Γ ; i.e.,

$$\|u\|_{L^2(\Gamma)}^2 = \sum_{e \in \mathcal{E}} \|u|_e\|_{L^2(e)}^2 < \infty.$$

The inner product in this space will be denoted by $(u, v)_{L^2(\Gamma)} = \int_{\Gamma} u(x)v(x) dx$.

The Sobolev space $H^1(\Gamma) = \bigoplus_e H^1(e) \cap C^0(\Gamma)$ is the space of all continuous functions u on Γ , $u \in C^0(\Gamma)$, such that $u|_e$ belongs to $H^1(e)$ for each edge e , i.e.,

$$\|u\|_{H^1(\Gamma)}^2 = \sum_{e \in \mathcal{E}} \|u|_e\|_{H^1(e)}^2 < \infty.$$

We remind that

$$\|u\|_{H^1(e)}^2 = \int_e \left(\frac{du}{dx} \right)^2 dx + \int_e u^2 dx$$

and therefore

$$\|u\|_{H^1(\Gamma)}^2 = \int_{\Gamma} \left(\frac{du}{dx} \right)^2 dx + \int_{\Gamma} u^2 dx.$$

Because of the properties of the Sobolev spaces of functions of one variable, the functions belonging to $H^1(e)$ are continuous (Brezis, 2010, Chapter 8). This justifies the assumption in Definition 2.3 of restricting membership in $H^1(\Gamma)$ to the continuous functions. Moreover, the global continuity assumption implies automatically that the functions on all edges adjacent to a vertex v assume the same value at v .

The operators that we consider here are quite simple, but suitable for describing interesting dynamics on metric graphs. More specifically, besides the continuity of the functions on Γ , we will use ‘Neumann–Kirchhoff’ conditions on the vertices, i.e., denoting by \mathcal{E}_v the subset of \mathcal{E} comprising the edges incident to the vertex v , we impose the condition

$$\sum_{e \in \mathcal{E}_v} \frac{du}{dx}(v) = 0 \quad \forall v \in \mathcal{V}. \quad (2.1)$$

This corresponds to assuming Neumann conditions at all vertices. Strictly speaking, (2.1) requires that $u|_e \in H^2(e)$ for each edge $e \in \mathcal{E}$. As discussed below (see Theorem 3.2), this condition is usually satisfied in our setting. Conditions (2.1) express the *conservation of currents* if the metric graph Γ is viewed as an electrical network, hence the name. Note that to give a meaning to (2.1), we need to assume that the derivatives are taken in the directions away from the vertex, which we call the outgoing directions (see Berkolaiko & Kuchment, 2013). Moreover, we observe that (2.1) are the natural boundary conditions satisfied by the following one-dimensional Schrödinger-type operator:

$$\mathcal{H} : u(x) \mapsto -\frac{d^2u}{dx^2} + v(x)u(x). \quad (2.2)$$

The function $v(x)$ in (2.2) is a *potential*; throughout the article, we assume that $v \in L^\infty(\Gamma)$.

The operator (2.2) is defined for functions $u \in L^2(\Gamma)$ such that $u|_e \in H^2(e)$ for all $e \in \mathcal{E}$. We observe that $u|_e \in H^2(e) \Rightarrow u|_e \in C^1(\bar{e})$, thus, because of the regularity property of $u|_e$, the Neumann–Kirchhoff conditions (2.1) are well defined in the classical sense. However, for our purposes, it is convenient to introduce a weak form of (2.2), which requires that $u \in H^1(\Gamma)$. Here, we follow the treatment in Berkolaiko & Kuchment (2013, p. 25). From (2.1) and (2.2), we have that the bilinear form \mathfrak{h} corresponding to the Hamiltonian \mathcal{H} is

$$\mathfrak{h} : H^1(\Gamma) \times H^1(\Gamma) \longrightarrow \mathbb{R}, \quad \mathfrak{h}(u, g) = \sum_{e \in \mathcal{E}} \left\{ \int_e \frac{du}{dx} \frac{dg}{dx} dx + \int_e u(x)g(x)v(x) dx \right\}. \quad (2.3)$$

The corresponding energy functional is given by

$$J[u] = \frac{1}{2} \sum_{e \in \mathcal{E}} \left\{ \int_e \left(\frac{du}{dx} \right)^2 + u(x)^2 v(x) dx \right\} \quad \forall u \in H^1(\Gamma). \quad (2.4)$$

Furthermore, if $v(x) \geq v_0$ for some constant $v_0 > 0$, then the symmetric bilinear form (2.3) is continuous and coercive on $H^1(\Gamma)$, with coercivity constant $\gamma_0 = \min(1, v_0)$ and continuity constant $\gamma_1 = \max(1, \|v\|_\infty)$, thus $J[\cdot]$ is strictly convex and has a unique minimum in $H^1(\Gamma)$ (see Theorem 3.2 below). Moreover, for $f \in L^2(\Gamma)$, the Euler equation for $J[u]$:

$$\mathfrak{h}(u, g) = \int_\Gamma f(x)g(x) dx \quad \forall g \in H^1(\Gamma) \quad (2.5)$$

has a unique solution (Lax–Milgram Theorem).

We observe that even though here we are taking into account only Neumann conditions, it is possible (if required) to fix the values of the functions on a subset of the vertices that will become the Dirichlet boundary vertices. These will play the same role of Dirichlet boundary conditions in the classical sense. More general conditions at the vertices can be imposed and we refer to Berkolaiko & Kuchment (2013) for a deeper discussion.

Hereafter, we define a quantum graph as follows:

DEFINITION 2.4 A *quantum graph* is a metric graph equipped with the Hamiltonian operator \mathcal{H} defined by the operator (2.2) subject to the conditions (2.1) at the vertices.

Although this definition is more restrictive than the one found, e.g., in Berkolaiko & Kuchment (2013), it is adequate for our purposes.

Finally, among our goals is the analysis of parabolic problems on metric graphs. In this case, we assume that the functions we use also depend on a second variable t representing time (see Raviart *et al.*, 1983), i.e.,

$$u(x, t) : \Gamma \times [0, T] \longrightarrow \mathbb{R},$$

and that they are elements of suitable Bochner spaces, specified below.

DEFINITION 2.5 Let V denote either $L^2(\Gamma)$ or $H^1(\Gamma)$. Let $C^0([0, T]; V)$ be the space of functions $u(x, t)$ that are continuous in t with values in V , i.e., for each fixed value t^* of t we have that $u(\cdot, t^*) \in V$. This space is equipped with the norm

$$\|u\|_{C^0([0, T]; V)} = \sup_{0 \leq t \leq T} \|u(\cdot, t)\|_V.$$

Let $L^2([0, T]; V)$ be the space of functions $u(x, t)$ that are square integrable in t for the dt measure with values in V , i.e., for each fixed value t^* of t we have that $u(\cdot, t) \in V$. This space is equipped with the norm

$$\|u\|_{L^2([0, T]; V)} = \left(\int_0^T \|u(\cdot, t)\|_V^2 dt \right)^{\frac{1}{2}}$$

and scalar product

$$(u, g)_{L^2([0, T]; V)} = \int_0^T (u(\cdot, t), g(\cdot, t))_V dt.$$

We note that all these definitions can be easily modified to deal with self-adjoint operators acting on spaces of complex-valued functions, as required, e.g., in quantum–mechanical applications.

3. Finite element approximation of quantum graphs

On each edge of the quantum graph, it is possible to use the classical one-dimensional finite element method. Let e be a generic edge identified by two vertices, which we denote by v_a and v_b . The first step is to subdivide the edge in n_e intervals of length h_e , with $n_e \geq 2$. The points

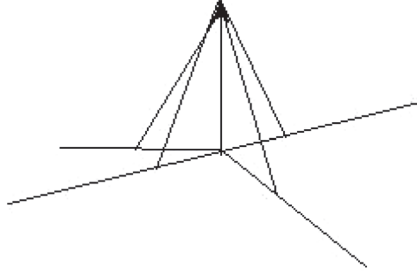
$$\left\{ \mathcal{V}_e = \{x_j^e\}_{j=1}^{n_e-1} \right\} \cup \{v_a^e\} \cup \{v_b^e\}$$

form a chain linking vertex v_a^e to vertex v_b^e lying on the edge e . The internal points x_j^e are said to be the *nodes* of the discretization. For each of the internal nodes, we denote by $\{\psi_j^e\}_{j=1}^{n_e-1}$ the standard hat basis functions

$$\psi_j^e(x^e) = \begin{cases} 1 - \frac{|x_j^e - x^e|}{h_e} & \text{if } x_{j-1}^e \leq x^e \leq x_{j+1}^e \\ 0 & \text{otherwise,} \end{cases} \quad (3.1)$$

where we have set $x_0^e = v_a^e$ and $x_{n_e}^e = v_b^e$. We also define the *neighboring set* \mathcal{W}_v of a vertex $v \in \mathcal{V}$ as the union of all the sets $[v_a^e, x_1^e]$ and $[x_{n_e}^e, v_b^e]$:

$$\mathcal{W}_v = \left\{ \bigcup_{e \in \{e \in \mathcal{E}_v \text{ s.t. } v_a^e = v\}} [v, x_1^e] \right\} \cup \left\{ \bigcup_{e \in \{e \in \mathcal{E}_v \text{ s.t. } v_b^e = v\}} [x_{n_e-1}^e, v] \right\}.$$


 FIG. 1. Example of ψ_j at a vertex v .

Note that \mathscr{W}_v is itself a (star shaped) metric graph because, as a sub-graph, it inherits the metric properties of the original graph. We introduce the functions $\phi_v(x)$ with support $\text{supp}(\phi_v(x)) = \mathscr{W}_v$ such that

$$\phi_v(x)|_{\mathscr{W}_v \cap e} = \begin{cases} 1 - \frac{|x_v^e - x^e|}{h_e} & \text{if } x^e \in \mathscr{W}_v \cap e; e \in \mathcal{E}_v \\ 0 & \text{otherwise,} \end{cases} \quad (3.2)$$

where x_v^e is either 0 or ℓ_e depending on the direction on the edge and its parametrization. Fig. 1 describes a simple example. The functions ψ_j^e are a basis for the finite-dimensional space

$$V_h^e = \left\{ w \in H_0^1(e); w|_{[x_j^e, x_{j+1}^e]} \in P_1 \quad j = 0, 1, \dots, n_e - 1 \right\},$$

where P_1 is the space of linear functions. Globally, we construct the space

$$V_h(\Gamma) = \left(\bigoplus_{e \in \mathcal{E}} V_h^e \right) \oplus \text{span} \{ \phi_v \}_{v \in \mathcal{V}}$$

and we have

$$V_h(\Gamma) \subset C^0(\Gamma).$$

This is a finite-dimensional space of functions that belong to $H^1(\Gamma)$.

Any function $w_h \in V_h(\Gamma)$ is then a linear combination of the ψ_j^e and ϕ_v :

$$w_h(x) = \sum_{e \in \mathcal{E}} \sum_{j=1}^{n_e-1} \alpha_j^e \psi_j^e(x) + \sum_{v \in \mathcal{V}} \beta_v \phi_v(x) \quad x \in \Gamma.$$

When we approximate the variational equation

$$\mathfrak{h}(u, g) = \int_{\Gamma} fg \, dx \quad \forall g \in H^1(\Gamma), \quad (3.3)$$

with $f \in L^2(\Gamma)$, on V_h we test on all the ψ_k^e, ϕ_v and, because of $\mathcal{W}_v \cap \mathcal{W}_z = \emptyset$ if $v \neq z$, we have the following finite-dimensional (discrete) system of equations:

$$\left. \begin{aligned} \mathfrak{h}_h(w_h, \psi_k^e) &= \sum_{e \in \mathcal{E}} \sum_{j=1}^{n_e-1} \alpha_j^e \left\{ \int_e \frac{d\psi_j^e}{dx} \frac{d\psi_k^e}{dx} dx + \int_e \psi_j^e \psi_k^e v(x) dx \right\} \\ &\quad + \sum_{v \in \mathcal{V}} \beta_v \left\{ \int_{\mathcal{W}_v} \frac{d\phi_v}{dx} \frac{d\psi_k^e}{dx} dx + \int_{\mathcal{W}_v} \phi_v \psi_k^e v(x) dx \right\} \\ &= \int_{\text{supp}(\psi_k^e)} f \psi_k^e dx \quad \forall \psi_k^e; \\ \mathfrak{h}_h(w_h, \phi_v) &= \sum_{e \in \mathcal{E}} \sum_{j=1}^{n_e-1} \alpha_j^e \left\{ \int_{\mathcal{W}_v} \frac{d\psi_j^e}{dx} \frac{d\phi_v}{dx} dx + \int_{\mathcal{W}_v} \psi_j^e \phi_v v(x) dx \right\} \\ &\quad + \beta_v \left\{ \int_{\mathcal{W}_v} \frac{d\phi_v}{dx} \frac{d\phi_v}{dx} dx + \int_{\mathcal{W}_v} \phi_v \phi_v v(x) dx \right\} \\ &= \int_{\mathcal{W}_v} f \phi_v dx \quad \forall \phi_v. \end{aligned} \right\} \quad (3.4)$$

Under the hypotheses of the previous section (continuity and coercivity of the Hamiltonian), the system (3.4) has a unique solution in $V_h(\Gamma)$. Henceforth, we will denote with $u_h \in V_h(\Gamma)$ the solution of the discrete variational equation (3.4).

The following theorem gives the properties and the structure of the discretized Hamiltonian \mathbf{H} corresponding to the system (3.4).

THEOREM 3.1 Let us denote by

$$\mathbf{u}_{\mathcal{E}} = \begin{bmatrix} \mathbf{u}^1 \\ \vdots \\ \mathbf{u}^M \end{bmatrix}, \quad \text{where} \quad \mathbf{u}^e = \begin{bmatrix} u_1^e \\ \vdots \\ u_{n_e-1}^e \end{bmatrix} \quad (e = 1, \dots, M),$$

and by

$$\mathbf{u}_{\mathcal{V}} = \begin{bmatrix} u_1 \\ \vdots \\ u_N \end{bmatrix}$$

the values that give

$$u_h(x) = \sum_{e \in \mathcal{E}} \sum_{j=1}^{n_e-1} u_j^e \psi_j^e(x) + \sum_{v \in \mathcal{V}} u_v \phi_v(x) \quad x \in \Gamma.$$

Let $\mathbf{f}_{\mathcal{E}}$ and $\mathbf{f}_{\mathcal{V}}$ be the vectors

$$\mathbf{f}_{\mathcal{E}} = \begin{bmatrix} \mathbf{f}^1 \\ \vdots \\ \mathbf{f}^M \end{bmatrix}, \quad \text{where} \quad \mathbf{f}^e = \begin{bmatrix} f_1^e \\ \vdots \\ f_{n_e-1}^e \end{bmatrix},$$

$$\mathbf{f}_{\mathcal{V}} = \begin{bmatrix} f_1 \\ \vdots \\ f_N \end{bmatrix},$$

with $f_k^e = \int_{\text{supp}(\psi_k^e)} f \psi_k^e dx$ and $f_v = \int_{\mathcal{W}_v} f \phi_v dx$. Then, the linear system (3.4) can be written as

$$\mathbf{H} \begin{bmatrix} \mathbf{u}_{\mathcal{E}} \\ \mathbf{u}_{\mathcal{V}} \end{bmatrix} = \begin{bmatrix} \mathbf{f}_{\mathcal{E}} \\ \mathbf{f}_{\mathcal{V}} \end{bmatrix},$$

where \mathbf{H} has the following structure:

$$\mathbf{H} = \begin{bmatrix} \mathbf{H}_{11} & \mathbf{H}_{12} \\ \mathbf{H}_{12}^T & \mathbf{H}_{22} \end{bmatrix} = \begin{bmatrix} \mathbf{A} & \mathbf{B} \\ \mathbf{B}^T & \mathbf{G} \end{bmatrix} + \begin{bmatrix} \mathbf{V} & \mathbf{C} \\ \mathbf{C}^T & \mathbf{F} \end{bmatrix} = \mathbf{L} + \mathbf{M}_v. \quad (3.5)$$

Moreover,

- (i) \mathbf{H} is a symmetric positive definite matrix.
- (ii) Both \mathbf{A} and \mathbf{V} are block diagonal matrices and each diagonal block (of dimension $n_e - 1$) can be permuted into a symmetric tridiagonal nonsingular matrix corresponding to the edge e .
- (iii) The entries of \mathbf{A} are given by

$$\int_e \frac{d\psi_j^e}{dx} \frac{d\psi_k^e}{dx} dx = \begin{cases} 2/h_e & \text{if } j = k \\ -1/h_e & \text{if } |j - k| = 1 \\ 0 & \text{otherwise,} \end{cases} \quad (3.6)$$

whereas the entries of \mathbf{B} are given by

$$\int_{\mathcal{W}_v} \frac{d\phi_v}{dx} \frac{d\psi_k^e}{dx} dx = \begin{cases} -1/h_e & \text{if } e \in \mathcal{W}_v \cap \mathcal{E}_v \neq \emptyset \\ 0 & \text{otherwise.} \end{cases} \quad (3.7)$$

Moreover, \mathbf{G} is diagonal and the entries of \mathbf{G} are given by

$$\int_{\mathcal{W}_v} \frac{d\phi_v}{dx} \frac{d\phi_v}{dx} dx = \sum_{e \in \mathcal{E}_v} 1/h_e. \quad (3.8)$$

- (iv) The entries of \mathbf{V} are $\int_e \psi_j^e \psi_k^e v(x) dx$, the entries of \mathbf{C} are $\int_{\mathcal{W}_v} \phi_v \psi_k^e v(x) dx$, the entries of \mathbf{F} are $\int_{\mathcal{W}_v} \phi_v \phi_v v(x) dx$ and \mathbf{F} is diagonal. Moreover, we have

$$\left| \int_{\mathcal{W}_v \cap e} \phi_v \psi_k^e v(x) dx \right| \leq \chi h_e \quad \text{and} \quad \left| \int_{\mathcal{W}_v} \phi_v \phi_v v(x) dx \right| \leq \chi \sum_{e \in \mathcal{E}_v} h_e, \quad (3.9)$$

where $\chi = \max_{\Gamma} |v(x)|$.

