

# **On Spline Galerkin Methods for Singular Integral Equations with Piecewise Continuous Coefficients**

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**Summary.** This article analizes the convergence of the Galerkin method with polynomial splines on arbitrary meshes for systems of singular integral equations with piecewise continuous coefficients in  $L^2$  on closed or open Ljapunov curves. It is proved that this method converges if and, for scalar equations and equidistant partitions, only if the integral operator is strongly elliptic (in some generalized sense). Using the complete asymptotics of the solution, we provide error estimates for equidistant and for special nonuniform partitions.

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## **Introduction**

Singular integral equations arising from numerous boundary value problems in aerodynamics, elasticity and thermoelasticity, electromagnetics, vibration theory and many other fields of engineering have the form

$$
c(t) x(t) + \frac{d(t)}{\pi i} \int\limits_{\Gamma} \frac{x(\tau)}{\tau - t} d\tau + \int\limits_{\Gamma} K(t, \tau) x(\tau) d\tau = y(t), \quad t \in \Gamma, \tag{0.1}
$$

see e.g. Anderssen [2], Kupradse [11], Michlin, Prößdorf [12], Muskhelishvili [13] and Parton, Perlin [15]. In this equation  $\Gamma$  is a closed or open plane curve, c, d, K and y are given functions (in general matrix - or vector-valued and discontinuous),  $x$  is the unknown function and the first integral is defined as a Cauchy principal value.

In many practical computations with Eq. (0.1) one uses spline approximations for the unknown function  $x$  on  $\Gamma$ . The two most popular discretization schemes are the Galerkin and collocation methods. If  $\Gamma$  is a closed curve and  $c$ and d are continuous functions, then convergence and error estimates for Galerkin and collocation methods using smooth polynomial splines follow from recent studies by Arnold, Saranen, Stephan and Wendland in [4, 23, 27], by Nedelec in [14] and by Elschner, Schmidt and the authors in [8, 16, 17, 19, 21, 22, 24, 25, 26] (see also the surveys just given by Elschner, Prößdorf  $\lceil 18 \rceil$ and by Wendland [29]). Convergence results on the spline collocation method in the space  $L^2$  for Eq. (0.1) with piecewise continuous coefficients c and d have most likely been first obtained by the authors in [19]. They were then generalized by Schmidt in [25]. For  $\Gamma$  a finite interval, convergence and error analysis for Galerkin's method with (smooth as well as weighted) splines has recently been developed by Elschner [9] for the case of continuous  $c$  and  $d$  such that  $\text{Re}\left\{c(t)+d(t)\right\} > 0$  on *F*. The special case where *K* is real-valued,  $c \equiv 1$  and *d* is purely imaginary has earlier been considered by Thomas [28].

It is the aim of this paper to give a convergence and error analysis for the spline Galerkin method in the space  $L^2$  if c and d are piecewise continuous matrix functions and Eq.  $(0.1)$  is strongly elliptic (in the sense of the subsequent definition). For sake of brevity, all results of this paper are proved for the case  $K=0$ , since they easily extend to the complete equation (0.1). (see, e.g. [10]).

We now introduce several notational conventions.  $\mathbb{C}^m(m \in \mathbb{N})$  denotes mdimensional complex Euclidean space with the usual scalar product

$$
[x, y] := \sum_{i=1}^{m} x_i \bar{y}_i, \quad \forall x, y \in \mathbb{C}^m,
$$
  

$$
x = (x_1, \dots, x_m)^T, \quad y = (y_1, \dots, y_m)^T.
$$

 $\mathbb{C}^{m \times m}$  stands for the set of all complex-valued  $m \times m$  matrices.

Let  $\Gamma$  be a closed or open oriented plane Liapunov curve having a regular parametrization

$$
\Gamma := \{t = \gamma(s) : s \in [0, 1]\}, \quad \gamma : [0, 1] \to \mathbb{C};
$$

for *F* closed, we assume  $\gamma(0) = \gamma(1)$ . By  $L^2(\Gamma, \mathbb{C}^m)$  we denote the Hilbert space of all square Lebesgue-integrable  $\mathbb{C}^m$ -valued functions on  $\Gamma$  with scalar product

$$
(f, g) := \int\limits_{\Gamma} [f(t), g(t)] \, |dt|, \quad \forall \, f, g \in L^2(\Gamma, \mathbb{C}^m).
$$

The symbol  $PC(\Gamma, \mathbb{C}^{m \times m})$  designates the space of all  $\mathbb{C}^{m \times m}$ -valued functions a on  $\Gamma$  which are piecewise continuous in the following sense: for each  $t \in \Gamma$  the finite limits  $a(t+0)$ : = lim  $a(\tau)$  (with respect to the orientation of  $\Gamma$ ) exist and  $\tau \rightarrow t \pm 0$ 

a is discontinuous at most at a finite number of points  $t \in \Gamma$ .  $C(\Gamma, \mathbb{C}^{m \times m})$  is the subspace of all continuous  $\mathbb{C}^{m \times m}$ -valued functions on  $\Gamma$ .

In  $L^2(\Gamma, \mathbb{C}^m)$  we consider singular integral operators of the form

$$
A = aP_r + bQ_r \tag{0.2}
$$

with coefficients *a*,  $b \in PC(\Gamma, \mathbb{C}^{m \times m})$ . Here  $P_r$  and  $Q_r$  denote the operators

$$
P_T
$$
 := 1/2(I + S\_T),  $Q_T$  := 1/2(I - S\_T),

I the indentity operator and  $S_r$  the Cauchy singular operator

$$
(S_T x)(t) = 1/\pi i \int\limits_{\Gamma} \frac{x(\tau)}{\tau - t} d\tau (t \in \Gamma).
$$

That  $A \in \mathcal{L}(L^2(\Gamma, \mathbb{C}^m))$  (see [12]) is well established and the operator defined by the left-hand side of Eq. (0.1) with  $K \equiv 0$  has the form (0.2), where  $a = c + d$  and  $b=c-d$ .