Proof. (i) This follows from the coercivity and the continuity of the self-adjoint Hamiltonian.

(ii) The nodes on the edges will describe a chain path between two vertices, therefore the block corresponding to the edge can be permuted into a tridiagonal matrix ordering the internal nodes by increasing or decreasing distance from one of the vertices defining the edge. Because the supports of the hat functions corresponding to internal nodes of distinct edges have empty intersections, we have a block diagonal structure for both \mathbf{A} and \mathbf{V} .

(iii) The values of the entries in \mathbf{A} can be easily computed, observing that the derivative of each hat function ψ_j^e on the interval $[x_j^e, x_{j+1}^e]$ is a constant equal to $\pm 1/h_e$, depending on the orientation we choose on the edge. Taking into account the restrictions of the ϕ_v on the incident edges and the orientation chosen on the edge e , we have that on the interval $[x_v^e, x_1^e]$ (or $[x_{n_e-1}^e, x_v^e]$) the derivatives of ψ_1^e (respectively, $\psi_{n_e-1}^e$) and $\phi_v|_e$ are constants of opposite sign. Hence, we have that the integrals are equal to $-1/h_e$ for the nonzero entries in \mathbf{B} . The expression (3.8) for the diagonal entries of \mathbf{G} is straightforward.

(iv) The bounds (3.9) follow from $0 \leq \psi_j^e \leq 1$ and $0 \leq \phi_v \leq 1$.

This completes the proof. \square

Next, we have the following error estimate.

THEOREM 3.2 Let $f \in L^2(\Gamma)$ and $v \in L^\infty(\Gamma)$, with $v(x) \geq v_0 > 0$ on Γ . Then the approximate solution u_h of (3.3) satisfies

$$\|u - u_h\|_{H^1(\Gamma)} \leq \sqrt{\frac{\gamma_1}{\gamma_0}} \Theta C \hat{h} \sum_{e \in \mathcal{E}} \|u|_e\|_{H^2(e)}, \quad (3.10)$$

where $\hat{h} = \max_e h_e$, C is a constant independent of u and \hat{h} and $\Theta = \mathbf{vol}(\Gamma) = \sum_{e \in \mathcal{E}} \ell_e$, the volume of Γ .

Proof. As a consequence of the assumptions, the bilinear form

$$\mathfrak{h}(z, w) = \int_{\Gamma} \frac{dz}{dx} \frac{dw}{dx} dx + \int_{\Gamma} v(x) z w dx$$

on $H^1(\Gamma)$ is coercive and continuous, i.e.,

$$\begin{cases} \gamma_0 \|z\|_{H^1(\Gamma)}^2 \leq \mathfrak{h}(z, z) \\ |\mathfrak{h}(z, w)| \leq \gamma_1 \|z\|_{H^1(\Gamma)} \|w\|_{H^1(\Gamma)}. \end{cases}$$

From these inequalities, we have that the energy norm $\|z\|^2 = \mathfrak{h}(z, z)$ and the $H^1(\Gamma)$ norm are equivalent. From the fact that $V_h(\Gamma) \subset H^1(\Gamma)$, it follows that $z = u - u_h$ minimizes the energy norm of z in $V_h(\Gamma)$ and, thus,

$$\gamma_0 \|u - u_h\|_{H^1(\Gamma)}^2 \leq \mathfrak{h}(z, z) \leq \mathfrak{h}(u - u_h^I, u - u_h^I) \leq \gamma_1 \|u - u_h^I\|_{H^1(\Gamma)}^2,$$

where u_h^I is the interpolant of u in the nodes and the vertices. We observe that under our assumptions, the solution to (3.3) satisfies $u|_e \in H^2(e)$ on each edge $e \in \mathcal{E}$ (see Brezis, 2010, Chapter 8). On each edge

$e = (\mathbf{v}_a, \mathbf{v}_b)$, the ψ_j^e , ($j = 1, \dots, n_e - 1$) and the restrictions to e of $\phi_{\mathbf{v}_a}$ and $\phi_{\mathbf{v}_b}$ form a basis for the finite element approximation of $H^1(e)$. Thus, we have the standard error estimate

$$\|u|_e - u_h|_e\|_{H^1(e)} \leq \sqrt{\frac{\gamma_1}{\gamma_0}} \ell_e C h_e \|u|_e\|_{H^2(e)} \quad \forall e \in \mathcal{E}, \quad (3.11)$$

where C is independent of u and h_e (see [Raviart et al., 1983](#), Chapter 3.2; [Brenner & Scott, 2002](#), Section 4.4). Letting now $\hat{h} = \max_e h_e$, the sum of all the local errors and of their upper bounds gives the global upper bound (3.10) for $\|u - u_h\|_{H^1(\Gamma)}$ (see [Dupont & Scott, 1980](#), Theorem 7.1). \square

3.1 Error estimates for the Neumann–Kirchhoff conditions

In general, unfortunately, the Neumann–Kirchhoff conditions cannot be exactly satisfied for any value of $h > 0$; for the special case of a single edge (i.e., an interval), see the discussion in [Raviart et al. \(1983, p. 71\)](#). However, they are asymptotically satisfied as the discretization is refined.

The following result provides an upper bound on how much the discrete solution u_h can deviate from satisfying the Neumann–Kirchhoff conditions at a given vertex \mathbf{v} of Γ .

THEOREM 3.3 If $f \in L^2(\Gamma)$, then for any vertex \mathbf{v} of Γ with neighboring set $\mathcal{W}_{\mathbf{v}}$ and degree $d_{\mathbf{v}}$, the finite element solution u_h of (3.4) satisfies

$$\left| \sum_{e \in \mathcal{V}_{\mathbf{v}}} \frac{du_h}{dx^e} \right| \leq 2 \mathbf{vol}(V_{\mathbf{v}}) \left(\|f\|_{L^2(\Gamma)} + \chi \|u\|_{L^2(\Gamma)} \right) \leq 2d_{\mathbf{v}} \hat{h} \left(\|f\|_{L^2(\Gamma)} + \chi \|u\|_{L^2(\Gamma)} \right), \quad (3.12)$$

where $\hat{h} = \max_e h_e$ and $\chi = \max_{\Gamma} |v(x)|$.

Proof. The entry in the right-hand side corresponding to a $\mathbf{v} \in \mathcal{V}$ is given by

$$f_{\mathbf{v}} = \int_{\mathcal{W}_{\mathbf{v}}} f \phi_{\mathbf{v}} dx.$$

Taking into account that the supports of each $\phi_{\mathbf{v}}|_e$ have length h_e and, using the triangle and Cauchy–Schwarz inequalities, we get

$$|f_{\mathbf{v}}| \leq \mathbf{vol}(V_{\mathbf{v}}) \|f\|_{L^2(\Gamma)} \leq \hat{h} d_{\mathbf{v}} \|f\|_{L^2(\Gamma)}. \quad (3.13)$$

Taking into account that on $\mathcal{W}_{\mathbf{v}} \cap e$ the function $u_h(x)$ is linear, and thus differentiable, we obtain from (3.7) and (3.8) that

$$[\mathbf{B}^T \mathbf{u}_{in} + \mathbf{G} \mathbf{u}_{\mathcal{V}}]_{\mathbf{v}} = \sum_{e \in \mathcal{E}_{\mathbf{v}}} \frac{du_h}{dx^e}(\mathbf{v}).$$

Similarly, from (3.1), we have that the entry corresponding to \mathbf{v} in $[\mathbf{C}^T \mathbf{u}_{in} + \mathbf{F} \mathbf{u}_{\mathcal{V}}]$ satisfies

$$[\mathbf{C}^T \mathbf{u}_{in} + \mathbf{F} \mathbf{u}_{\mathcal{V}}]_{\mathbf{v}} = \int_{\mathcal{W}_{\mathbf{v}}} u_h \phi_{\mathbf{v}} v(x) dx.$$

Therefore, from (3.13) we have that at each vertex v and for ϕ_v ,

$$\sum_{e \in \mathcal{W}_v} \frac{du_h}{dx^e}(v) = \int_{\mathcal{W}_v} (f - u_h v(x)) \phi_v \, dx,$$

and thus from (3.9), we have

$$\left| \sum_{e \in V_v} \frac{du_h}{dx^e}(v) \right| \leq 2 \mathbf{vol}(V_v) \left(\|f\|_{L^2(\Gamma)} + \chi \|u\|_{L^2(\Gamma)} \right) \leq 2d_v \hat{h} \left(\|f\|_{L^2(\Gamma)} + \chi \|u\|_{L^2(\Gamma)} \right).$$

This completes the proof. \square

Inequality (3.12) shows that, in general, the Neumann–Kirchhoff interface conditions are satisfied only in the limit for $\hat{h} \searrow 0$.

REMARK 3.4 Taking into account (3.13) and (3.12), for vertices v with a large degree $d_v \gg 1$ (so-called *hubs*), the Neumann–Kirchhoff conditions can be poorly satisfied even for a reasonably small value of h . This suggests that in these cases, an adaptive choice of the mesh points should be made. In particular, the analysis suggests that on the edges having an end point corresponding to a hub, the mesh node x_1^e (or $x_{n_e-1}^e$) in \mathcal{W}_v should be chosen much closer to x_v^e than the distance between two consecutive nodes on e , i.e., $|x_1^e - x_v^e| = \tilde{h}_e$ (or $|x_{n_e-1}^e - x_v^e| = \tilde{h}_e$), where \tilde{h} is such that $\tilde{h}_e d_v \leq \hat{h}$.

4. Extended graphs

The discretization of a metric graph Γ leads naturally to a new graph \mathcal{G} , which we refer to as the *extended graph*, having as vertices the vertices of Γ plus all the discretization nodes x_j^e and as edges all the intervals cut out by the discretization on each edge.

Hereafter, given a set of matrices $\{\mathbf{Y}_k\}_{k=1}^K$, we will denote by $\mathbf{blkdiag}(\{\mathbf{Y}_k\}_{k=1}^K)$ the block diagonal matrix obtained using the \mathbf{Y}_k as diagonal blocks. We remark that we do not require that the blocks be square.

4.1 Construction of the coefficient matrices

It is possible to build the matrix \mathbf{H} using an extended incidence matrix $\tilde{\mathbf{E}}$ obtained from the incidence matrix \mathbf{E} used to describe the original graph Γ . In particular, as $\mathbf{H} = \mathbf{L} + \mathbf{M}$, where \mathbf{L} is the stiffness matrix and \mathbf{M} is the mass matrix (or more generally the matrix \mathbf{M}_v describing the discretized potential $v(x)$, see (3.5)) on the extended graph \mathcal{G} , we will focus on the construction of \mathbf{L} and \mathbf{M} using \mathbf{E} . Let us define the matrices

$$\mathbf{E}^+ = \frac{1}{2}(\mathbf{E} + |\mathbf{E}|) \quad \text{and} \quad \mathbf{E}^- = \frac{1}{2}(\mathbf{E} - |\mathbf{E}|),$$

where $|\mathbf{E}|$ denotes the entry-wise absolute value of \mathbf{E} . Note that $\mathbf{E} = \mathbf{E}^+ + \mathbf{E}^-$.

LEMMA 4.1 Let $\mathbf{E} \in \mathbb{R}^{N \times M}$ be the incidence matrix describing a graph with N vertices and M edges without loops. We have

$$\begin{aligned} \mathbf{E}^+ (\mathbf{E}^+)^T + \mathbf{E}^- (\mathbf{E}^-)^T &= \mathbf{D}, \\ \mathbf{E}^+ (\mathbf{E}^-)^T + \mathbf{E}^- (\mathbf{E}^+)^T &= -\mathbf{Ad}, \end{aligned} \quad (4.1)$$

where \mathbf{D} and \mathbf{Ad} are the degree and adjacency matrix of the graph, respectively.

Proof. First note that the rows of both \mathbf{E}^+ and \mathbf{E}^- are mutually orthogonal, since we separate the incoming edges from the out-going edges in \mathbf{E} . Thus, both $\mathbf{E}^+ (\mathbf{E}^+)^T$ and $\mathbf{E}^- (\mathbf{E}^-)^T$ are diagonal matrices. Every row in each matrix corresponds to a vertex in the graph and will have a number of nonzeros equal to the number of edges that, respectively, have that vertex as the origin (\mathbf{E}^+) or ending (\mathbf{E}^-). Therefore, the sum of the two diagonal matrices will result in the total number of edges incident to each vertex. The second relation is a straightforward consequence of $\mathbf{E}\mathbf{E}^T = \mathbf{L}_r = \mathbf{D} - \mathbf{Ad}$ and $\mathbf{E} = \mathbf{E}^+ + \mathbf{E}^-$. \square

Consider now the $n_e - 1$ nodes interior to the generic edge e . Because of the chain structure of them, to each edge we associate the $(n_e - 1) \times n_e$ matrix

$$\mathbf{E}_e = \begin{bmatrix} -1 & 1 & & & & \\ & -1 & 1 & & & \\ & & & \ddots & \ddots & \\ & & & & -1 & 1 \end{bmatrix}.$$

Note that, strictly speaking, this is not an incidence matrix, because of the first and last column. Next, we define the block diagonal matrix $\tilde{\mathbf{E}} = \mathbf{blkdiag}(\{\mathbf{E}_e\}_{e \in \mathcal{E}})$ of size $\tilde{n} \times (\tilde{n} + M)$, where $\tilde{n} = \sum_{e \in \mathcal{E}} (n_e - 1)$. If we denote by $\mathbf{e}_1^{n_e}$ and by $\mathbf{e}_{n_e}^{n_e}$, respectively, the first and the last column of the identity matrix \mathbf{I}_{n_e} , $e \in \mathcal{E}$, then we can build the $N \times (\tilde{n} + M)$ matrix $\hat{\mathbf{E}}$ by expanding each column \mathbf{E}_j^+ of \mathbf{E}^+ in an $N \times n_{e_j}$ block equal to

$$\hat{\mathbf{E}}_j^+ = \mathbf{E}_j^+ \otimes (\mathbf{e}_1^{n_{e_j}})^T, \quad (4.2)$$

and by expanding each column \mathbf{E}_j^- of \mathbf{E}^- in an $N \times n_{e_j}$ block equal to

$$\hat{\mathbf{E}}_j^- = \mathbf{E}_j^- \otimes (\mathbf{e}_{n_{e_j}}^{n_{e_j}})^T. \quad (4.3)$$

Then, by adding these matrices, we can construct the new matrix

$$\hat{\mathbf{E}} = [\hat{\mathbf{E}}_1^+ + \hat{\mathbf{E}}_1^-, \dots, \hat{\mathbf{E}}_M^+ + \hat{\mathbf{E}}_M^-]. \quad (4.4)$$

Finally, the incidence matrix of the extended graph \mathcal{G} is given by

$$\tilde{\mathbf{E}} = \begin{bmatrix} \tilde{\mathbf{E}} \\ \hat{\mathbf{E}} \end{bmatrix} \in \mathbb{R}^{(\tilde{n}+N) \times (\tilde{n}+M)}.$$

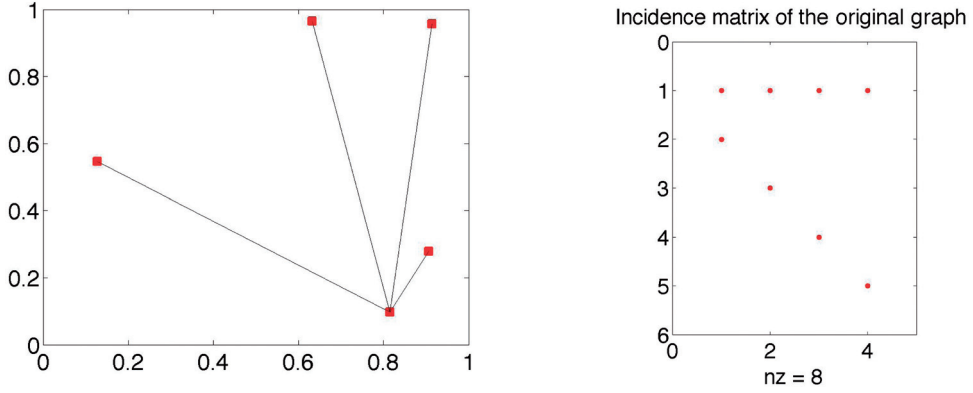


FIG. 2. Example of a simple planar metric graph and its incidence matrix.