We call operator  $A=aP_r+bQ_r$  strongly elliptic<sup>1</sup>, if there exist a compact operator  $T \in \mathscr{L}(L^2(\Gamma, \mathbb{C}^m))$  and an invertible function  $\Theta \in PC(\Gamma, \mathbb{C}^{m \times m})$ , discontinuous at most at the points of discontinuity of a and b, such that  $A = \Theta(A_0)$  $+ T$ ), where  $A_0$  has a positive real part - i.e.

$$
\operatorname{Re}(A_0 f, f) \ge \varepsilon(f, f), \quad \forall f \in L^2(\Gamma, \mathbb{C}^m)
$$

with  $\varepsilon$ (const) > 0.

In [20] we have shown that A is strongly elliptic if and only if

$$
\forall t \in \Gamma: \exists b(t \pm 0)^{-1} \quad \text{and} \quad \exists C \in \mathbb{C}^{m \times m}: \text{Re } C > 0, \text{Re } C(b^{-1}a)(t \pm 0) > 0. \tag{B}
$$

(If  $\Gamma = (\alpha, \beta)$  is an open curve, then condition (B) need only be satisfied from the right or from the left, according as we consider  $\alpha$  or  $\beta$  respectively.) In Lemma 1.2 further conditions equivalent to (B) for special cases are formulated.

Now let  $\{t_0=\gamma(s_0), t_1=\gamma(s_1),...,t_N=\gamma(s_N)\}\subseteq \Gamma(0=s_0 be a$ given set of points such that a and b are continuous on  $\Gamma \setminus \{t_0, ..., t_N\}$ . A sequence  $\{\Lambda_k\}_{k\in\mathbb{N}}$   $(\Lambda_k=\{\sigma_0^k,\ldots,\sigma_{n_k}^k\}, Q=\sigma_0^k<\sigma_1^k<\ldots\sigma_{n_k}^k=1)$  of partitions of [0, 1] is called admissible if  $\{s_0, ..., s_N\} \subseteq A_k(k \in \mathbb{N})$  and if

 $h_{A_n}$ : = max { $\sigma_{i+1}^k-\sigma_i^k$  |  $i=0, 1, ..., n_k-1$ }  $\rightarrow 0$  (k  $\rightarrow \infty$ ).

Let  $PS_d(\Lambda_k, \mathbb{C})$  (deN) denote the space of all  $\varphi \in PC(\Gamma, \mathbb{C})$  such that  $\varphi \circ \gamma$  is (d  $-1$ ) times continuously differentiable on  $[0, 1] \setminus \{s_0, ..., s_N\}$  and the restriction of  $\varphi \circ \gamma$  to  $[\sigma_i^k, \sigma_{i+1}^k]$  is a polynomial of degree not exceeding  $d$  - i.e. the functions  $\varphi \in PS_{d}(A_{k}, \mathbb{C})$  are splines with maximal smoothness  $\sigma_i^k \in \Lambda_k \setminus \{s_0, ..., s_N\}$  and with no smoothness condition at  $s_i(i=0, ..., N)$ .<br> $PS_d(\Lambda_k, \mathbb{C}^m)$  stands for the space of vector-valued functions f  $PS_d(\hat{A}_k, \mathbb{C}^m)$  stands for the space of vector-valued functions f  $=(f_1, ..., f_m) \in PC(\Gamma, \mathbb{C}^m)$  having components  $f_j \in PS_d(\Lambda_k, \mathbb{C})$ . Obviously,  $PS_0(\Lambda_k, \mathbb{C}^m)$  is the space of step functions.

The Galerkin method for the equation  $Ax = y$  may now be formulated as follows: Find  $x_{4k} \in PS_d(\Lambda_k, \mathbb{C}^m)$  satisfying

$$
(Ax_{\Lambda_k}, \varphi) = (y, \varphi) \qquad (\forall \varphi \in PS_d(A_k, \mathbb{C}^m)). \tag{0.3}
$$

We say that the Galerkin method with respect to the admissible sequence  ${A_k}_{k \in \mathbb{N}}$  is convergent for the operator A, if (0.3) is uniquely solvable for sufficiently large k and if  $x_{4k}$  converges to  $x = A^{-1}y$  in  $L^2(\Gamma, \mathbb{C}^m)$  for any  $y \in L^2(\Gamma, \mathbb{C}^m)$ .

We proceed in this paper as follows. In Sect. 1 we examine the Galerkin method (0.3) with respect to any admissible sequence  $\{\Lambda_k\}_{k\in\mathbb{N}}$  for the operator A defined by (0.2). By means of the discrete commutator property, given in Lemma 1.1, it is proved that this method is convergent in  $L^2(\Gamma, \mathbb{C}^m)$  if A is strongly elliptic and invertible. In the special case  $\Gamma = (0, 1)$ ,  $m = 1$  and  $a, b \in C$ 

<sup>&</sup>lt;sup>1</sup> See Stephan and Wendland [27], who introduced this concept in case of smooth functions *O,a,b* 

[0, 1], Theorem 1.1 yields a result of Elschner ([9], Theorem 3.1). In Sect. 2 we show that strong ellipticity of  $A$  is necessary for the Galerkin method to be convergent, provided the partitions are equidistant and  $m = 1$ . For continuous coefficients a and b and a closed curve  $\Gamma$ , this fact was established in [24] by means of a different technique. As a corollary to our proof, which is based on certain localization techniques and on the method of associated operators (see [19]), we obtain that the strong ellipticity of  $A$  is sufficient for the stability of the collocation method with piecewise linear splines. Section 3 deals with estimating the error of the Galerkin method  $(0.3)$  by means of utilizing the complete asymptotics of the solutions of (0.1) which has been established in [9]. We show for instance, that an asymptotic order of convergence of  $O(n_k^{-d-1})$  can be achieved on special nonuniform partitions. This fact was obtained in [28] and [9] by using weighted continuous splines on nonuniform meshes for continuous coefficients a and b and  $\Gamma = (0, 1)$ .

All results of Sects. 1 and 2 are valid for singular integral operators of the form  $P_r a + Q_r b$ ; one need only to take the adjoints of the operators  $a^* P_r$  $+b^*Q_r$ .