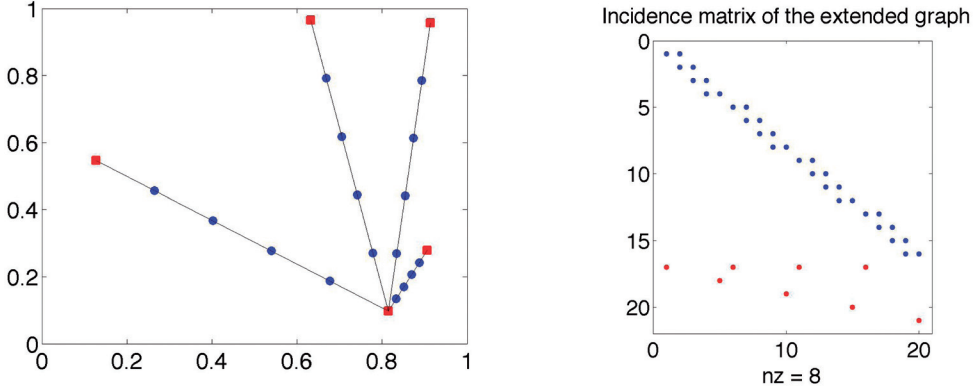


FIG. 3. Example of the extension of the graph given in Fig. 2 when four-node chains are added internally to each edge and its incidence matrix.

In Fig. 2, we give an example of a simple planar graph and its incidence matrix. In Fig. 3, we show the extended graph and its corresponding incidence matrix when four extra nodes are added internally to each edge.

Next, we introduce the following block diagonal matrix of weights:

$$\mathbf{W} = \text{blkdiag} \left(\left\{ \frac{1}{h_e} \mathbf{I}_{n_e} \right\}_{e \in \mathcal{E}} \right).$$

The matrix

$$\mathbf{L} = \tilde{\mathbf{E}} \mathbf{W} \tilde{\mathbf{E}}^T = \begin{bmatrix} \mathbf{A} & \mathbf{B} \\ \mathbf{B}^T & \mathbf{G} \end{bmatrix}, \quad (4.5)$$

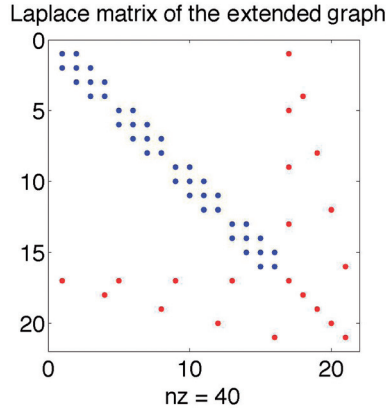


FIG. 4. Pattern of the matrix \mathbf{L} , where the red bullets correspond to the original vertices and the blue ones to the internal nodes on each edge. See online version for colour.

where the leading principal block has order \tilde{n} and is block diagonal, is the stiffness matrix corresponding to the Laplace operator $-\frac{d^2}{dx^2}$ on Γ . In Fig. 4 we show the nonzero pattern of \mathbf{L} , where the blue bullets correspond to the internal nodes on each edge and the red ones are the original vertices.

The matrix \mathbf{M}_v , describing the discretized potential (which reduces to the mass matrix \mathbf{M} for $v(x) \equiv 1$) requires the values of the integrals

$$\int_e \psi_j^e \psi_k^e v(x) dx, \quad \int_{\mathcal{W}_v} \phi_v \psi_k^e v(x) dx \quad \text{and} \quad \int_{\mathcal{W}_v} \phi_v \phi_v v(x) dx,$$

which can be computed numerically using either the trapezoidal formula or Simpson's rule (Raviart *et al.*, 1983). In the first case, we need to have the values of the function $v(x)$ on all the vertices of the extended graph, i.e., in each vertex in the set

$$\widehat{\mathcal{V}} = \mathcal{V} \cup \left(\bigcup_e \mathcal{V}_e \right).$$

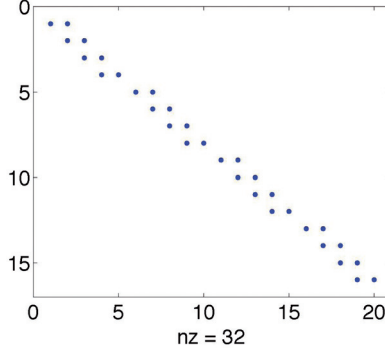
In the second case, we also need the values of $v(x)$ at each vertex in the set

$$\mathcal{V}_{\text{mid}} = \bigcup_e \widetilde{\mathcal{V}}_e,$$

where

$$\widetilde{\mathcal{V}}_e = \left\{ \frac{x_{j+1}^e + x_j^e}{2} \right\}_{j=1}^{n_e-2} \cup \left\{ \frac{x_a^e + x_1^e}{2} \right\} \cup \left\{ \frac{x_b^e + x_{n_e-1}^e}{2} \right\}.$$

Let $\mathbf{K}_e \in \mathbb{R}^{(n_e-1) \times (n_e-1)}$ be the diagonal matrices $\mathbf{K}_e = \mathbf{diag}(v(\mathcal{V}_e))$ for all $e \in \mathcal{E}$, i.e., the matrices having the values of $v(x)$ on the set \mathcal{V}_e on the main diagonal; let $\mathbf{K}_{\mathcal{V}} \in \mathbb{R}^{\tilde{n} \times \tilde{n}}$ be the diagonal matrix $\mathbf{K}_{\mathcal{G}} = \mathbf{blkdiag}(\{\mathbf{K}_e\}_{e \in \mathcal{E}})$ and $\mathbf{K}_{\mathcal{V}} \in \mathbb{R}^{N \times N}$ be the diagonal matrix $\mathbf{K}_{\mathcal{V}} = \mathbf{diag}(v(\mathcal{V}))$ with the values of

FIG. 5. Example of a simple planar metric graph: $\tilde{\mathbf{E}}$.

$\nu(x)$ on the vertices of the original graph Γ on the main diagonal. Finally, we can assemble the diagonal matrix $\mathbf{K}_1 \in \mathbb{R}^{(\tilde{n}+N) \times (\tilde{n}+N)}$:

$$\mathbf{K}_1 = \mathbf{blkdiag}(\mathbf{K}_{\mathcal{E}}, \mathbf{K}_{\mathcal{V}}).$$

Next, we form the diagonal matrix $\mathbf{K}_2 \in \mathbb{R}^{(\tilde{n}+M) \times (\tilde{n}+M)}$ given by $\mathbf{K}_2 = \mathbf{blkdiag}(\nu(\tilde{\mathcal{Y}}_2))$ and the matrix

$$\widehat{\mathbf{W}} = \mathbf{blkdiag}(\{h_e \mathbf{I}_{n_e}\}_{e \in \mathcal{E}}).$$

Therefore, the mass matrix \mathbf{M} and the potential part \mathbf{M}_ν of the Hamiltonian computed by the Simpson quadrature rule are, respectively,

$$\mathbf{M} = \frac{1}{6} \left(|\tilde{\mathbf{E}}| \widehat{\mathbf{W}} |\tilde{\mathbf{E}}|^T + \mathbf{diag} \left(\left\{ (|\tilde{\mathbf{E}}| \widehat{\mathbf{W}} |\tilde{\mathbf{E}}|^T)_{i,i} \right\}_{i=1}^{\tilde{n}+M} \right) \right) \text{ and}$$

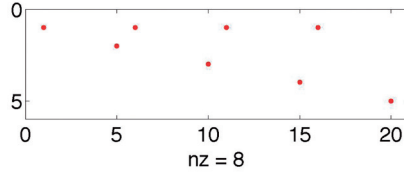
$$\mathbf{M}_\nu = \frac{1}{6} \left(|\tilde{\mathbf{E}}| \mathbf{K}_2 \widehat{\mathbf{W}} |\tilde{\mathbf{E}}|^T + \mathbf{K}_1 \mathbf{diag} \left(\left\{ (|\tilde{\mathbf{E}}| \widehat{\mathbf{W}} |\tilde{\mathbf{E}}|^T)_{i,i} \right\}_{i=1}^{\tilde{n}+M} \right) \right).$$

If instead the simple trapezoidal rule is used for the computation of the mass matrix and the potential then we will take the diagonal part of the two previous matrices. Although the Simpson rule, being exact for quadratic polynomials, gives ‘exact’ values for the entries of \mathbf{M} (and also of \mathbf{M}_ν when $\nu(x) = \nu$ with $\nu > 0$ a constant), the trapezoidal rule will only give approximations that are accurate up to an error of $\mathcal{O}(h^2)$. Clearly, in the limit $h \rightarrow 0$, the bound (3.10) remains valid.

We observe that the mass matrix \mathbf{M} has a block structure that matches the block structure of \mathbf{L} , Theorem 3.1:

$$\mathbf{M} = \begin{bmatrix} \mathbf{V} & \mathbf{C} \\ \mathbf{C}^T & \mathbf{F} \end{bmatrix}. \quad (4.6)$$

In both cases, the block structure pattern is a straightforward consequence of the partitioning of the rows in $\tilde{\mathbf{E}}$. In Figs 5 and 6, we give the structures of $\tilde{\mathbf{E}}$ and $\widehat{\mathbf{E}}$ for the simple graph in Fig. 2.


 FIG. 6. Example of a simple planar metric graph: $\widehat{\mathbf{E}}$.

In matrices \mathbf{L} and \mathbf{M} the $N \times \tilde{n}$ blocks \mathbf{B}^T and \mathbf{C}^T are given by

$$\mathbf{B}^T = \widehat{\mathbf{E}}\widehat{\mathbf{E}}^T\mathbf{W}_1 \quad \text{and} \quad \mathbf{C}^T = |\widehat{\mathbf{E}}| |\widehat{\mathbf{E}}^T|\widehat{\mathbf{W}}_1,$$

where \mathbf{W}_1 and $\widehat{\mathbf{W}}_1$ are the first blocks corresponding to $\widehat{\mathbf{E}}^T$ (and its absolute value) of \mathbf{W} and $\widehat{\mathbf{W}}$. Because of the lower triangular structure of $\widehat{\mathbf{E}}^T$ and $|\widehat{\mathbf{E}}|^T$, the pattern of \mathbf{B}^T and \mathbf{C}^T will be the same of $\widehat{\mathbf{E}}$ (see Figs 5 and 6). Therefore, we have that both $\mathbf{B}^T\mathbf{B}$ and $\mathbf{C}^T\mathbf{C}$ are diagonal.

Finally, we observe that $\mathbf{G} = \widehat{\mathbf{E}}\mathbf{W}_2\widehat{\mathbf{E}}^T$, where \mathbf{W}_2 is the second diagonal block of \mathbf{W} corresponding to $\widehat{\mathbf{E}}$, is also a diagonal matrix. Similar considerations show that \mathbf{F} is also a diagonal matrix. In particular, for the Laplace case with $n_e = n$ and $\ell_e = \ell$ for all $e \in \mathcal{E}$, we have that

$$\mathbf{B}^T\mathbf{B} = \frac{1}{h}\mathbf{D} = \mathbf{G} \quad \text{and} \quad \mathbf{C}^T\mathbf{C} = \frac{h}{3}\mathbf{D} = \mathbf{F}. \quad (4.7)$$

Furthermore, we have $h_e = h$ and

$$\mathbf{B}^T = \frac{1}{h} \left(\mathbf{E}^+ \otimes (\mathbf{e}_1^n)^T + \mathbf{E}^- \otimes (\mathbf{e}_n^n)^T \right) (\mathbf{I}_M \otimes \mathbf{E}_e^T), \quad (4.8)$$

$$\mathbf{C}^T = \frac{h}{6} \left(|\mathbf{E}^+| \otimes (\mathbf{e}_1^n)^T + |\mathbf{E}^-| \otimes (\mathbf{e}_n^n)^T \right) (\mathbf{I}_M \otimes |\mathbf{E}_e|^T) = \frac{-h^2}{6}\mathbf{B}^T. \quad (4.9)$$

In Section 4.2 we will describe the relation between the underlying domain decomposition techniques and the original graph Laplacian.

REMARK 4.2 If in the original graph all the ℓ_e are equal, all the n_e are equal to n and $v(x) = v$ with v a constant in the Hamiltonian (2.2), then the extended graph will produce a matrix \mathbf{H}

$$\mathbf{H} = \mathbf{L} + v\mathbf{M}.$$

Moreover, if the Hamiltonian (2.2) is given by

$$-\frac{d}{dx} \left(\alpha(x) \frac{du}{dx} \right) \quad \alpha(x) \geq \alpha_0 > 0, \quad \alpha(x) \in L^\infty(\Gamma)$$

and $v(x)|_e = v_e$ with v_e constant on each edge e (the Hamiltonian is still continuous and coercive), it suffices to modify the weight matrix \mathbf{W} as

$$\mathbf{W} = \text{blkdiag} \left(\left\{ \frac{1}{h_e} \text{diag}(\{w_j\}_{j=1}^{n_e}) \right\}_{e \in \mathcal{E}} \right),$$

where

$$w_j = \frac{1}{h_e} \int_{(j-1)h_e}^{jh_e} \alpha(x) \, dx,$$

the matrix $\widehat{\mathbf{W}}$ as

$$\widehat{\mathbf{W}} = \text{blkdiag}(\{h_e \nu_e \mathbf{I}_{n_e}\}_{e \in \mathcal{E}}),$$

and to proceed as described above to obtain the matrix \mathbf{H} approximating the infinite-dimensional Hamiltonian \mathcal{H} . We finally remark that the block structure of \mathbf{H} (Theorem 3.1) will be determined by the corresponding block structure of \mathbf{L} and \mathbf{M} .

4.2 Extended graph and domain decomposition

In this section, we investigate the structure of the Schur complement system arising from the elimination of the nodes internal to each edge. We begin by observing that, because of the node ordering scheme used, the matrix \mathbf{H} is partitioned as in a (nonoverlapping) domain decomposition approach, where individual edges $e \in \mathcal{E}$ of the original metric graph are the subdomains and the vertices $v \in \mathcal{V}$ are the interfaces between them.

The following theorem shows that, under certain conditions, the Schur complement of the discretization of the second derivative operator coincides with the graph Laplacian of the underlying combinatorial graph, Γ .

THEOREM 4.3 Assume that $n_e = n$ for all $e \in \mathcal{E}$ and that all edges in the metric graph Γ have the same length $\ell_e = \ell = 1$. Assume further that $v(x) \equiv 0$. Then,

$$\mathbf{S} := \mathbf{G} - \mathbf{B}^T \mathbf{A}^{-1} \mathbf{B} = \mathbf{L}_\Gamma.$$

Proof. Under our hypotheses, we have

$$\mathbf{L} = n \widetilde{\mathbf{E}} \widetilde{\mathbf{E}}^T,$$

see (4.5). Because of the partition of the rows of $\widetilde{\mathbf{E}}$ into nodes internal to the edges (corresponding to rows of $\widetilde{\mathbf{E}}$) and the vertices of the original graph (corresponding to rows of $\widehat{\mathbf{E}}$), we have

$$\begin{aligned} \mathbf{A} &= n \widetilde{\mathbf{E}}^T, \\ \mathbf{B} &= n \widetilde{\mathbf{E}} \widehat{\mathbf{E}}^T, \\ \mathbf{G} &= n \widehat{\mathbf{E}} \widehat{\mathbf{E}}^T = n \mathbf{D}. \end{aligned} \tag{4.10}$$

The nonzero pattern of the Schur complement

$$\mathbf{S} = \mathbf{G} - \mathbf{B}^T \mathbf{A}^{-1} \mathbf{B}$$

coincides with that of

$$\mathbf{B}^T \mathbf{A}^{-1} \mathbf{B} = n \widehat{\mathbf{E}} \left(\widetilde{\mathbf{E}}^T (\widetilde{\mathbf{E}} \widetilde{\mathbf{E}}^T)^{-1} \widetilde{\mathbf{E}} \right) \widehat{\mathbf{E}}^T.$$

We point out that the matrix $n \bar{\mathbf{E}}^T (\bar{\mathbf{E}} \bar{\mathbf{E}}^T)^{-1} \bar{\mathbf{E}}$ is block diagonal and that each of the M blocks can be easily computed for our simple Hamiltonian:

$$n \bar{\mathbf{E}}^T (\bar{\mathbf{E}} \bar{\mathbf{E}}^T)^{-1} \bar{\mathbf{E}} = \mathbf{I}_M \otimes \mathbf{T}, \quad (4.11)$$

with

$$\mathbf{T} = n \mathbf{I}_n - \mathbb{1}_n \mathbb{1}_n^T,$$

where $\mathbb{1}_j$ is the vector of all ones of dimension j . Moreover, from (4.4), the matrix $\widehat{\mathbf{E}}$ is a stretching of the incidence matrix of the original graph. Under our current hypotheses, it is given by

$$\widehat{\mathbf{E}} = \mathbf{E}^+ \otimes (\mathbf{e}_1^n)^T + \mathbf{E}^- \otimes (\mathbf{e}_n^n)^T, \quad (4.12)$$

where $\mathbf{e}_1^n \in \mathbb{R}^n$ and $\mathbf{e}_n^n \in \mathbb{R}^n$ are, respectively, the first and the last column of the identity matrix. The Schur complement is given by

$$\begin{aligned} \mathbf{S} &= \mathbf{G} - \left(\mathbf{E}^+ \otimes (\mathbf{e}_1^n)^T + \mathbf{E}^- \otimes (\mathbf{e}_n^n)^T \right) (\mathbf{I}_M \otimes \mathbf{T}) \left(\mathbf{E}^+ \otimes (\mathbf{e}_1^n)^T + \mathbf{E}^- \otimes (\mathbf{e}_n^n)^T \right)^T \\ &= \mathbf{G} - \left(\mathbf{E}^+ \otimes (\mathbf{e}_1^n)^T + \mathbf{E}^- \otimes (\mathbf{e}_n^n)^T \right) \left((\mathbf{E}^+)^T \otimes \mathbf{T} \mathbf{e}_1^n + (\mathbf{E}^-)^T \otimes \mathbf{T} \mathbf{e}_n^n \right) \\ &= \mathbf{G} - \left(\mathbf{E}^+ (\mathbf{E}^+)^T \otimes (\mathbf{e}_1^n)^T \mathbf{T} \mathbf{e}_1^n + \mathbf{E}^- (\mathbf{E}^-)^T \otimes (\mathbf{e}_n^n)^T \mathbf{T} \mathbf{e}_n^n \right. \\ &\quad \left. + \mathbf{E}^+ (\mathbf{E}^-)^T \otimes (\mathbf{e}_1^n)^T \mathbf{T} \mathbf{e}_n^n + \mathbf{E}^- (\mathbf{E}^+)^T \otimes (\mathbf{e}_n^n)^T \mathbf{T} \mathbf{e}_1^n \right). \end{aligned}$$

Moreover, from the following relations:

$$(\mathbf{e}_1^n)^T \mathbf{T} \mathbf{e}_1^n = (\mathbf{e}_n^n)^T \mathbf{T} \mathbf{e}_n^n = n - 1 \quad \text{and} \quad (\mathbf{e}_n^n)^T \mathbf{T} \mathbf{e}_1^n = (\mathbf{e}_1^n)^T \mathbf{T} \mathbf{e}_n^n = -1,$$

and from Lemma 4.1, we have, taking into account that $\mathbf{G} = n \widehat{\mathbf{E}} \widehat{\mathbf{E}}^T$, that

$$\begin{aligned} \mathbf{S} &= \mathbf{G} - (n - 1) \left(\mathbf{E}^+ (\mathbf{E}^+)^T + \mathbf{E}^- (\mathbf{E}^-)^T \right) + \left(\mathbf{E}^- (\mathbf{E}^+)^T + \mathbf{E}^+ (\mathbf{E}^-)^T \right) \\ &= n \mathbf{D} - (n - 1) \mathbf{D} - \mathbf{A} \mathbf{d} = \mathbf{E} \mathbf{E}^T = \mathbf{L}_\Gamma. \end{aligned}$$

This completes the proof. \square

In Fig. 7, we display the sparsity pattern of the product $\mathbf{A}^{-1} \mathbf{B}$ for the simple example considered in Figs 2 and 3.