# **1. The Convergence of the Galerkin Method for Strongly Elliptic Operators**

Let  $P_{A_k}$  denote the orthogonal projection of  $L^2(\Gamma, \mathbb{C}^m)$  onto  $PS_d(A_k, \mathbb{C}^m)$ . Then (0.3) is equivalent to

$$
P_{\underline{A}_k} A P_{\underline{A}_k} x_{\underline{A}_k} = P_{\underline{A}_k} y.
$$

**Lemma 1.1.** *If*  $\{\Delta_k\}$  *is an admissible sequence of partitions and*  $f \in PC(\Gamma, \mathbb{C}^{m \times m})$ *is continuous on*  $\Gamma \setminus \{t_0, t_1, \ldots, t_N\}$ , then

$$
||(I - P_{A_k})f P_{A_k}|| \to 0 (k \to \infty)
$$
  

$$
||P_{A_k} f (I - P_{A_k})|| \to 0 (k \to \infty).
$$

*Proof.* We only need to prove  $\|(I-P_{A_k})fP_{A_k}\| \to 0$ , since the other assertion follows by taking the operator adjoints.

$$
f = (f_{ij})_{i,j=1}^m \in PC(\Gamma, \mathbb{C}^{m \times m})
$$

implies

$$
(I - P_{A_k}) f P_{A_k} = ((I - P_{A_k}) f_{ij} P_{A_k})_{i,j=1}^m,
$$

where  $P_{4k}$  in the entries of the last matrix are the orthogonal projections of  $L^2(\Gamma, \mathbb{C})$  onto  $PS_d(\Lambda_k, \mathbb{C})$ . Therefore, we have to prove  $\|(I-P_{\Lambda_k})fP_{\Lambda_k}\| \to 0$  for the case  $m = 1$  only.

Now we introduce a new scalar product in  $L^2(\Gamma, \mathbb{C})$  by

$$
(f,g)_1:=\int\limits_0^1(f\circ\gamma)(s)(\overline{g}\circ\gamma)(s)\,ds.
$$

Clearly, the norms  $||f|| = (f, f)^{1/2}$  and  $||f||_1 = (f, f)^{1/2}$  are equivalent, and thus it suffices to show  $\|(I-P_{A_k})fP_{A_k}\|_1\to 0$ , where  $\|\cdot\|_1$  denotes the operator norm corresponding to the norm  $\|\cdot\|_1$  in  $L^2(\Gamma, \mathbb{C})$ . Let  $\tilde{P}_{A_k}$  be the orthogonal projection of  $L^2(\Gamma, \mathbb{C})$  onto  $PS_d(A_k, \mathbb{C})$  corresponding to  $(\cdot, \cdot)_1$ . Since  $\sup \|P_{4k}\|_1 < \infty$ , we obtain (see [1]) k

$$
\begin{aligned} \| (I - P_{A_k}) f P_{A_k} g \|_1 &= \| f P_{A_k} g - P_{A_k} (f P_{A_k} g) \|_1 \\ &\leq C \, \| (f P_{A_k} g) - \tilde{P}_{A_k} (f P_{A_k} g) \|_1 \\ &= C \, \| (I - \tilde{P}_{A_k}) f \tilde{P}_{A_k} P_{A_k} g \|_1, \end{aligned}
$$

where  $C$  is a positive constant. Consequently, we only have to show  $||(I-\tilde{P}_{4k})f\tilde{P}_{4k}||_1 \rightarrow 0$ . Hence, we may assume  $\Gamma = [0, 1].$ 

Let  $f^{i}(i=0, ..., N-1)$  be the function vanishing on  $[0, 1] \setminus [s_i, s_{i+1}]$  and equal to f on  $[s_i, s_{i+1}]$ , and  $P_{4k}^i$  be the orthogonal projection onto the subspace of all functions  $f \in PS_d(A_k, \mathbb{C})$  vanishing on  $[0, 1] \setminus [s_i, s_{i+1}]$ . Thus

$$
(I - P_{A_k}) f P_{A_k} = \sum_{i=0}^{N-1} (I - P_{A_k}^i) f^i P_{A_k}^i
$$

and we only need to demonstrate  $||(I-P_{A_k}^i)f^iP_{A_k}^i|| \rightarrow 0$   $(i=0, ..., N-1)$ . We may therefore assume:  $f \in C([0,1],\mathbb{C})$  and  $P_{4k}$  is the orthogonal projection of  $L^2([0, 1], \mathbb{C})$  onto  $S_d(\Lambda_k, \mathbb{C})$ .  $S_d(\Lambda_k, \mathbb{C})$  is the subspace of all  $f \in L^2([0, 1], \mathbb{C})$ , which are  $(d-1)$  times continuously differentiable and coincide with a polynomial of degree not exceeding d on every interval  $[\sigma_i^k, \sigma_{i+1}^k]$  ( $i=0, ..., n_k-1$ ).

We set

$$
\sigma_{-r}^k := -1/2r(\sigma_1^k - \sigma_0^k) \qquad (r = 1, 2, ..., d)
$$

and

$$
\sigma_{n_k+r}^k := 1 + 1/2 r (\sigma_{n_k}^k - \sigma_{n_k-1}^k) \qquad (r = 1, 2, ..., d).
$$

A base  $\{\psi_i\}_{i=-d}^{n_k}$  of  $S_d(\Lambda_k, \mathbb{C})$  can now be given as follows (see [6]):

$$
g(s, t) := (s - t)^{d}_{+}
$$
  

$$
\psi_{i}(t) := \sqrt{(d+1)} \sqrt{(\sigma_{i+d+1}^{k} - \sigma_{i}^{k})} g(\sigma_{i}^{k}, \dots, \sigma_{i+d+1}^{k}; t) \qquad (t \in [0, 1])
$$

wherein  $g(\sigma_i^k, \ldots, \sigma_{i+d+1}^k; t)$  represents the divided difference of  $g(s, t)$  at  $\sigma_i^k, \sigma_{i+1}^k, \ldots, \sigma_{i+d+1}^k$ . It is well known that  $\psi_i \geq 0$  and  $\psi_i$  has support  $[\sigma_i^k, \sigma_{i+d+1}^k]$  (see [6]). Moreover, there exists a positive constant D such that

$$
D^{-1}\left\|\sum_{i=-d}^{n_{k}-1}\xi_{i}\psi_{i}\right\| \leq \left(\sum_{i=-d}^{n_{k}-1}|\xi_{i}|^{2}\right)^{1/2} \leq D\left\|\sum_{i=-d}^{n_{k}-1}\xi_{i}\psi_{i}\right\| \tag{1.1}
$$