Hence, we have shown that upon elimination of the internal nodes on the edges, the resulting Schur complement reduces to the combinatorial Laplacian associated with the original graph. In particular, \mathbf{S} is sparse. This fact is important enough to warrant further comment. Indeed, since the inverse of an irreducible tridiagonal matrix is full (Meurant, 1992), *a priori* one could have expected the Schur complement $\mathbf{S} = \mathbf{G} - \mathbf{B}^T \mathbf{A}^{-1} \mathbf{B}$ to incur significant fill-in, similar to what happens when solving discretized elliptic PDEs in two-dimensional and three-dimensional domains. The fact that \mathbf{S} is sparse has important implications when solving the discretized equations. This will be discussed further in Section 7. We observe that the Schur complement remains sparse even if we drop the assumption that all edges on the metric graph Γ have equal length, and that the same number of interior nodes is used in discretizing each edge. In this case, however, it is no longer true that \mathbf{S} coincides with the graph Laplacian \mathbf{L}_Γ .

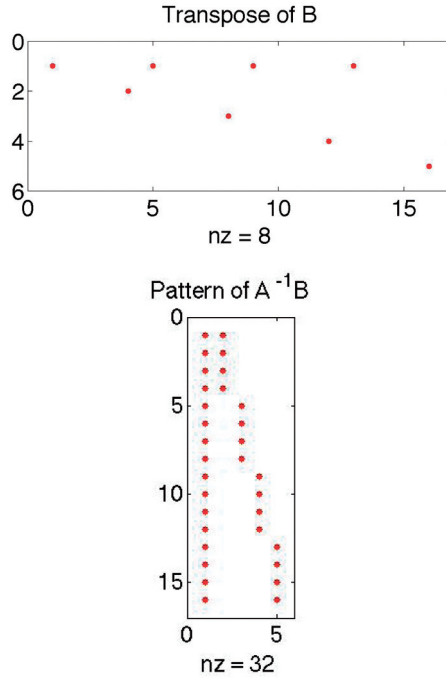


FIG. 7. Example of a simple planar metric graph: \mathbf{B}^T (top) and $\mathbf{A}^{-1}\mathbf{B}$ (bottom).

REMARK 4.4 In Theorem 4.3, we have assumed that $v(x) \equiv 0$. In the general case, where we have a more complex Hamiltonian, the Schur complement will contain additional information; indeed, the values in a vertex v now take into account the contributions of the solutions on the edges incident to it. However, the above argument shows that the nonzero pattern of the resulting \mathbf{S} will still coincide with the nonzero pattern of \mathbf{L}_Γ . Additional details are given in the appendix.

5. Conditioning of the mass matrix

In this section, we examine the conditioning of the mass matrix \mathbf{M} in terms of h and of the maximum degree of Γ . It is well known (Wathen, 1987) that for a number of different types of finite elements in one, two and three dimensions, the condition number of the mass matrix is independent of h , and that a simple diagonal scaling of the mass matrix yields a well-conditioned matrix. This fact will be useful when we discuss the solution of parabolic problems. In our situation, the condition number of \mathbf{M} is also independent of h . However, in the case of complex graphs, the spectrum can be quite spread out. With the scaling adopted in this article, and assuming for simplicity a constant h throughout the graph, the eigenvalues range between $\mathcal{O}(h)$ and $\mathcal{O}(d_{\max}h)$, where d_{\max} is the maximum degree of a vertex in Γ . More precisely, we have the following result.

THEOREM 5.1 Let Γ be a graph having N vertices and M edges and with at least one node v_i of degree $d_{v_i} > 6$. Let \mathbf{M} be the mass matrix relative to a piecewise linear, continuous finite-element approximation

of the $L^2(\Gamma)$ norm by using n nodes on each edge of Γ , and $h = \frac{1}{n+1}$. Then the spectrum $\mathbf{sp}(\mathbf{M})$ of \mathbf{M} satisfies

$$\mathbf{sp}(\mathbf{M}) \subset \frac{h}{6} [\mathcal{O}(1), \mathcal{O}(d_{\max})].$$

Proof. The proof is a simple application of Gershgorin's Theorem. Let $\tilde{\mathbf{M}} = 6(n+1)\mathbf{M}$, then we have that the diagonal entries of $\tilde{\mathbf{M}}_{ii}$ for $i = 1, \dots, nM$ are equal to 4 and the corresponding off-diagonal sum of the entries is 2 (note that the entries of \mathbf{M} are non-negative). The last N diagonal entries of $\tilde{\mathbf{M}}_{ii}$, $i = nM+1, \dots, nM+N$ are equal to $2d_{v_i}$, and the sum of the corresponding off-diagonal entries is equal to d_{v_i} . Therefore, under the hypothesis $d_{v_i} > 6$, the spectrum of $\tilde{\mathbf{M}}$ lays in the union of two disjoint sets and at least one eigenvalue lays in the union of the circles corresponding to the highest degree. \square

REMARK 5.2 If the highest degree of a vertex in Γ is much larger than the others, then one eigenvalue of $\tilde{\mathbf{M}}$ is of the same order of it. Numerical experiments on graphs with a power law degree distribution indicate that the condition number of the mass matrix $\tilde{\mathbf{M}}$ grows like $\mathcal{O}(d_{\max})$ as N , and therefore d_{\max} increases. We note that this phenomenon has no analog in the 'standard' finite element theory because of the requirement that the triangles or the tetrahedra must satisfy the minimum angle condition (Raviart *et al.*, 1983, p. 109), which imposes a bound on the maximum number of edges incident to a vertex, independent of N .

6. A generalized eigenvalue problem

In this section, we analyse the generalized eigenvalue problem

$$\mathbf{H}\mathbf{w} = \lambda\mathbf{M}\mathbf{w}. \tag{6.1}$$

The finite element discretization of the eigenvalue problem $\mathcal{H}u = \lambda u$ leads to algebraic eigenvalue problems of the form (6.1). In particular, for $v(x) = 0$ the eigensolutions (λ, \mathbf{w}) of (6.1) are approximations of the eigenvalues and eigenfunctions of the simple Hamiltonian $\mathcal{H} = -\frac{d^2}{dx^2}$ (we assume Neumann–Kirchhoff conditions throughout). This operator has discrete spectrum and a complete orthonormal basis of eigenfunctions (as a special case of Berkolaiko & Kuchment, 2013, Theorem 3.1.1). Note that $\mathbf{H} = \mathbf{L}$ is positive semidefinite, and the kernel of \mathbf{H} is spanned by the vector of all ones. Our motivation for the study of (6.1) is that it plays a central role in the solution of diffusion problems on Γ . The spectrum of a quantum graph is also of fundamental interest in scattering theory, photonics, quantum chaos, and so forth (Berkolaiko & Kuchment, 2013). For any sufficiently small, fixed h , only the leftmost part of the spectrum of \mathcal{H} , which is unbounded above, is well approximated by the corresponding eigenvalues of (6.1). However, in many problems, these are the only eigenvalues of interest.

In the study of diffusion processes on combinatorial graphs Γ , the behavior of the solution is essentially determined by the spectrum of the graph Laplacian \mathbf{L}_Γ , which, to a large extent, reflects structural properties of Γ . An important question is therefore to determine whether the eigenvalues of the graph Laplacian bear any relation to the spectrum of the Laplace operator $\mathcal{H} = -\frac{d^2}{dx^2}$, with Neumann–Kirchhoff conditions on the quantum graph constructed from Γ . The spectral analysis of quantum graphs is generally quite challenging (see Berkolaiko & Kuchment, 2013). A simple result is the following.

THEOREM 6.1 Let μ_j and \mathbf{q}_j , $j = 1, \dots, (n-1)M$ be the eigenvalues of the symmetric positive definite pencil $(\mathbf{H}_{11}, \mathbf{V})$, i.e.,

$$\mathbf{H}_{11}\mathbf{q}_j = \mu_j\mathbf{V}\mathbf{q}_j \quad \forall j.$$

We have that $\lambda_* \notin \left\{ \mu_j \right\}_{j=1}^{(n-1)M}$ is an eigenvalue for

$$\mathbf{H}\mathbf{x} = \lambda\mathbf{M}\mathbf{x}$$

if and only if λ_* is a root of the algebraic (rational) equation $\det \mathbf{T}(\lambda) = 0$, where

$$\mathbf{T}(\lambda) = \mathbf{H}_{22} - \lambda\mathbf{F} - (\mathbf{H}_{12} - \lambda\mathbf{C})^T (\mathbf{H}_{11} - \lambda\mathbf{V})^{-1} (\mathbf{H}_{12} - \lambda\mathbf{C}), \quad (6.2)$$

i.e., if and only if there exists $\mathbf{y} \neq \mathbf{0}$ such that $\mathbf{T}(\lambda_*)\mathbf{y} = \mathbf{0}$.

Proof. Given the generalized eigenvalue problems for the symmetric positive definite pencils (\mathbf{H}, \mathbf{M}) and $(\mathbf{H}_{11}, \mathbf{V})$, if λ is not an eigenvalue of the pencil $(\mathbf{H}_{11}, \mathbf{V})$, the matrix $\mathbf{H}_{11} - \lambda\mathbf{V}$ is nonsingular. If we partition the eigenvectors of the pencil (\mathbf{H}, \mathbf{M}) conformally to the block structure of \mathbf{H} and \mathbf{M} , we have

$$\mathbf{H}_{11}\mathbf{q}_1 + \mathbf{H}_{12}\mathbf{q}_2 = \lambda(\mathbf{V}\mathbf{q}_1 + \mathbf{C}\mathbf{q}_2), \quad (6.3)$$

$$\mathbf{H}_{12}^T\mathbf{q}_1 + \mathbf{H}_{22}\mathbf{q}_2 = \lambda(\mathbf{C}^T\mathbf{q}_1 + \mathbf{F}\mathbf{q}_2). \quad (6.4)$$

Thus, we have from (6.3) that

$$(\mathbf{H}_{11} - \lambda\mathbf{V})\mathbf{q}_1 = (\lambda\mathbf{C} - \mathbf{H}_{12})\mathbf{q}_2$$

and

$$\mathbf{q}_1 = (\mathbf{H}_{11} - \lambda\mathbf{V})^{-1} (\lambda\mathbf{C} - \mathbf{H}_{12})\mathbf{q}_2.$$

Substituting \mathbf{q}_1 in (6.4), we have

$$\left(\mathbf{H}_{12}^T (\mathbf{H}_{11} - \lambda\mathbf{V})^{-1} (\lambda\mathbf{C} - \mathbf{H}_{12}) + \mathbf{H}_{22} \right) \mathbf{q}_2 = \lambda \left(\mathbf{C}^T (\mathbf{H}_{11} - \lambda\mathbf{V})^{-1} (\lambda\mathbf{C} - \mathbf{H}_{12}) + \mathbf{F} \right) \mathbf{q}_2.$$

Rearranging terms, we obtain the desired result. \square

We note that this result is not new, being known in the context of algebraic substructuring methods for large-scale linear eigenvalue problems (see, e.g., [Bekas & Saad, 2005](#)) and in the approximation of certain resonance problems ([Bindel & Hood, 2013](#)). What is of interest here is that Theorem 6.1 has several continuous (i.e., infinite-dimensional) counterparts described in [Kuchment \(2004, 2005\)](#). For self-adjoint Hamiltonians, a λ_* which is not in the spectrum of the Hamiltonian restricted to any edge of Γ will be an eigenvalue if and only if it is a root of an algebraic equation obtained by imposing certain conditions at each vertex of Γ , just as in Theorem 6.1 above.

It is important to remark that for diffusion problems, the difference between the global eigenvalues of \mathbf{L} on the extended graph and the eigenvalues of the combinatorial graph Laplacian \mathbf{L}_Γ can make the quantum graph version of the problem more challenging and richer insofar the behavior of the solution on the graph, seen as a quantum graph, is more difficult to predict *a priori*. This is due to the possible occurrence of ‘Dirichlet eigenvalues’ associated with the edges, namely eigenvalues of the pencil $(\mathbf{H}_{11}, \mathbf{V})$, which are also eigenvalues of the pencil (\mathbf{H}, \mathbf{M}) (see Kuchment, 2004). Also, it is clear that while the spectrum of \mathbf{L}_Γ is necessarily bounded for a fixed N , the spectrum of the discrete Hamiltonian is unbounded above as h is refined and N is kept fixed, reflecting the unboundedness of the infinite-dimensional Hamiltonian \mathcal{H} .

Some numerical comparisons between the eigenvalues of \mathbf{L}_Γ and those of the simple Hamiltonian $\mathcal{H} = -\frac{d^2}{dx^2}$ for a few small graphs can be found in Section 8.2, where the behavior of the eigenfunctions of \mathcal{H} is also shown.

7. Solution of the discretized equations

In this section, we discuss the solution of problems arising from the finite element discretization of quantum graphs, namely:

- (1) Systems of linear algebraic equations $\mathbf{H}\mathbf{u}_h = \mathbf{f}_h$.
- (2) Initial value problems for systems of linear ODEs.

Problem (1) arises in particular when solving variational problems of the form (3.3) using finite element methods. Problem (2) arises from the semidiscretization of (the weak form of) parabolic PDEs on Γ :

Given $u_0 \in H^1(\Gamma)$ and $f \in L^2([0, T]; L^2(\Gamma))$ find $u \in L^2([0, T]; H^1(\Gamma)) \cap C^0([0, T]; H^1(\Gamma))$, such that

$$\begin{cases} \frac{\partial u}{\partial t} - \frac{\partial^2 u}{\partial x^2} + mu = f & \text{on } \Gamma \times [0, T], \\ u(x, 0) = u_0 & x \in \Gamma, \end{cases} \quad (7.1)$$

where $m \geq 0$ (see also Definition 2.5). We further assume that $u(\cdot, t)$ satisfies the Neumann–Kirchhoff conditions on the vertices v of Γ , for all $t \in [0, T]$.

The need for solving large linear systems on the extended graph \mathcal{G} also arises in the solution of problem (7.1) by fully implicit methods, and in the solution of generalized eigenvalue problems of the form (6.1) when shift-and-invert methods are used.

7.1 Solution of linear algebraic systems

We assume again for simplicity of notation that each edge $e \in \mathcal{E}$ contains the same number $n - 1$ of internal nodes. For solving linear systems of the form $\mathbf{H}\mathbf{u}_h = \mathbf{f}_h$ corresponding to the extended graph \mathcal{G} , we make use of the block LU factorization

$$\mathbf{H} = \begin{bmatrix} \mathbf{H}_{11} & \mathbf{H}_{12} \\ \mathbf{H}_{12}^T & \mathbf{H}_{22} \end{bmatrix} = \begin{bmatrix} \mathbf{H}_{11} & \mathbf{O} \\ \mathbf{H}_{12}^T & \mathbf{S} \end{bmatrix} \begin{bmatrix} \mathbf{I}_{(n-1)M} & \mathbf{H}_{11}^{-1}\mathbf{H}_{12} \\ \mathbf{O} & \mathbf{I}_N \end{bmatrix}. \quad (7.2)$$

We recall that when the potential $v(x)$ is positive on the (metric) graph Γ , the matrix \mathbf{H} is guaranteed to be positive definite. In particular, both \mathbf{H}_{11} and the Schur complement

$$\mathbf{S} = \mathbf{H}_{22} - \mathbf{H}_{12}^T \mathbf{H}_{11}^{-1} \mathbf{H}_{12}$$

are symmetric and positive definite (SPD).

The factorization (7.2) corresponds to the following block elimination procedure frequently used in domain decomposition algorithms. First, the following $N \times N$ reduced linear system is solved:

$$\mathbf{S} \mathbf{u}_h^v = \mathbf{f}_h^v - \mathbf{H}_{12}^T \mathbf{H}_{11}^{-1} \mathbf{f}_h^e \equiv \mathbf{c}_h, \quad (7.3)$$

where \mathbf{u}_h^v and \mathbf{f}_h^v are the values of the discrete solution \mathbf{u}_h and external load \mathbf{f}_h at the vertices $v \in \mathcal{V}$, and \mathbf{f}_h^e are the values of \mathbf{f}_h on the nodes internal to the edges $e \in \mathcal{E}$.

Next, the values of the solution \mathbf{u}_h^e at the internal nodes are obtained by solving the $(n_e - 1)M \times (n_e - 1)M$ linear system

$$\mathbf{H}_{11} \mathbf{u}_h^e = \mathbf{f}_h^e - \mathbf{H}_{12} \mathbf{u}_h^v, \quad (7.4)$$

where \mathbf{f}_h^e is the vector containing the values of \mathbf{f}_h at the nodes internal to the edges.

As already observed, the coefficient matrix \mathbf{S} of (7.3) is sparse, and it can be constructed explicitly. The resulting linear system can be solved either by a direct solver (sparse Cholesky factorization) or by an iterative method, such as PCG. In the case of PCG, the explicit assembling of the Schur complement can be avoided. If the M diagonal blocks of \mathbf{H}_{11} are factored at the outset, then each matrix–vector product with \mathbf{S} involves, at each PCG iteration, a diagonal scaling (with \mathbf{H}_{22}), two sparse matrix–vector products (with \mathbf{H}_{12} and \mathbf{H}_{12}^T) and M independent tridiagonal solves (using the precomputed factorizations) at each iteration. Although this procedure is more expensive than a matrix–vector product with the assembled \mathbf{S} , it has the advantage that it can be more easily performed in parallel in a distributed environment.

The linear system (7.4), which consists of M completely uncoupled tridiagonal linear systems, can be solved in $O(n_e M)$ time (or less if parallelism is exploited), essentially the cost of one CG iteration on (7.3) if the Schur complement is not explicitly assembled. Note that in practice, moderate values of n_e may suffice for a sufficiently accurate solution. On the other hand, M may be very large for a large graph Γ .

Clearly, the critical step is the solution of the reduced system (7.3). A sparse Cholesky factorization may be appropriate if N is not too large, but one should keep in mind that in the case of complex graphs such as small-world and power law graphs (see Newman, 2010), the Cholesky factor will frequently incur enormous amounts of fill-in, regardless of the ordering used. This may discourage the use of a sparse direct solver even for moderate values of N ; see Benzi & Kuhlemann (2013) for an example from computational biology. Thus, here we focus instead on iterative solvers, particularly on PCG.¹

As is well known, the rate of convergence of the CG algorithm depends on the distribution of the eigenvalues of the coefficient matrix, and the key to rapid convergence is the choice of an appropriate preconditioner. A wide variety of algebraic preconditioners is available, including preconditioners based on splittings of the coefficient matrix, incomplete factorizations and algebraic multigrid (AMG) methods

¹ Chebyshev semi-iteration is also quite attractive for problems, where bounds on the extreme eigenvalues of \mathbf{S} are available (see Benzi & Kuhlemann, 2013).

(Benzi, 2002; Saad, 2003). With relatively simple graph topologies, incomplete Cholesky preconditioning or AMG can be expected to give good results; however, for such problems, a sparse direct solver may be the best choice, especially if the original graph Γ is planar. Generally, in the case of large, complex graphs, however, where direct solvers are not an option, incomplete Cholesky factorization turns out to be not competitive (Benzi & Kuhlemann, 2013). It turns out, however, that when Γ is a graph with a power law degree distribution, information about the eigenvalue distribution of the graph Laplacian \mathbf{L}_Γ is available, which suggests the use of a simple diagonal preconditioner. In particular, it is known (Chung & Lu, 2006) that under certain conditions on Γ , the nonzero eigenvalues of the *normalized Laplacian*

$$\widehat{\mathbf{L}}_\Gamma = \mathbf{D}^{-\frac{1}{2}} \mathbf{L}_\Gamma \mathbf{D}^{-\frac{1}{2}} = \mathbf{I}_N - \mathbf{D}^{-\frac{1}{2}} \mathbf{A} \mathbf{D}^{-\frac{1}{2}}$$

lie in a small interval $(1 - \delta, 1 + \delta)$ with $0 < \delta < 1$ independent of N . This guarantees that diagonally preconditioned CG will converge rapidly, with a rate that is independent of the number of vertices N in Γ , when applied to linear systems involving the graph Laplacian, and we expect a similar behavior when solving systems of the form (7.3) whenever \mathbf{S} is in some sense ‘close’ to \mathbf{L}_Γ . We will return to this topic in Section 8.

Another simple preconditioner that can be easily implemented is the first-degree polynomial preconditioner given by

$$\mathbf{P}^{-1} = \mathbf{D}^{-1} + \mathbf{D}^{-1} (\mathbf{D} - \mathbf{S}) \mathbf{D}^{-1} \approx \mathbf{S}^{-1}. \quad (7.5)$$

This approximation is obtained from the identity

$$\mathbf{S}^{-1} = (\mathbf{I}_N - \mathbf{D}^{-1} (\mathbf{D} - \mathbf{S}))^{-1} \mathbf{D}^{-1}$$

by truncating to first order the Neumann series expansion

$$(\mathbf{I}_N - \mathbf{D}^{-1} (\mathbf{D} - \mathbf{S}))^{-1} = \sum_{k=0}^{\infty} (\mathbf{D}^{-1} (\mathbf{D} - \mathbf{S}))^k.$$

This preconditioner has been found to be effective in solving certain Laplacian-type linear systems stemming from complex network analysis, see Salkuyeh *et al.* (2014). It is only slightly more expensive than the diagonal one. For a typical sparse problem averaging about five nonzeros per row in \mathbf{S} , the cost of each PCG iteration is about 50% higher than with diagonal scaling, but convergence is faster. In Section 8, we present the results of numerical experiments using diagonal and polynomial preconditioning on a few different types of graphs.

7.2 Solution of systems of ODEs

Semidiscretization of the weak form of (7.1) leads to initial value problems of the form

$$\mathbf{M} \dot{\mathbf{u}}_h = -\mathbf{H} \mathbf{u}_h + \mathbf{f}_h \quad \mathbf{u}_h(0) = \mathbf{u}_{h,0}, \quad (7.6)$$

where the dot denotes differentiation with respect to time. Let us consider for simplicity the case of the *heat equation* with $f = 0$ and m independent of time. Then, $\mathbf{f}_h = \mathbf{0}$ and the solution of (7.6) can be given explicitly in terms of matrix exponentials:

$$\mathbf{u}_h(t) = \exp(-t \mathbf{M}^{-1} \mathbf{H}) \mathbf{u}_{h,0} \quad \forall t \in [0, T]. \quad (7.7)$$

Approximations to the exact semidiscrete solution (7.7) can be obtained in various ways. Here, we consider time-stepping methods and exponential integrators based on Krylov subspace approximations. Besides explicit methods, which we do not discuss here due to the well-known restrictions on the time step required for stability, the most commonly used time-stepping methods are backward Euler and Crank–Nicolson. These methods require at each step the solution of linear systems with coefficient matrices

$$\mathbf{I} + \Delta t \mathbf{M}^{-1} \mathbf{H} \quad \text{and} \quad \mathbf{I} + \frac{\Delta t}{2} \mathbf{M}^{-1} \mathbf{H}, \quad (7.8)$$

respectively, where $\Delta t > 0$ is the time step.

On the other hand, Krylov-based exponential integrators compute directly, for any value of $t \in (0, T]$, an approximation of (7.7) by projecting the problem on a suitable low-dimensional subspace and forming the exponential of the (small) projected matrix explicitly. This is equivalent to a polynomial approximation of (7.7). Both Lanczos- and Arnoldi-type methods can be used, depending on whether the underlying problem is symmetric (Hermitian in the complex case) or not. Here, we note that the matrix $\mathbf{M}^{-1} \mathbf{H}$, which occurs in both (7.7) and (7.8), is generally nonsymmetric, apparently precluding the use of the more efficient Lanczos-based methods. However, since both \mathbf{H} and \mathbf{M} are SPD, the product $\mathbf{M}^{-1} \mathbf{H}$ is symmetrizable. For example, if $\mathbf{M} = \mathbf{R}^T \mathbf{R}$ is the Cholesky factorization of \mathbf{M} with \mathbf{R} upper triangular, introducing the new variable $\mathbf{v}_h = \mathbf{R} \mathbf{u}_h$ leads to the solution vector in (7.7), becoming

$$\mathbf{u}_h(t) = \mathbf{R}^{-1} \mathbf{v}_h(t) \quad \text{where} \quad \mathbf{v}_h(t) = \exp(-t \tilde{\mathbf{H}}) \mathbf{v}_{h,0} \quad \forall t \in [0, T], \quad (7.9)$$

with $\tilde{\mathbf{H}} = \mathbf{R}^{-T} \mathbf{H} \mathbf{R}^{-1}$ and $\mathbf{v}_{h,0} = \mathbf{R} \mathbf{u}_{h,0}$. Here $\tilde{\mathbf{H}}$ is symmetric, and Lanczos-based methods can be applied to approximate the solution. A similar symmetrization applied to the matrices in (7.8) allows one to use the PCG method to solve the linear systems arising from implicit methods, rather than the more expensive nonsymmetric Krylov methods. The price to pay for this is the need to compute the Cholesky factorization of \mathbf{M} and to perform triangular solves at each step or iteration. Unfortunately, for large graphs, the Cholesky factorization $\mathbf{M} = \mathbf{R}^T \mathbf{R}$ of the mass matrix can be prohibitive when the integrals that define \mathbf{M} are computed with the Simpson rule. In this case, however, it is possible to replace the mass matrix \mathbf{M} with a diagonal approximation, $\tilde{\mathbf{M}} \approx \mathbf{M}$, obtained, for example, by *lumping*. Symmetrization is then trivial, at the expense of an additional error, which, however, is $O(h)$ as $h \searrow 0$. The same holds true if the integrals in the mass matrix are approximated with the simple trapezoidal rule, in which case \mathbf{M} is diagonal.

In Section 8.3, we discuss the results of experiments comparing the two quadrature rules, where we apply a Krylov-based method (Afanasjew *et al.*, 2008; Güttel, 2017) to approximate the solution of a simple diffusion problem on different types of graphs.

8. Numerical experiments

In this section, we illustrate the results of numerical studies, including the solution of simple elliptic and parabolic PDEs and eigenvalue problems on graphs.

8.1 Solution of simple elliptic problems

We begin by showing some results of experiments using the PCG method to solve linear systems arising from the discretization of simple elliptic problems posed on scale-free graphs obtained using the

TABLE 1 *Number of vertices and edges for three scale-free graphs and in the corresponding extended graphs*

Graph	N	M	m	p
Γ_1	2000	3974	81480	79480
Γ_2	5000	9968	204360	199360
Γ_3	10000	19965	409300	399300

Barabási–Albert model (Barabasi & Albert, 1999). We are especially interested in seeing how PCG iterations scale with problem size. We consider two main situations that lead to linear systems of increasing size:

- The mesh size h is fixed, but the size of the graph Γ increases.
- The graph Γ is fixed, but h decreases.

In the first situation, we assume that the graph’s average degree is kept roughly constant. In the second situation, we are applying PCG to a reduced system of fixed size N , but *a priori* the number of iterations could grow since the condition number (more precisely, the eigenvalue distribution) of the Schur complement may worsen as $h \searrow 0$.

We use the Matlab toolkit CONTEST (Taylor & Higham, 2009) to generate scale-free graphs with different numbers N of vertices and M of edges, while keeping the average degree constant. In Table 1, we report the sizes of three graphs Γ generated using CONTEST together with the number of vertices ($m = (n_e - 1)M + N$), and number of edges ($p = n_e M$) of the corresponding extended graphs \mathcal{G} for the constant choice $h = \frac{1}{21}$ of the mesh size. All edges are assumed to have unit length. In particular, we do not consider here adaptive discretizations that take into account the presence of hubs (see Section 3.1), although this is not difficult to do. We recall that m is the order of the discrete Hamiltonian \mathbf{H} , whereas N is the order of the corresponding Schur complement \mathbf{S} . For these problems, the Schur complement has only about four nonzeros per row.

For the Hamiltonian, we consider the simple elliptic operator

$$\mathcal{H} u = \left(-\frac{d^2}{dx^2} + v \right) u, \quad (8.1)$$

where v is constant, with Neumann–Kirchhoff conditions at the vertices. The discrete Hamiltonian \mathbf{H} obtained applying one-dimensional linear finite elements to the weak form of (8.1) is then SPD for all $v > 0$. We choose the right-hand side $\mathbf{f}_h = \mathbf{e}_1$, corresponding to a unit load applied to the first-numbered vertex v_1 of the graph Γ .

In Table 2, we report iteration counts for the conjugate gradient method with no preconditioning, with diagonal preconditioning and with the polynomial preconditioner (7.5) for $v = 0.1$, using the three graphs of Table 1. In all cases, the initial guess is the zero vector and the stopping criterion is a reduction of the relative residual norm below $\sqrt{\text{eps}}$, where $\text{eps} \approx 2.2204 \cdot 10^{-16}$. In all cases, the relative 2-norm of the error is of the order of $\sqrt{\text{eps}} \approx 10^{-8}$.

TABLE 2 *PCG iteration counts for Schur complement system, different preconditioners ($\nu = 0.1$)*

Graph	No prec.	Diagonal	Polynomial
Γ_1	78	28	15
Γ_2	102	28	15
Γ_3	115	28	15

TABLE 3 *PCG iteration counts for Schur complement system, different values of h . Graph: Γ_1 ($\nu = 0.1$)*

h^{-1}	m	No preconditioning	Diagonal	Polynomial
21	81480	78	28	15
41	160960	78	28	15
81	319920	78	28	15
101	399400	78	28	15

From the results, one can see that, while the number of iterations increases without preconditioning for increasing graph size, it remains constant (and quite small) with both diagonal and polynomial preconditioning. Qualitatively similar results are observed for other values of ν , with the convergence being faster as ν increases.

Next, we fix the metric graph (using Γ_1) and we refine the discretization of the edges. In Table 3, we report iteration counts for four different values of h , corresponding to $n = 20, 40, 80, 100$ equally spaced nodes per edge. We also report the order m of the discrete Hamiltonian, \mathbf{H} . The order of the Schur complement \mathbf{S} is fixed ($N = 2000$). Although this Schur complement is so small that a sparse direct solver suffices, we are interested in the behavior of the PCG iteration as a function of h , expecting a qualitatively similar behavior for larger graphs. The results show that the convergence of the conjugate gradient algorithm, even in the absence of preconditioning, is completely h -independent, suggesting that even for rather fine meshes, the Schur complement is close to a well-conditioned matrix. As before, convergence is faster (slower) if ν is taken larger (respectively, smaller).

We conclude that for this class of quantum graphs, the reduced system approach, combined with diagonally preconditioned CG for solving the Schur complement system, results in rates of convergence that are both h - and N -independent, and thus it is *optimal*, in the sense that the total solution cost scales linearly in the number $m = (n_e - 1)M + N$ of degrees of freedom. This approach has also very high inherent parallelism, especially if the Schur complement is not explicitly assembled (except for the diagonal entries of \mathbf{S} if they are needed). Polynomial preconditioning may lead to slightly less work overall, but the difference is small.

The observed convergence behavior can be explained as follows. For $\nu > 0$ and $h > 0$ fixed and sufficiently small, the Schur complement \mathbf{S} is, up to a constant factor, a small perturbation of the combinatorial graph Laplacian, \mathbf{L}_Γ (see the appendix). Hence, we expect the convergence behavior of PCG applied to the Schur complement system to be close to that of PCG applied to a (consistent) linear system of the form $\mathbf{L}_\Gamma \mathbf{x} = \mathbf{b}$. We remark that, although this system is singular, the singularity is benign, the kernel being one-dimensional and spanned by a known vector since Γ is connected. In particular, the eigenvalue

$\lambda_1 = 0$ of \mathbf{L}_Γ plays no role in determining the rate of convergence of the conjugate gradient method. Now, the distribution of the extreme (nonzero) eigenvalues of the Laplacian of scale-free graphs is known. Quite a lot is known also about the nonzero eigenvalues of the normalized Laplacian $\widehat{\mathbf{L}}_\Gamma = \mathbf{D}^{-\frac{1}{2}} \mathbf{L}_\Gamma \mathbf{D}^{-\frac{1}{2}}$. Of course, since scale-free graphs obtained using the Barabási–Albert model have an element of randomness, these results are to be taken in a probabilistic sense.