(see [7]). We obtain

$$
\left\|f\sum_{i=-d}^{n_{k}-1}\xi_{i}\psi_{i}-P_{A_{k}}f\sum_{i=-d}^{n_{k}-1}\xi_{i}\psi_{i}\right\| \leq \left\|f\sum_{i=-d}^{n_{k}-1}\xi_{i}\psi_{i}-\sum_{i=-d}^{n_{k}-1}f(\sigma_{i}^{k})\xi_{i}\psi_{i}\right\|, \qquad (1.2)
$$

where  $f(\sigma_i^k) = f(0)$  ( $i = -d, -d + 1, ..., -1$ ). For  $t \in [\sigma_i^k, \sigma_{i+1}^k]$ ,

$$
\left| f(t) \sum_{i=-d}^{n_{k}-1} \xi_{i} \psi_{i}(t) - \sum_{i=-d}^{n_{k}-1} f(\sigma_{i}^{k}) \xi_{i} \psi_{i}(t) \right| = \left| \sum_{i=j-d}^{j+1} (f(t) - f(\sigma_{i}^{k})) \xi_{i} \psi_{i}(t) \right|
$$
  

$$
\leq \omega(f, dh_{A_{k}}) \sum_{i=j-d}^{j+1} |\xi_{i}| |\psi_{i}(t)|
$$
  

$$
\leq \omega(f, dh_{A_{k}}) \left| \sum_{i=-d}^{n_{k}-1} |\xi_{i}| |\psi_{i}(t) \right| \tag{1.3}
$$

with  $\omega(f, \delta) = \sup \{|f(t_1) - f(t_2)||t_1, t_2 \in [0, 1], |t_1 - t_2| < \delta\}$  (modulus of continuity).  $(1.1)$  and  $(1.3)$  yield

$$
\left\|f\sum_{i=-d}^{n_{k}-1}\xi_{i}\psi_{i}-\sum_{i=-d}^{n_{k}-1}f(\sigma_{i}^{k})\xi_{i}\psi_{i}\right\|\leq\omega(f,dh_{A_{k}})\left\|\sum_{i=-d}^{n_{k}-1}|\xi_{i}|\psi_{i}\right\|
$$
  

$$
\leq\omega(f,dh_{A_{k}})D^{2}\left\|\sum_{i=-d}^{n_{k}-1}\xi_{i}\psi_{i}\right\|
$$

and (1.2) implies

$$
\left\|(I-P_{A_k})f\sum_{i=-d}^{n_k-1}\xi_i\psi_i\right\|\leq D^2\omega(f,dh_{A_k})\left\|\sum_{i=-d}^{n_k-1}\xi_i\psi_i\right\|.
$$

Since  $h_{\Delta_k} \rightarrow 0$ , the proof is completed.

*Remark.* Lemma 1.1 was independently given in [16, 17] and in [3] for the case of smooth functions  $f$ . In these papers however, the authors considered arbitrary Sobolev norms instead of  $L^2$ -norms and consequently the partitions were assumed to be regular. (See also [25], where a detailed proof of the aforementioned result is given for the case of Sobolev norms.) The idea of the proof presented above is based on the method of derivation of Lemma 4.1 in [21].

**Corollary 1.1.** *If*  $\{A_k\}$  *is admissible, then*  $P_{A_k}$  *converges strongly to*  $I$  - *i.e.*  $P_{A_k}$   $\rightarrow$  *I. Proof.* Consider  $g = (g_1, g_2, \dots, g_m) \in C(\Gamma, \mathbb{C}^m)$  and set

$$
f:=(g_i\,\delta_{ij})_{i,j=1}^m\in C(\Gamma,\,\mathbb{C}^{m\times m}),\,h=(1,\,\ldots,\,1)\in C(\Gamma,\,\mathbb{C}^m).
$$

Lemma 1.1 implies  $(I-P_{A_k})g=(I-P_{A_k})fP_{A_k}h\rightarrow 0$  and further  $P_{A_k}\rightarrow I$ , since  $\sup \|P_{A_k}\| < \infty.$ 

It is now rather simple to prove

**Theorem 1.1.** If  $\{A_k\}$  is an arbitrary admissible sequence of partitions and A  $=aP_r+bQ_r\in\mathscr{L}(L^2(\Gamma,\mathbb{C}^m))$  is invertible and strongly elliptic, then the Galerkin *method with respect to the sequence*  $\{\Lambda_k\}$  *is convergent for A.* 

*Proof.* Since  $P_{A_k} \rightarrow I$ , the Galerkin method is convergent if and only if the approximate operators  $P_{4k}AP_{4k}$  are stable - i.e.

$$
||P_{A_k}AP_{A_k}\varphi|| \geq C ||\varphi|| \qquad (\forall \varphi \in PS_d(A_k, \mathbb{C}^m), \forall k \geq k_0),
$$

where C is a positive constant and  $k_0 \in \mathbb{N}$  (see, e.g. [10]). The strongly elliptic operator A admits the representation  $A = \Theta A_0 + T$ , where  $\Theta \in PC(\Gamma, \mathbb{C}^{m \times m})$  is

continuous on  $\Gamma \setminus \{t_0, \ldots, t_N\}$ , Re  $A_0 > 0$  and T is compact. Lemma 1.1 yields

$$
P_{A_k} \Theta^{-1} P_{A_k} P_{A_k} \Theta P_{A_k} = -P_{A_k} \Theta^{-1} (I - P_{A_k}) \Theta P_{A_k} + P_{A_k},
$$
  

$$
||P_{A_k} \Theta^{-1} (I - P_{A_k}) \Theta P_{A_k}|| \to 0.
$$

Thus  $P_{A_k} \Theta P_{A_k}$  are stable. Furthermore  $\text{Re } A_0 > 0$  implies that the operators  $P_{4\nu}A_0 P_{4\nu}$  are also stable. The stability of the operators

$$
P_{A_k} A P_{A_k} = (P_{A_k} \Theta P_{A_k}) (P_{A_k} A_0 P_{A_k}) + P_{A_k} \Theta (I - P_{A_k}) A_0 P_{A_k} + P_{A_k} T P_{A_k}
$$

follows from Lemma 1.1 and the fact that T is compact (see, e.g.  $[10]$ ).

In the following two cases, condition (B) may be easily verified.

#### **Lemma 1.2. [22, 20],**

a) If  $m=1$ , then (B) holds at  $\tau \in \Gamma$  if and only if  $b(\tau +0) \neq 0$  and

$$
\forall \mu \in [0, 1]: \{\mu(b^{-1}a)(\tau + 0) + (1 - \mu)(b^{-1}a)(\tau - 0)\} \notin (-\infty, 0]. \tag{1.4}
$$

/f b) If a and b are continuous at  $\tau \in \Gamma$ , then condition (B) holds at  $\tau$  if and only  $\forall \mu \in [0, 1]: \det(\mu a(\tau) + (1 - \mu) b(\tau)) \neq 0.$ 

# **2. A Necessary Condition for the Convergence of the Galerkin method. A Collocation Method**

In this section we assume all partitions to be equidistant - i.e.  $\sigma_{i+1}^k - \sigma_i^k = \sigma_{j+1}^k$  $-\sigma_i^k$  ( $\forall i, j=0, ..., n_k-1$ ). Actually, there exist admissible sequences of equidistant partitions In particular, we can choose for a given set  $\{t_0, ..., t_N\} \subseteq \Gamma$  a regular parametrization  $\gamma$  such that  $\gamma(i/N)=t_i$  and thus get an admissible sequence of equidistant partitions

$$
\Delta_k := \{ \sigma_0^k := 0, \, \sigma_1^k := 1/kN, \, \sigma_2^k := 2/kN, \, \dots, \, \sigma_{kN}^k = 1 \}.
$$

Now let us consider the Galerkin method (0.3) for the singular operator A.

**Theorem 2.1.** Let  $m = 1$  and  $a, b \in PC(\Gamma, \mathbb{C})$ . If  $\{A_k\}$  is an admissible sequence of *equidistant partitions and the Galerkin method with respect to this sequence is convergent for*  $A = aP_r + bQ_r \in \mathcal{L}(L^2(\Gamma, \mathbb{C}))$ *, then A is strongly elliptic.* 

Before we prove Theorem 2.1, we introduce the same notations as in [19] for the collocation method, and we construct to this end a base of  $PS_d(\Lambda_k, \mathbb{C}^m)$ (*m* $\geq$ 1). On every interval  $(s_i, s_{i+1})$  (*i*=0, ..., *N*-1) the elements of this base differ from the functions  $\psi_i$  of Sect. 1 by at most a constant factor, however, because our partitions are equidistant, we can use the representation of [5], Sect. 4.2.

Let  $\sigma_i = \sigma_i^k = i/n$  (*i* = 0, ..., *n*),  $n = n_k$  and  $t_i = \gamma(s_i)$ ,  $s_i = k_i/n$  (*i* = 0, ..., *N*). We denote the characteristic functions of  $[s_i, s_{i+1}]$  and  $[-1, 0]$  by  $\chi_i$  and  $\Pi$ respectively and further,  $\Pi_{d+1}$  stands for the  $(d+1)$ -fold convolution of  $\Pi$ . Thus the space  $S_d(\Lambda_k, \mathbb{C})$  (see Sect. 1) has the base

$$
\tilde{\psi}_j(s) := \Pi_{d+1}(ns-j) \qquad (j=1,\ldots,n+d).
$$

For  $t = \gamma(s)$  set

$$
\varphi_{1}(t) := (\chi_{0} \tilde{\psi}_{1})(s), \varphi_{2}(t) := (\chi_{0} \tilde{\psi}_{2})(s), \dots, \varphi_{k_{1}+d}(t) := (\chi_{0} \tilde{\psi}_{k_{1}+d})(s),
$$
  
\n
$$
\varphi_{k_{1}+d+1}(t) := (\chi_{1} \tilde{\psi}_{k_{1}+1})(s),
$$
  
\n
$$
\varphi_{k_{1}+d+2}(t) := (\chi_{1} \tilde{\psi}_{k_{1}+2})(s), \dots, \varphi_{k_{2}+2d}(t) := (\chi_{1} \tilde{\psi}_{k_{2}+d})(s),
$$
  
\n
$$
\vdots
$$
  
\n
$$
\varphi_{k_{j}+jd+1}(t) := (\chi_{j} \tilde{\psi}_{k_{j}+1})(s), \varphi_{k_{j}+jd+2}(t) := (\chi_{j} \tilde{\psi}_{k_{j}+2})(s), \dots,
$$
  
\n
$$
\varphi_{k_{j}+1+(j+1)d}(t) := (\chi_{j} \tilde{\psi}_{k_{j+1}+d})(s),
$$
  
\n
$$
\vdots
$$
  
\n
$$
\varphi_{k_{N-1}+(N-1)d+1}(t) := (\chi_{N-1} \tilde{\psi}_{k_{N-1}+1})(s),
$$
  
\n
$$
\varphi_{k_{N-1}+(N-1)d+2}(t) := (\chi_{N-1} \tilde{\psi}_{k_{N-1}+2})(s), \dots,
$$
  
\n
$$
\varphi_{n+Nd}(t) := (\chi_{N-1} \tilde{\psi}_{n+d})(s).
$$
 (2.1)

Every function  $\varphi \in PS_d(\Lambda_k, \mathbb{C}^m)$  admits a representation

$$
\varphi = \sum_{i=1}^{n+Nd} \xi_i \varphi_i, \quad \xi_i \in \mathbb{C}^m
$$

and there exist suitable positive constants  $D_1$ ,  $D_2$  such that out of (1.1) follows

$$
D_1 \left\| \sum_{i=1}^{n+Nd} \zeta_i \, \varphi_i \right\| \le n^{-1/2} \left( \sum_{i=1}^{n+Nd} |\zeta_i|^2 \right)^{1/2} \le D_2 \left\| \sum_{i=1}^{n+Nd} \zeta_i \, \varphi_i \right\| \tag{2.2}
$$

with  $|\xi_i|:=[\xi_i, \xi_i]^{1/2}$ .