We first consider the case of the normalized Laplacian. Note that, since the matrix $\mathbf{D}^{-\frac{1}{2}} \mathbf{A} \mathbf{D}^{-\frac{1}{2}}$ is symmetric and stochastic, the eigenvalues of $\widehat{\mathbf{L}}_\Gamma = \mathbf{I}_N - \mathbf{D}^{-\frac{1}{2}} \mathbf{A} \mathbf{D}^{-\frac{1}{2}}$ lie in the interval $[0, 2]$. Moreover, for scale-free graphs with sufficiently large minimum degree d_{\min} , the nonzero eigenvalues of the normalized Laplacian can be expected to fall with high probability for $N \rightarrow \infty$ in the interval

$$I = \left(1 - \frac{2}{\sqrt{w}}, 1 + \frac{2}{\sqrt{w}}\right), \quad \text{where } w = d_{\text{avg}}. \quad (8.2)$$

Here, d_{avg} denotes the average expected degree for Γ . See (Chung & Lu, 2006, Chapter 9) for the precise statement of this result. While the assumption on the minimum degree is quite restrictive, the conclusions of this theorem appear to hold in practice even for power law random graphs such as those considered here, for which d_{\min} is rather small. Our three power law graphs satisfy

$$0.83 < \frac{2}{\sqrt{w}} < 0.84,$$

so we expect the nontrivial eigenvalues of the normalized Laplacian to lie between

$$1 - \frac{2}{\sqrt{w}} \approx 0.16 \quad \text{and} \quad 1 + \frac{2}{\sqrt{w}} \approx 1.84. \quad (8.3)$$

Hence, the *effective condition number* $\kappa_2^{\text{eff}}(\widehat{\mathbf{L}}_\Gamma)$, defined as the ratio of the largest and the smallest nontrivial eigenvalues of $\widehat{\mathbf{L}}_\Gamma$, can be expected to satisfy

$$\kappa_2^{\text{eff}}(\widehat{\mathbf{L}}_\Gamma) \leq 1.84/0.16 = 11.5,$$

independently of N . Therefore, the conjugate gradient method with diagonal scaling can be expected to converge rapidly with a rate independent of N , which is what we observe in practice. Although this argument is not rigorous, we found in practice that it gives reliable estimates of the condition number of the normalized Laplacian of random power law graphs, and thus of the rate of convergence of PCG applied to the Schur complement system.

Now we turn to the case of the (unnormalized) Laplacian \mathbf{L}_Γ . Using, for example, Ostrowski's Theorem (a quantitative version of Sylvester's Law of Inertia, see Horn & Johnson, 2012, Theorem 4.5.9), one can easily show that the first nonzero eigenvalue of \mathbf{L}_Γ is related to the first nonzero eigenvalue of $\widehat{\mathbf{L}}_\Gamma$ by the following inequality:

$$\lambda_2(\mathbf{L}_\Gamma) \geq d_{\min} \cdot \lambda_2(\widehat{\mathbf{L}}_\Gamma), \quad (8.4)$$

where again d_{\min} denotes the minimum degree of any vertex in Γ . In our case, $d_{\min} = 2$ (by construction, see Taylor & Higham, 2009). It follows from (8.3) and (8.4) that for any of our three graphs Γ_i , the

smallest nonzero eigenvalue can be expected to satisfy

$$\lambda_2(\mathbf{L}_\Gamma) \geq 0.32, \quad \text{independent of } N. \quad (8.5)$$

Of course, the same holds for any power law graph with the same minimum and average degree. To four decimals, the smallest nonzero Laplacian eigenvalue for the three graphs used in our experiments was given by

- $\lambda_2 = 0.5259$ for Γ_1
- $\lambda_2 = 0.5338$ for Γ_2
- $\lambda_2 = 0.5257$ for Γ_3 .

Given the random nature of these graphs, we repeated the calculation several times, always getting similar values. Hence, although the lower bound (8.5) is slightly pessimistic, it does correctly predict that the smallest nonzero eigenvalue of a random power law graph can be expected to remain bounded away from zero as N increases. Since we saw (cf. Table 2) that without preconditioning the number of PCG iterations required to solve the Schur complement system increases with N , we predict that the largest eigenvalues of \mathbf{S} must increase with N , all else being constant. Once again, we replace \mathbf{S} with the graph Laplacian \mathbf{L}_Γ , for which analytical results are available. In [Elsässer \(2006\)](#), the following remarkable and *a priori* unexpected result is proved: In a power law graph, the upper portion of the spectrum of \mathbf{L}_Γ is distributed like the largest degrees of vertices of Γ (we refer to [Elsässer, 2006](#) for the precise statement).

Looking at the three graphs Γ_i ($i = 1, 2, 3$) used in our test, we obtained the following results. For Γ_1 , the five largest degrees are

$$185, 66, 63, 54, 45,$$

and the five largest eigenvalues of \mathbf{L}_{Γ_1} are approximately

$$186.03, 67.07, 64.06, 55.19, 45.99.$$

For Γ_2 , the five largest degrees are

$$227, 122, 90, 89, 88,$$

and the five largest eigenvalues of \mathbf{L}_{Γ_2} are approximately

$$228.03, 123.06, 91.14, 90.07, 89.01.$$

For Γ_3 , the five largest degrees are

$$430, 237, 147, 128, 98,$$

and the five largest eigenvalues of \mathbf{L}_{Γ_3} are approximately

$$431.02, 238.02, 148.04, 129.08, 99.00.$$

This shows that the theory developed in [Elsässer \(2006\)](#) is remarkably accurate. Hence, without preconditioning, the effective condition number $\kappa_2^{\text{eff}}(\mathbf{L}_\Gamma)$ of the Laplacian grows like $O(d_{\max})$ as $N \rightarrow \infty$, where d_{\max} denotes the maximum degree. As shown in the appendix, if $\nu < h^{-1}$ then in the limit for $\nu, h \rightarrow 0$ the Schur complement matrix \mathbf{S} reduces to the graph Laplacian \mathbf{L}_Γ . Hence, for ν and h sufficiently small (but fixed), we expect the iteration count of unpreconditioned CG applied to the Schur complement system to grow with N , as observed in our experiments. Fortunately, a simple diagonal scaling is sufficient to remove this dependency on N .

We stress the fact that the optimality of the Schur complement reduction approach with diagonally scaled CG is a phenomenon that has no analog in the usual two-dimensional or three-dimensional PDE setting.

8.2 Eigenvalues and eigenfunctions

It is instructive to compare numerically the eigenvalues of \mathbf{L}_Γ with the first few eigenvalues of the Hamiltonian $\mathcal{H} = -\frac{d^2}{dx^2}$ on the metric graph based on Γ and to investigate the behavior of the corresponding eigenfunctions. In order to be able to visualize the eigenfunctions, we consider at first only a few small, simple graphs: a cross (or star graph) with five vertices, a simple graphene-like structure with 12 vertices and a tree with 16 vertices. In all cases, we assume the edges have unit length, and we discretize the eigenvalue problem using linear finite elements with 100 internal discretization nodes on each edge, leading to a generalized eigenvalue problem of the form (6.1). Throughout this section, we assume Neumann–Kirchhoff conditions.

In Table 4, we report the eigenvalues of the graph Laplacian \mathbf{L}_Γ for each graph.

In Fig. 8, we display the six smallest eigenpairs obtained approximating $-u'' = \lambda u$ on the cross graph via (6.1). Note that the first nonzero eigenvalue $\lambda_2 = 2.4675$ has multiplicity three.

In Fig. 9, we display the computed approximations to the 10 smallest eigenpairs for the small graphene graph consisting of two hexagons connected by an edge. Note the peculiar behavior of the eigenfunctions. More complex graphene-like graph models show similar behavior, but, unfortunately, the cluttered displays become difficult to visualize.

Finally, we display the results for the same simple Hamiltonian $\mathcal{H}u = -u''$ on a binary tree with an extra vertex connected to the root. In Fig. 10, we display the first eight leftmost eigenpairs and, in Fig. 11, the next eight eigenpairs. As in the previous examples, we note the existence of repeated eigenvalues (reflecting the symmetry of the underlying graphs) and the fact that, as expected, the eigenfunctions corresponding to higher eigenvalues become increasingly oscillatory.

We note that for these small, highly regular graphs, the smallest eigenvalues of the (discretized) Hamiltonian tend to mimic the behavior of the eigenvalues of the graph Laplacian \mathbf{L}_Γ . It is therefore natural to ask whether this is true in general.

To answer this question, we performed some eigenvalue computations on a few graphs with irregular, complex topologies. For the experiments we used three graphs, one synthetic and the other two representing real-world networks. The first is a Barabási–Albert graph (preferential attachment model) with 2000 vertices and 3974 edges, of the same kind as the graph Γ_1 described in Table 1. The second one is a graph describing a social network of drug users in Colorado Springs, where the edges indicate pairs of individuals that have exchanged needles in the last three months ([Estrada, 2011](#), p. 122). The third graph represents protein–protein interactions in beer yeast ([Newman, 2010](#), p. 89). We denote the three networks Γ_1 , Γ_d and Γ_y , respectively. The number of vertices/edges for the networks Γ_d and Γ_y are 616/2012 and 2224/6609, respectively.

TABLE 4 *Eigenvalues of \mathbf{L}_Γ for each of the three graphs*

Cross	Graphene	Tree
0	0	0
1	0.1578	0.0968
1	1	0.2679
1	1	0.2679
5	1	0.4965
	1.4931	1
	3	1
	3	1
	3	1
	3.5069	1.7356
	4	2.1939
	4.8422	3.5767
		3.7321
		3.7321
		4.7093
		5.1912

In Table 5, we report the six smallest eigenvalues of the graph Laplacian \mathbf{L}_Γ for each of these three graphs, together with approximations to the six smallest eigenvalues of the simple Hamiltonian $\mathcal{H} = -\frac{d^2}{dx^2}$, with Neumann–Kirchhoff conditions obtained by solving the discrete problem (6.1). Again, each edge is assumed to have unit length, and 100 interior nodes are used to discretize \mathcal{H} on each edge. The eigenvalues appear to have converged to an approximation with at least three accurate significant digits. These results show that it is difficult to make general statements, as the relationship between the small eigenvalues of \mathcal{H} and those of \mathbf{L}_Γ appears to be very much graph dependent. This fact confirms the observation made in Section 6 that the diffusion dynamics could be rather different for the combinatorial graph Γ and for the corresponding metric graph, even for very simple PDEs. This is discussed further in the next section.

8.3 Parabolic problems

Here, we consider approximating the simple diffusion-type equation

$$\frac{\partial u}{\partial t} = \frac{\partial^2 u}{\partial x^2} \quad (x, t) \in \Gamma \times (0, T]; \quad u(x, 0) = u_0(x) \quad x \in \Gamma, \quad (8.6)$$

where $u_0(x)$ is an initial temperature distribution. As always in this article, Neumann–Kirchhoff conditions are imposed at the vertices. Note that as $t \rightarrow \infty$, the solution $u(x, t)$ of (8.6) must decay to a steady state, everywhere constant solution at an exponential rate.

Finite element discretization of the second derivative in (8.6) leads to a system of linear ODEs of the form (7.6) with $\mathbf{f}_i = \mathbf{0}$, the solution of which is given by (7.7) with $\mathbf{H} = \mathbf{L}$. For the initial temperature

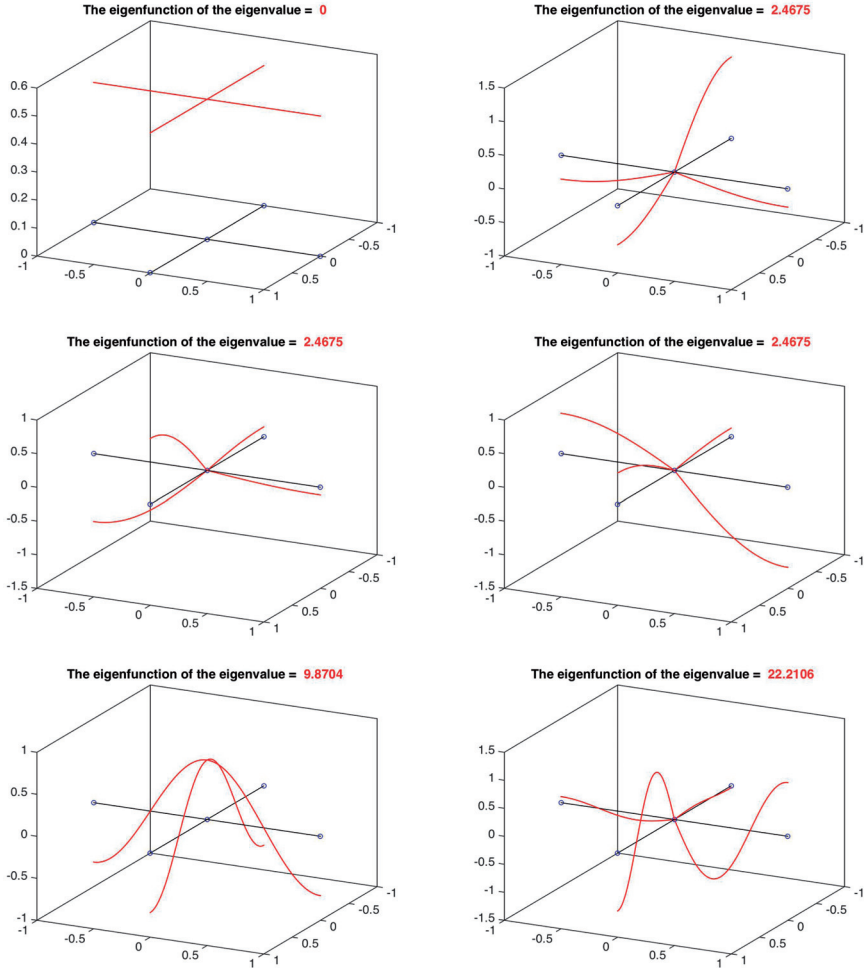


FIG. 8. Eigenvalues–eigenfunctions for $-u'' = \lambda u$ on a star graph ($n = 100$ internal points).

distribution, we have chosen the function $u_0(x)$ that is linear on each edge and such that $u_0(x|_e) = x$. Moreover, u_0 is normalized so that $\|u_0\|_{L^2(\Gamma)} = 1$.

We have used the package in Güttel (2017) to approximate the action of the matrix exponential in (7.7) at times $t_1 = 0.0002$, $t_2 = 0.0004$, \dots , $t_{10} = 0.002$ for four different graphs: the graphs Γ_1 , Γ_d and Γ_y from the previous section and a graphene-like lattice Γ_g consisting of 200 contiguous hexagons. Each edge in each graph is assumed to have unit length, and 20 interior nodes are used to discretize each edge. In the case of the graphene-like graph, which has 840 vertices and 1210 edges, the extended graph has 25040 nodes. Because the steady state is approached quickly for all graphs, we limit ourselves to a very small time interval; slower decay can be obtained by premultiplying the diffusion operator in (8.6) by a small diffusivity coefficient, but doing so does not change the relative rate of decay obtained for different graphs.

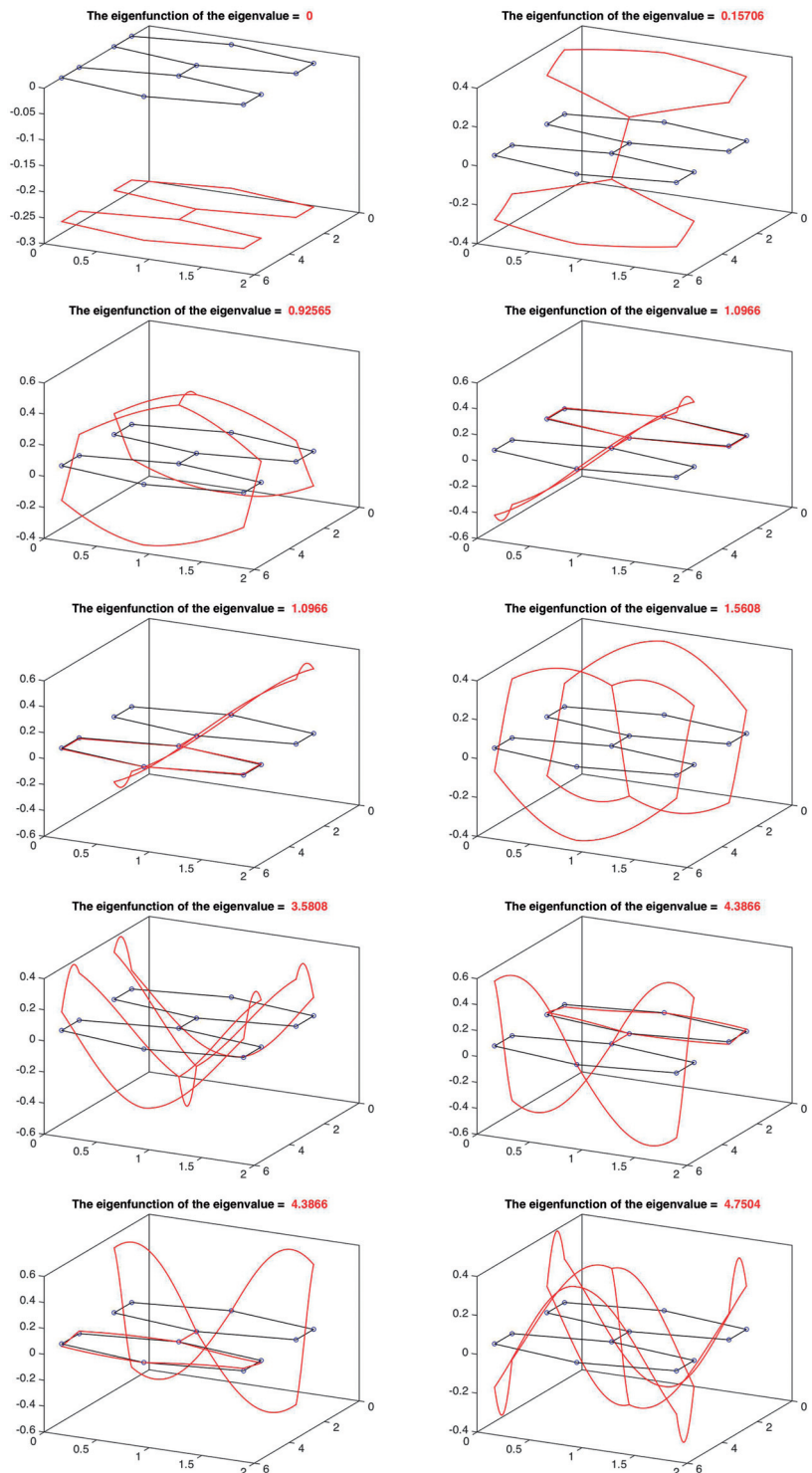


FIG. 9. Eigenvalues–eigenfunctions for $-u'' = \lambda u$ on a simple graphene graph ($n = 100$ internal points).