We now fix  $\tau = \gamma(\sigma) \in \Gamma$  and derive the subsequent necessary condition at point  $\tau$  (see Lemma 2.2). Using this condition, we shall obtain that condition (B) is necessary if  $m=1$ . Moreover, for  $m\geq 1$  and  $d=0$ , we get a stability condition for a certain collocation method.

Inequality (2.2) allows us to identify the operators of  $\mathscr{L}(PS<sub>d</sub>(A<sub>k</sub>, \mathbb{C}^m))$  with matrices over finite dimensional Euclidean spaces. For this purpose, we introduce associated operators analogously to [19]. Let  $l^2(\mathbb{C}^m)$  be the Hilbert space of the following sequences:

$$
l^{2}(\mathbb{C}^{m}) := \{ \xi := \{ \xi_{s} \}_{s=-\infty}^{\infty} \mid \xi_{s} \in \mathbb{C}^{m}, \|\xi\| < \infty \},
$$
  

$$
\|\xi\| := (\xi, \xi)^{1/2}, \quad (\xi, \eta) := \sum_{s=-\infty}^{\infty} [\xi_{s}, \eta_{s}].
$$

We introduce projections P, Q,  $P_k \in \mathcal{L}(l^2(\mathbb{C}^m))$  by

$$
P\xi := \tilde{\xi}, \quad \tilde{\xi}_s := \begin{cases} \xi_s & \text{if } s \ge 0 \\ 0 & \text{if } s < 0, \end{cases}
$$
\n
$$
Q := I - P.
$$

For an  $i=i(\sigma, k)$  ( $i\in\{0, 1, ..., n\}$ ) and a  $j=j(\sigma, k)$  ( $j\in\{0, ..., N\}$ ), we have  $\sigma\in[i/n,$  $(i+1)/n$  and  $k<sub>i</sub> \leq i < k<sub>i+1</sub>$ . We set

$$
P_k \xi := \xi, \qquad \xi_s := \begin{cases} \xi_s & \text{if } -(i+jd) \le s \le (n-i) + (N-j) \, d-1 \\ 0 & \text{otherwise} \end{cases}
$$

and define

$$
E_k: \text{im } P_k \to PS_d(\Lambda_k, \mathbb{C}^m),
$$
  
\n
$$
E_k \xi := \sqrt{n} \{ \xi_{-(i+jd)} \varphi_1 + \xi_{-(i+jd)+1} \varphi_2 + \dots
$$
  
\n
$$
+ \xi_0 \varphi_{i+jd+1} + \xi_1 \varphi_{i+jd+2} + \dots + \xi_{(n-i)+(N-j)d-1} \varphi_{n+Nd} \}.
$$

(Thus  $P_k$ ,  $E_k$  are chosen in such a manner that  $E_k^{-1} \varphi_{i+jd+1} = {\delta_{i,0}}_{i=-\infty}^{\infty}$ ; however,  $\varphi_{i+jd+1}$  is an element of the base of  $PS_d(A_k, \mathbb{C}^m)$ , the support of which contains  $\tau = \gamma(\sigma)$ .) To every operator  $B_k \in \mathcal{L}(PS_d(A_k, \mathbb{C}^m))$  we associate the operator  $\tilde{B}_k = E_k^{-1} B_k E_k \in \mathcal{L}(\text{im } P_k)$ , which can be identified with  $\tilde{B}_k P_k \in \mathscr{L}(l^2(\mathbb{C}^m)).$ 

If the Galerkin method converges, then the operators

$$
A_k := P_{A_k} A P_{A_k} \in \mathcal{L}(PS_d(\Lambda_k, \mathbb{C}^m))
$$

are stable. By (2.2) the same fact holds for the sequences  $\{\tilde{A}_{k}\}\$  and  $\{\tilde{A}_{k}^{*}\}\$ , where  $\tilde{A}_k^*$  denotes the adjoint operator in  $\mathscr{L}(l^2(\mathbb{C}^m))$ . Hence the strong limit of  $\{\tilde{A}_k\}$ , the adjoint of which is the strong limit of  $\{\tilde{A}_{k}^{*}\}\)$ , is invertible.

Now we examine the limit of  $\{\tilde{A}_k\}$ . The convergence of  $\{\tilde{A}_k^*\}$  is derivable in the same way. We set

$$
c := 1/2(a+b), \quad d := 1/2(a-b), \quad c_k := P_{A_k} c P_{A_k}, \quad d_k := P_{A_k} d P_{A_k}, \quad S_k := P_{A_k} S_{\Gamma} P_{A_k}.
$$

Then, by Lemma 1.1,  $\{\tilde{A}_k\}$  has the same limit as  $\{\tilde{B}_k: = \tilde{c}_k + \tilde{d}_k\tilde{S}_k\}.$ 

**Lemma 2.1.** (i) *For*  $\sigma \in (0, 1)$  *the following statements are valid:* 

a)  $\{\tilde{c}_k\}$  *converges strongly to*  $(c(\tau+0)P + c(\tau-0)Q + T_1)$ *, where*  $T_1 \in \mathcal{L}(l^2(\mathbb{C}^m))$ *is compact.* 

b)  $\{\tilde{S}_k\}$  *converges strongly to*  $(C(\rho) + T_2)$ *, where*  $T_2 \in \mathcal{L}(l^2(\mathbb{C}^m))$  *is compact* and  $C(\rho)$  denotes the operator of convolution generated by the following function *(see* [19]):

$$
\rho(s) := -\frac{\sum_{r \in \mathbb{Z}} \frac{\text{sign}(r+1/2)}{(s+2r\pi)^{2(d+1)}}}{\sum_{r \in \mathbb{Z}} \frac{1}{(s+2r\pi)^{2(d+1)}}} (s \in [0, 2\pi]).
$$

(ii) *For*  $\sigma = 0$  *or*  $\sigma = 1$  *the following statements are valid:* 

c)  $\{\tilde{c}_k\}$  *converges strongly to*  $(c(\tau) P + T_1)$  *or*  $(c(\tau) Q + T_1)$ *, where*  $T_1 \in \mathcal{L}(\text{im } P)$ *and*  $T_1 \in \mathcal{L}(\text{im }Q)$  are compact respectively.

d)  $\{\tilde{S}_k\}$  converges strongly to  $(PC(\rho)P+T_2)$  or  $(QC(\rho)Q+T_2)$ , where  $T_2 \in \mathscr{L}(\text{im }P)$  and  $T_2 \in \mathscr{L}(\text{im }Q)$  are compact respectively.

*Proof.* We prove (i). Assertion (ii) is demonstrated analogously.

Let  $\sigma \in (0, 1)$  and denote by  $w(c, h)$  the modulus of continuity of  $c \circ \gamma$ :

$$
w(c, h) := \sup \{|c(\gamma(x)) - c(\gamma(y))| \mid x, y \in [0, 1], |x - y| \leq h, s_j \notin [x, y](j = 1, ..., N)\}.
$$

Further,  $G_k$  is the following matrix, which will be identified with the corresponding operators in im  $P_k$  or  $l^2(\mathbb{C}^m)$ :

$$
G_k := (n(\varphi_{s+i+jd+1}, \varphi_{r+i+jd+1}))_{r, s=-(i+jd)}^{(n-i)+(N-j)d-1}.
$$

Writing down the matrix of  $c_k$  corresponding to the base  $\{\varphi_r\}$ , we obtain

$$
\tilde{c}_k = G_k^{-1} \left( n(c \varphi_{r+i+jd+1}, \varphi_{s+i+jd+1}) \right)_{s, r=-\,(i+jd)}^{(n-i)+(N-j)d-1}.
$$

An easy computation yields (choose any  $\sigma'_{r+i+id+1}$  esupp  $\varphi_{r+i+id+1}$ )

$$
\tilde{c}_k = G_k^{-1} (\delta_{r,s} c(\sigma'_{r+i+jd+1}))_{s,r=-\,(i+jd)}^{(n-i)+(N-j)d-1} G_k + O(w(c, 1/n)) \tag{2.3}
$$

where  $O(w(c, 1/n))$  denotes a matrix, whose operator norm is majorized by a constant multiple of  $w(c, 1/n)$ . This term thus becomes negligible in our further considerations. Inequality (2.2) leads to

$$
D_2^{-2}\left(\sum |\xi_s|^2\right) \leq n\left(\sum \xi_s \varphi_s, \sum \xi_r \varphi_r\right) \leq D_1^{-2}\left(\sum |\xi_s|^2\right)
$$
  

$$
D_2^{-2} \|\{\xi_s\}\|^2 \leq \left(\{\xi_s\}, G_k\{\xi_s\}\right) \leq D_1^{-2} \|\{\xi_s\}\|^2.
$$

Hence,  $G_k$  is invertible and  $\{G_k\}$ ,  $\{G_k^{-1}\}$  are uniformly bounded.

Now let u,  $v \in \mathbb{Z}$ ,  $d < |u|$ ,  $|v| < \min\{|i-k_s|-d| | s=0, ..., N; k_s+i\}$ ,  $u_1 = u+d$ ,  $v_1$  $v = v + d$  for u,  $v < 0$  and  $i = k_j$ . For u,  $v > 0$  or  $i \neq k_j$ , we set  $u_1 = u$ ,  $v_1 = v$ . If  $i \neq k_j$ or  $u > 0$ , then

$$
\varphi_{u+i+jd+1} = \varphi_{k_j+jd+(i-k_j+u+1)} = \chi_j \tilde{\psi}_{k_j+(i-k_j+u+1)} = \tilde{\psi}_{i+u_1+1}.
$$

If  $i = k_i$  and  $u < 0$ , then

$$
\varphi_{u+i+jd+1} = \varphi_{k_{j-1}+(j-1)d+k_j-k_{j-1}+u+d+1}
$$
  
=  $\chi_{j-1} \tilde{\psi}_{k_{j-1}+(i-k_{j-1}+u+d+1)} = \tilde{\psi}_{i+u_1+1},$ 

therefore,

$$
(G_k)_{u,v} = n(\varphi_{v+i+jd+1}, \varphi_{u+i+jd+1})
$$
  
=  $n \int_0^1 \tilde{\psi}_{i+u_1+1}(s) \tilde{\psi}_{i+v_1+1}(s) |\gamma'(s)| ds$   
=  $\int_{\mathbb{R}} \Pi_{d+1}(s-(i+u_1+1)) \Pi_{d+1}(s-(i+v_1+1)) |\gamma'(s/n)| ds.$ 

Consequently,  $(G_k)_{u,v} = 0$  for  $u \cdot v < 0$  and by holding u and v fixed, we obtain

$$
\lim_{k \to \infty} (G_k)_{u, v} = |\gamma'(\sigma)| \int_{\mathbb{R}} \Pi_{d+1}(s - (u_1 - v_1)) \Pi_{d+1}(s) ds
$$
  
=  $|\gamma'(\sigma)| \int_{\mathbb{R}} e^{-i(u_1 - v_1)s} |F \Pi_{d+1}(s)|^2 ds,$  (2.5)

where  $FH_{d+1}$  denotes the Fourier transform of  $H_{d+1}$ . This means

$$
F\Pi_{d+1} = (\sqrt{2\pi})^d (F\Pi)^{d+1},
$$
  
( $F\Pi$ )( $s$ ) =  $\sqrt{2/\pi} \frac{\sin s/2}{s} e^{-i \cdot s/2},$   
 $|F\Pi_{d+1}(s)|^2 = 1/2\pi \left(\frac{\sin s/2}{s/2}\right)^{2(d+1)}.$ 