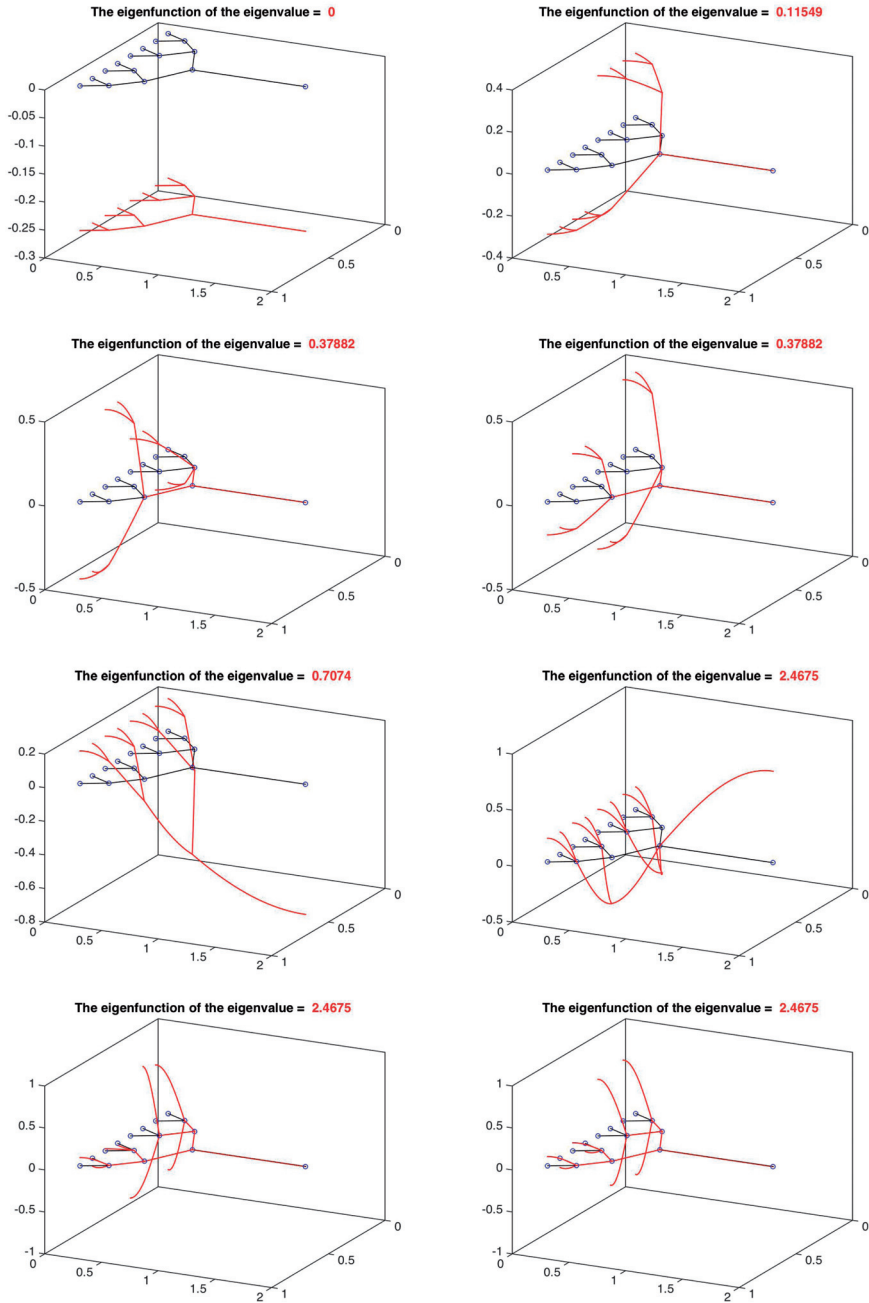


FIG. 10. Eigenvalue–eigenfunctions for $-u'' = \lambda u$ on a simple tree ($n = 100$ internal points).

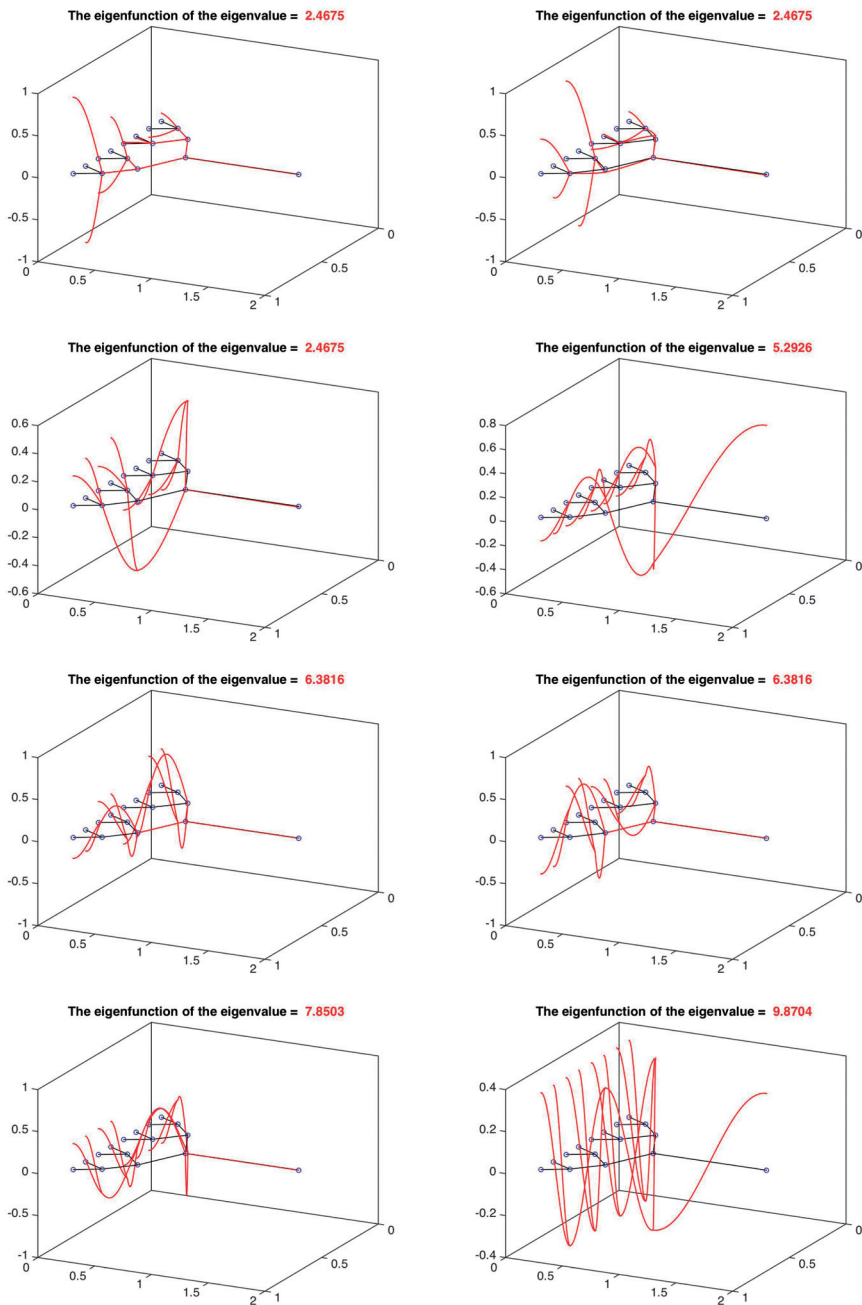


FIG. 11. Eigenvalues–eigenfunctions for $-u'' = \lambda u$ on a simple tree ($n = 100$ internal points).

TABLE 5 *The smallest six eigenvalues of \mathbf{L}_Γ and of $\mathcal{H} = -\frac{d^2}{dx^2}$ for each of the three graphs*

Eigenvalue	Γ_1		Γ_d		Γ_y	
	\mathbf{L}_Γ	\mathcal{H}	\mathbf{L}_Γ	\mathcal{H}	\mathbf{L}_Γ	\mathcal{H}
λ_1	0	0	0	0	0	0
λ_2	0.5227	0.3374	0.0107	0.0035	0.0600	0.0664
λ_3	0.5335	0.3453	0.0148	0.0073	0.0727	0.0765
λ_4	0.5431	0.3486	0.0317	0.0095	0.0904	0.0826
λ_5	0.5502	0.3522	0.0410	0.0157	0.1177	0.0909
λ_6	0.5572	0.3559	0.0617	0.0370	0.1226	0.0918

Consider the discrete solution (7.7). The rate of decay to steady state is governed primarily by the first nonzero eigenvalue of the matrix pencil (\mathbf{H}, \mathbf{M}) (i.e., of the matrix $\mathbf{M}^{-1}\mathbf{H}$). To see this, let m be the total number of nodes on the extended graph \mathcal{G} and $\mathbf{q}_1, \mathbf{q}_2, \dots, \mathbf{q}_m$ be the normalized eigenvectors corresponding to the generalized eigenvalues $\lambda_1 < \lambda_2 \leq \dots \leq \lambda_m$. Then, the solution (7.7) is given, for all times $t \geq 0$, by

$$\mathbf{u}_h(t) = e^{-t\lambda_1}(\mathbf{q}_1^T \mathbf{u}_{h,0})\mathbf{q}_1 + e^{-t\lambda_2}(\mathbf{q}_2^T \mathbf{u}_{h,0})\mathbf{q}_2 + \dots + e^{-t\lambda_m}(\mathbf{q}_m^T \mathbf{u}_{h,0})\mathbf{q}_m. \quad (8.7)$$

In our case $\mathbf{H} = \mathbf{L}$, hence $\lambda_1 = 0$ and $\mathbf{q}_1 = \frac{1}{\sqrt{m}}\mathbb{1}$ (where $\mathbb{1}$ is the vector of all ones), therefore $\mathbf{u}_h(t) \rightarrow \frac{1}{m}(\mathbb{1}^T \mathbf{u}_{h,0})\mathbb{1}$ as $t \rightarrow \infty$. In other words, the solution tends to a state of thermal equilibrium, where the temperature is the same everywhere on \mathcal{G} and is given by the space average of the initial solution. It is also clear from (8.7) that the rate of convergence depends primarily on the magnitude of the smallest nonzero eigenvalue, λ_2 , since all the terms corresponding to larger eigenvalues tend to zero faster.

The six smallest eigenvalues of the Hamiltonian for the graphs Γ_1 , Γ_d and Γ_y have been reported in Table 5. For the graphene-like graph Γ_g , the corresponding values are $\lambda_1 = 0$, $\lambda_2 = 0.0014$, $\lambda_3 = 0.0056$, $\lambda_4 = 0.0126$, $\lambda_5 = 0.0160$ and $\lambda_6 = 0.0176$. Hence, one would expect the convergence to steady state to take longer on the graphene-like graph than on the other graphs, with the Barábasi–Albert graph Γ_1 exhibiting the fastest rate of decay ($\lambda_2 = 0.3374$). This is also intuitive in view of the fact that Γ_g is a two-dimensional lattice with large diameter, whereas the Barábasi–Albert graph is a small-world graph with small diameter, and so diffusion should take place faster on the latter graph than on the former (with the other two graphs occupying somewhat intermediate positions).

However, things are not quite so simple. It is important to keep in mind that we are dealing with solutions to a partial differential equation, and that the actual decay behavior may be different when measured by different, nonequivalent norms. Recall (see, e.g., Brezis, 2010, Chapter 11) that the solution $u(t)$ of problem (8.6) satisfies at each time t the relation

$$\frac{1}{2} \|u(t)\|_{L^2(\Gamma)}^2 + \int_0^t (\mathcal{H}u(s), u(s))_{L^2(\Gamma)} ds = \frac{1}{2} \|u(0)\|_{L^2(\Gamma)}^2, \quad (8.8)$$

and if $u(t) \in C^1(\Gamma)$ for all t , we have

$$\frac{d}{dt} \left(\frac{1}{2} \|u(t)\|_{L^2(\Gamma)}^2 \right) + (\mathcal{H}u(s), u(s))_{L^2(\Gamma)} = 0 \quad (8.9)$$

(where, for ease of notation, we have suppressed the dependence of u on x). We point out that the finite element approximation gives the relations

$$\mathbf{M}\dot{\mathbf{u}}_h = -\mathbf{H}\mathbf{u}_h \quad \mathbf{u}_h(0) = \mathbf{u}_{h,0},$$

and therefore the semidiscrete solutions $\mathbf{u}_h(t)$ will satisfy for each h

$$\frac{d}{dt} \left(\frac{1}{2} \|\mathbf{u}_h(t)\|_{L^2(\Gamma)}^2 \right) + \mathbf{u}_h(t)^T \mathbf{H}\mathbf{u}_h(t) = 0. \quad (8.10)$$

We also note that in our case $(\mathcal{H}u(s), u(s))_{L^2(\Gamma)} = \|u(s)\|_{H^1(\Gamma)}^2$, the square of the H^1 (semi)norm. This quantity has the physical meaning of an energy. Since the quadratic form $\mathbf{u}_h(t)^T \mathbf{H}\mathbf{u}_h(t)$ is the discrete H^1 seminorm (squared), (8.10) implies that large energy solutions at a time t will show a faster decrease of the $L^2(\Gamma)$ norm.

In Fig. 12, we display for our test problems both the H^1 seminorm and the L^2 -norm of the solutions (squared). These plots show that the decay behavior is different in different norms. In particular, the approach to equilibrium is fastest for the graphene-like graph and slowest for the Barábasi–Albert graph when measured in the L^2 -norm, but the situation is reversed when decay is measured in terms of energy, with the decay of the H^1 seminorm now being noticeably slower for the highly regular graphene-like graph (note the semilogarithmic scale used for the energy plot). This observation is consistent with our remarks above.

As a final note, one of the purposes of these experiments was to compare the use of two variants of the mass matrix \mathbf{M} , the one obtained from the Simpson rule and the diagonal approximation of it obtained using the trapezoidal rule. To generate the results shown in Fig. 12, we made use of the trapezoidal rule. The results for the Simpson rule were found to be virtually identical and therefore are not shown. The essential equivalence of the two variants of the mass matrix is in line with what one would expect from a simple error analysis of the two quadrature rules. The Simpson rule, however, requires computing the Cholesky factorization of the mass matrix and two triangular solves at each step of the Krylov subspace method (see the discussion in Section 7.2). In contrast, the trapezoidal rule does not necessitate any of this, since it leads to a diagonal approximation of the mass matrix, resulting in computing times that are orders of magnitude smaller.

9. Conclusions and future work

In this article, we have introduced and analysed a linear finite element method for the discretization of elliptic, parabolic and eigenvalue problems posed on graphs, with special attention to the important case of Neumann–Kirchhoff vertex conditions. The structure and main properties of the resulting stiffness and mass matrices have been carefully described, together with the associated notion of *extended graph*. Numerical linear algebra aspects have been discussed, as well as the numerical integration of simple parabolic PDEs on graphs. The effect of graph topology on the convergence behavior of iterative solvers has been discussed and illustrated by numerical experiments, showing that a combination of Schur complement reduction (closely related to a nonoverlapping domain decomposition approach) with diagonally

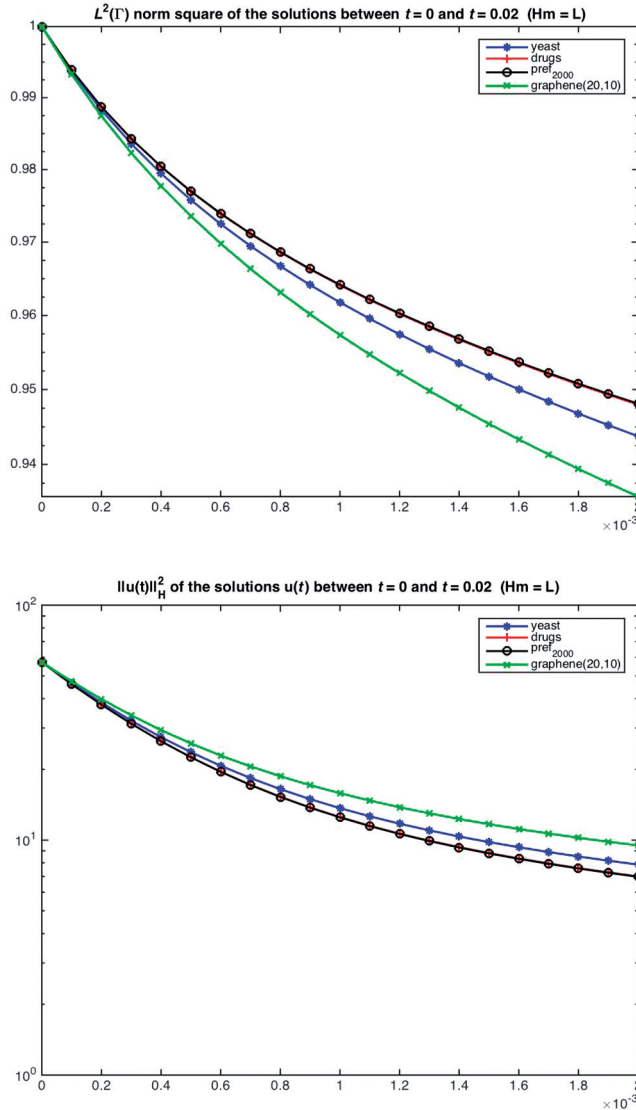


FIG. 12. $L^2(\Gamma)$ norm squared and energy of solution $u(x, t)$ to the heat equation as a function of time ($n = 20$ internal nodes). Note the logarithmic scale on the vertical axis in the second plot.

scaled CG results in optimal solvers for scale-free graphs. This approach has also very high inherent parallelism. Not surprisingly, we have found that the solution of PDEs on graphs, particularly complex networks, can lead to new phenomena that are not typically observed when solving PDEs posed on more typical two-dimensional and three-dimensional domains.

The numerical analysis of PDEs on graphs and networks is still in its infancy, and much work remains to be done in this area. Future work should address the numerical solution of more complex types of

differential equations on graphs. These include hyperbolic problems (especially nonlinear conservation laws, which are important for the description of transport phenomena on networks, as well as for the propagation of shocks), Schrödinger-type equations, systems of PDEs (such as the Dirac equations, which are important in the modeling of graphene) and nonlocal PDEs of fractional type. In particular, the influence of the underlying graph topology on the discretization and solver behavior should be investigated for these more complex PDEs.

Finally, in this article we have focused on a simple linear continuous finite element discretization, which allowed us to relate the discretized Hamiltonian to standard linear operators associated with graphs, such as the incidence matrix and the graph Laplacian. It would also be of interest to compare the efficacy of different discretization strategies (discontinuous or higher order finite elements, finite differences, finite volumes and spectral methods) in solving PDEs on graphs.

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Appendix

We assume that we approximate by uniform linear finite elements the simple Hamiltonian

$$\mathcal{H}u = -\frac{d^2u}{dx^2} + vu$$

($v = \text{constant}$), with $\ell_e = 1$ and $n_e = n \forall e \in \mathcal{E}$, so that the block

$$\mathbf{H}_{11} = \mathbf{A} + \mathbf{V} = \mathbf{I} \otimes \left(\frac{1}{h} \widehat{\mathbf{T}} + hv \widehat{\mathbf{M}} \right)$$

is the combination of the two tridiagonal matrices of order $n - 1$:

$$\widehat{\mathbf{T}} = \text{tridiag} \{-1, 2, -1\}, \quad \widehat{\mathbf{M}} = \frac{1}{6} \text{tridiag} \{1, 4, 1\},$$

the block

$$\mathbf{H}_{12} = \mathbf{B} + \mathbf{C} = \left(\frac{-1}{h} + \frac{vh}{6} \right) |\mathbf{B}|,$$

and

$$\mathbf{H}_{22} = \mathbf{G} + \mathbf{F} = \left(\frac{1}{h} + \frac{vh}{3} \right) \mathbf{D}.$$

Both the matrices $\widehat{\mathbf{T}}$ and $\widehat{\mathbf{M}}$ are diagonalized by the symmetric, orthonormal matrix \mathbf{U} (Meurant, 1992):

$$\mathbf{U} = (u_{i,j}) \quad \left\{ u_{i,j} = \sqrt{2h} \sin(ij\pi h) \right\}_{i,j=1}^{n-1}.$$

$$\mathbf{U} \widehat{\mathbf{T}} \mathbf{U} = \Lambda \quad \text{and} \quad \mathbf{U} \widehat{\mathbf{M}} \mathbf{U} = \frac{1}{6} (2\mathbf{I} + \Lambda),$$

where $\lambda_i = 4 \sin^2 \left(i \frac{\pi}{2} h \right)$. Moreover, we have

$$\mathbf{U} \mathbf{H}_{11} \mathbf{U} = \Theta = \frac{1}{h} \Lambda + \frac{hv}{6} (2\mathbf{I} + \Lambda).$$

We observe the following useful relations:

$$\sin(i\pi h) = (-1)^{i-1} \sin(in\pi h), \tag{A.1}$$

$$\sin(i\pi h) = 2 \cos \left(i \frac{\pi}{2} h \right) \sin \left(i \frac{\pi}{2} h \right), \tag{A.2}$$

$$\sin(i\pi h)^2 = \frac{1}{4} (4 - \lambda_i) \lambda_i. \tag{A.3}$$

The Schur complement \mathbf{S} of the Hamiltonian in our specific case is

$$\begin{aligned} \mathbf{S} &= \mathbf{G} + \mathbf{F} - (\mathbf{B} + \mathbf{C})^T (\mathbf{A} + \mathbf{V})^{-1} (\mathbf{B} + \mathbf{C}) \\ &= \left(\frac{1}{h} + \frac{vh}{3} \right) \mathbf{D} - \left(\frac{-1}{h} + \frac{vh}{6} \right)^2 |\widehat{\mathbf{E}} \widehat{\mathbf{E}}^T| \left(\mathbf{I} \otimes (\mathbf{U} \Theta^{-1} \mathbf{U}) \right) |\widehat{\mathbf{E}} \widehat{\mathbf{E}}^T|. \end{aligned}$$

Taking into account that

$$\bar{\mathbf{E}}^T = \mathbf{I}_M \otimes \mathbf{E}_e^T,$$

with $\mathbf{E}_e \in \mathbb{R}^{(n-1) \times n}$, and in view of (4.2), (4.3) and (4.4) with $n_{e_j} = n$, we have

$$\begin{aligned} |\widehat{\mathbf{E}}\bar{\mathbf{E}}^T| &= (\mathbf{E}^+ \otimes (\mathbf{e}_1^n)^T + |\mathbf{E}^-| \otimes (\mathbf{e}_n^n)^T)(\mathbf{I}_M \otimes |\mathbf{E}_e|^T) \\ &= (\mathbf{E}^+ \otimes (\mathbf{e}_1^{n-1})^T + |\mathbf{E}^-| \otimes (\mathbf{e}_{n-1}^{n-1})^T). \end{aligned}$$

Finally, we have

$$\begin{aligned} \mathbf{S} &= \left(\frac{1}{h} + \frac{\nu h}{3}\right) \mathbf{D} \\ &\quad - \left(\frac{-1}{h} + \frac{\nu h}{6}\right)^2 (\mathbf{E}^+ \otimes (\mathbf{e}_1^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) + |\mathbf{E}^-| \otimes (\mathbf{e}_{n-1}^{n-1})^T)((\mathbf{E}^+)^T \otimes \mathbf{e}_1^{n-1} + |\mathbf{E}^-|^T \otimes \mathbf{e}_{n-1}^{n-1}). \end{aligned}$$

By expanding the products, we have

$$\begin{aligned} \mathbf{S} &= \left(\frac{1}{h} + \frac{\nu h}{3}\right) \mathbf{D} - \left(\frac{-1}{h} + \frac{\nu h}{6}\right)^2 \\ &\quad \cdot \left[(\mathbf{E}^+(\mathbf{E}^+)^T) \otimes (\mathbf{e}_1^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_1^{n-1} \right. \\ &\quad + (\mathbf{E}^+|\mathbf{E}^-|^T) \otimes (\mathbf{e}_1^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_{n-1}^{n-1} \\ &\quad + (|\mathbf{E}^-|(\mathbf{E}^+)^T) \otimes (\mathbf{e}_{n-1}^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_1^{n-1} \\ &\quad \left. + (|\mathbf{E}^-||\mathbf{E}^-|^T) \otimes (\mathbf{e}_{n-1}^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_{n-1}^{n-1} \right]. \end{aligned}$$

Because of the symmetry of \mathbf{U} and from (A.1), (A.2) and (A.3), we have

$$\begin{aligned} (\mathbf{e}_{n-1}^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_1^{n-1} &= (\mathbf{e}_1^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_{n-1}^{n-1}, \\ (\mathbf{e}_{n-1}^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_{n-1}^{n-1} &= (\mathbf{e}_1^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_1^{n-1}. \end{aligned}$$

Thus, from Lemma 4.1 and the property $|\mathbf{E}^-| = -\mathbf{E}^-$, it follows that

$$\begin{aligned} \mathbf{S} &= \left(\frac{1}{h} + \frac{\nu h}{3}\right) \mathbf{D} - \left(\frac{-1}{h} + \frac{\nu h}{6}\right)^2 \left[(\mathbf{E}^+\mathbf{E}^{+T} + \mathbf{E}^-\mathbf{E}^{-T}) \otimes (\mathbf{e}_1^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_1^{n-1} \right. \\ &\quad \left. - (\mathbf{E}^+\mathbf{E}^{-T} + \mathbf{E}^-\mathbf{E}^{+T}) \otimes (\mathbf{e}_{n-1}^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_1^{n-1} \right] \\ &= \left(\frac{1}{h} + \frac{\nu h}{3}\right) \mathbf{D} - \left(\frac{-1}{h} + \frac{\nu h}{6}\right)^2 \left[\mathbf{D} \otimes (\mathbf{e}_1^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_1^{n-1} + \mathbf{A}\mathbf{d} \otimes (\mathbf{e}_{n-1}^{n-1})^T (\mathbf{U}\Theta^{-1}\mathbf{U}) \mathbf{e}_1^{n-1} \right]. \end{aligned}$$

Moreover, we have

$$\begin{aligned} \mathbf{S} &= \left[\left(\frac{1}{h} + \frac{\nu h}{3} \right) - \left(\frac{-1}{h} + \frac{\nu h}{6} \right)^2 (\mathbf{e}_1^{n-1})^\top (\mathbf{U}\Theta^{-1}\mathbf{U})\mathbf{e}_1^{n-1} \right] \mathbf{D} \\ &\quad - \left[\left(\frac{-1}{h} + \frac{\nu h}{6} \right)^2 (\mathbf{e}_{n-1}^{n-1})^\top (\mathbf{U}\Theta^{-1}\mathbf{U})\mathbf{e}_{n-1}^{n-1} \right] \mathbf{Ad}. \end{aligned}$$

In the general case $\nu > 0$, it follows that the nonzero pattern of \mathbf{S} will always coincide with that of \mathbf{L}_T .

LEMMA A.1

$$\left(\frac{1}{h} + \frac{\nu h}{3} \right) - \left(\frac{-1}{h} + \frac{\nu h}{6} \right)^2 (\mathbf{e}_1^{n-1})^\top (\mathbf{U}\Theta^{-1}\mathbf{U})\mathbf{e}_1^{n-1} \lesssim 1 + \frac{(n-1)\nu}{3\pi^2 + \nu} + \mathcal{O}(h). \quad (\text{A.4})$$

If $\nu < h^{-1}$, then we have the estimate

$$\left(\frac{-1}{h} + \frac{\nu h}{6} \right)^2 (\mathbf{e}_{n-1}^{n-1})^\top (\mathbf{U}\Theta^{-1}\mathbf{U})\mathbf{e}_{n-1}^{n-1} \approx \frac{-1}{1 + \frac{\nu}{3\pi^2}} + \mathcal{O}(h). \quad (\text{A.5})$$

Proof. We have

$$\begin{aligned} & - \left(\frac{-1}{h} + \frac{\nu h}{6} \right)^2 (\mathbf{e}_1^{n-1})^\top (\mathbf{U}\Theta^{-1}\mathbf{U})\mathbf{e}_1^{n-1} \\ &= -\frac{1}{h^2} \left(1 - \frac{\nu h^2}{6} \right)^2 \sum_{j=1}^{n-1} \frac{2h^2 \sin^2(j\pi h)}{\left(1 + \frac{\nu h^2}{6} \right) \lambda_j + \frac{\nu h^2}{3}} = - \left(1 - \frac{\nu h^2}{6} \right)^2 \sum_{j=1}^{n-1} \frac{\left[2 \cos^2 \left(j \frac{\pi h}{2} \right) \right] \lambda_j}{\left(1 + \frac{\nu h^2}{6} \right) \lambda_j + \frac{\nu h^2}{3}} \\ &= - \left(1 - \frac{\nu h^2}{6} \right)^2 \sum_{j=1}^{n-1} \frac{\left[2 \cos^2 \left(j \frac{\pi h}{2} \right) \right]}{\left(1 + \frac{\nu h^2}{6} \right) + \frac{\nu h^2}{3\lambda_j}} \leq - \left(1 - \frac{\nu h^2}{6} \right)^2 \sum_{j=1}^{n-1} \frac{\left[2 \cos^2 \left(j \frac{\pi h}{2} \right) \right]}{\left(1 + \frac{\nu h^2}{6} \right) + \frac{\nu h^2}{3\lambda_1}} \\ &\approx - \left(1 - \frac{\nu h^2}{6} \right)^2 \sum_{j=1}^{n-1} \frac{\left[2 \cos^2 \left(j \frac{\pi h}{2} \right) \right]}{\left(1 + \frac{\nu h^2}{6} \right) + \frac{\nu}{3\pi^2}} = - \frac{\left(1 - \frac{\nu h^2}{6} \right)^2}{1 + \frac{\nu h^2}{6} + \frac{\nu}{3\pi^2}} (n-1). \end{aligned}$$

Here, we have used the identity

$$\sum_{j=1}^{n-1} 2 \cos^2 \left(\frac{j\theta}{2} \right) = \sum_{j=1}^{n-1} [\cos(j\theta) + 1] = \sum_{j=1}^{n-1} \cos(j\theta) + n - 1$$

(with $\theta = \pi h = \pi/n$) and the Dirichlet kernel formula

$$1 - 2 \sum_{j=1}^{n-1} \cos(j\theta) = \frac{\sin\left(\left(n - \frac{1}{2}\right)\theta\right)}{\sin\left(\frac{\theta}{2}\right)},$$

which together yield

$$\sum_{j=1}^{n-1} 2 \cos^2\left(\frac{j\theta}{2}\right) = n - 1.$$

Therefore, we have (A.4).

From (A.1) and expanding $\lambda^{-1} = \csc^2(j\pi h/2)/4$ to the first-order ($j\pi h/2 < \pi$), we obtain that

$$\begin{aligned} \left(\frac{-1}{h} + \frac{\nu h}{6}\right)^2 (\mathbf{e}_{n-1}^{n-1})^\top (\mathbf{U}\Theta^{-1}\mathbf{U})\mathbf{e}_1^{n-1} &= \left(1 - \frac{\nu h^2}{6}\right)^2 \sum_{j=1}^{n-1} \frac{(-1)^j 2 \cos^2(j\frac{\pi}{2}h)}{\left(1 + \frac{\nu h^2}{6}\right) + \frac{\nu h^2}{3\lambda_j}} \\ &\approx \left(1 - \frac{\nu h^2}{6}\right)^2 \sum_{j=1}^{n-1} \frac{(-1)^j 2 \cos^2(j\frac{\pi}{2}h)}{\left(1 + \frac{\nu h^2}{6}\right) + \frac{\nu}{3} \left[\frac{1}{j\pi} + \frac{j\pi h^2}{48}\right]^2} \\ &\approx \left(1 - \frac{\nu h^2}{6}\right)^2 \sum_{j=1}^{n-1} \frac{(-1)^j 2 \cos^2(j\frac{\pi}{2}h)}{1 + \frac{\nu}{3\pi^2} \left[1 + \mathcal{O}\left(\frac{\pi^2 h}{48}\right)\right]^2} \approx \frac{1}{1 + \frac{\nu}{3\pi^2}} \sum_{j=1}^{n-1} (-1)^j 2 \cos^2\left(j\frac{\pi}{2}h\right) + \mathcal{O}(h). \end{aligned}$$

Finally, we obtain (A.5) by observing that $\sum_{j=1}^{n-1} (-1)^j 2 \cos^2(j\frac{\pi}{2}h) = -1$. □

REMARK A.2 We observe that the estimate (A.5) is only a crude approximation of the real value of the nonzero entries of \mathbf{S} (corresponding to the entries of \mathbf{Ad}) when $\nu \geq h^{-1}$. In this case, we can only say that the off-diagonal entries in the Schur complement will go to zero at least as fast as ν^{-1} for $\nu \rightarrow \infty$. However, numerically we observe that they are much closer to zero than the estimated value in (A.5).