Substituting the last equation in (2.5) yields

$$
\lim_{k \to \infty} (G_k)_{u, v} = 1/2\pi \int_0^{2\pi} e^{-i(u_1 - v_1)s} \zeta(s) ds,
$$
  

$$
\zeta(s) := 4^{(d+1)} |\gamma'( \sigma) | (\sin s/2)^{2(d+1)} \sum_{r \in \mathbb{Z}} \frac{1}{(s + 2r\pi)^{2(d+1)}},
$$
  

$$
\lim_{k \to \infty} (G_k)_{u, v} = 0 \quad \text{(for } u \cdot v < 0).
$$

For  $|u|, |v| \leq d$ , we compute  $\lim_{k \to \infty} (G_k)_{u,v}$  in an analogous way. Thereby we have in Eq. (2.5) the product of  $\Pi_{d+1}$  with some characteristic function instead of  $H_{d+1}$ . Since only  $(2d+1)$  diagonals of  $G_k$  have non-zero elements, one easily concludes that  $G_k$  converges strongly to  $(C(\zeta)+T_3)$ , where  $T_3 \in \mathcal{L}(l^2(\mathbb{C}^m))$  is a suitable operator of finite range. The appearance of operator  $T_3$  is explained by the fact that (2.5) does not hold for  $i=k_j$  and |u|, |v| < d; however  $T_3=0$  if  $i+k_i$ . Out of sup  $||G_k^{-1}|| < \infty$  and  $G_k \rightarrow (C(\zeta) + T_3)$  follows that  $G_k^{-1}$  converges strongly to  $(C(\zeta^{-1}) + T_4)$ , where  $T_4 \in \mathcal{L}(l^2(\mathbb{C}^m))$  is compact. Now a) is an easy consequence of (2.3).

For  $d < |u|$ ,  $|v| < \min\{|i - k_s| - d | s = 0, ..., N; k_s + i\}$  we obtain analogously to (2.3) and (2.5)

$$
D_k := (n(S_T \varphi_{r+i+jd+1}, \varphi_{s+i+jd+1}))_{s, r=-(i+jd)}^{(n-i)+(N-j)d-1}, \quad \tilde{S}_k = G_k^{-1} D_k
$$
  
\n
$$
\lim_{k \to \infty} (D_k)_{u, v} = |\gamma'(\sigma)| \int_R (S_R \Pi_{d+1}(\cdot - (v_1 - u_1))(s) \Pi_{d+1}(s) ds
$$
  
\n
$$
= -|\gamma'(\sigma)| \int_R e^{-i(u_1 - v_1)s} |F \Pi_{d+1}(s)|^2 \text{ sign } s ds
$$
  
\n
$$
= 1/2\pi \int_0^{2\pi} e^{-i(u_1 - v_1)s} \vartheta(s) ds,
$$
 (2.6)

where

$$
\vartheta(s) := -4^{(d+1)} |\gamma'(\sigma)| (\sin s/2)^{2(d+1)} \sum_{r \in \mathbb{Z}} \frac{\operatorname{sign}(r+1/2)}{(s+2r\pi)^{2(d+1)}}.
$$

The techniques of the proofs of Lemma 3.1 and Lemma 2.3 in [19] show that  $D_k \to C(3)$  for  $i \neq k_j$  (i.e.,  $\tau \notin \{t_0, ..., t_N\}$ ) and

$$
D_k \to (PC(\vartheta)P + QC(\vartheta)Q + PC(\vartheta_1)Q + QC(\vartheta_2)P + T_5) \quad \text{for } i = k_j,
$$

where  $\theta_1(s) = e^{ids}\theta(s)$ ,  $\theta_2(s) = e^{-ids}\theta(s)$ . The finite range operator  $T_5$  appears here, since (2.6) does not hold for  $i=k_i$  and |u|,  $|v| < d$ . For  $i=k_i$  and  $|u|$ ,  $|v| < d$ we have a formula analogous to (2.6), where again a product of  $\Pi_{d+1}$  with some characteristic function appears instead of  $H_{d+1}$ . Now  $PC(\vartheta_1)Q = PC(\vartheta)Q$  $+PC(\vartheta_3)Q$ , where

$$
\vartheta_3(s) = (e^{ids} - 1) \vartheta(s).
$$

Since  $\theta_3$  is a  $2\pi$ -periodic and continuous function on **R**, the operator  $PC(\theta_3)Q$ is compact. Consequently, we obtain  $D_k \rightarrow (C(3) + T_6)$ , where  $T_6 \in \mathcal{L}(l^2(\mathbb{C}^m))$  is compact. In any case, we have derived  $S_k = G_k^{-1}D_k \rightarrow (C(\rho) + T_2)$  for a suitable compact  $T_2 \in \mathcal{L}(l^2(\mathbb{C}^m))$ , and the proof of Lemma 2.1 is complete.

After defining

$$
A^{\tau} := P(c(\tau + 0) + d(\tau + 0) C(\rho)) + Q(c(\tau - 0) + d(\tau - 0) C(\rho)) \quad (\tau \in \Gamma),
$$
  
\n
$$
B^{\tau} := P(c(\tau + 0) + d(\tau + 0) C(\rho)) P \quad (\tau = \gamma(0)),
$$
  
\n
$$
C^{\tau} := Q(c(\tau - 0) + d(\tau - 0) C(\rho)) Q \quad (\tau = \gamma(1)),
$$

Lemma 2.1 yields  $\tilde{A}_k \rightarrow (A^{\tau} + T)$  for  $\sigma \in (0, 1)$ ,  $\tilde{A}_k \rightarrow (B^{\tau} + T)$  for  $\sigma = 0$  and  $\tilde{A}_k \rightarrow (C^{\tau})$ +T) for  $\sigma$ =1, where T is compact. If further  $\tau \notin \{t_0, ..., t_N\}$  or  $d = 0$ , then the proof of Lemma 2.1 yields  $T=0$ . The limit of  $\{\tilde{A}_k\}$  established above and the considerations preceeding Lemma 2.1 lead to

Lemma 2.2 (i) *Let F be a closed curve. If the Galerkin method with respect to an admissible sequence of equidistant partitions*  $\{\Lambda_k\}$  *converges for the operator A, then the operators*  $A^{\dagger}(\tau \in \Gamma)$  *are Fredholm operators with index 0.* 

(ii) *Let F be open. If the Galerkin method with respect to an admissible sequence of equidistant partitions*  $\{\Lambda_k\}$  *converges for the operator A, then the operators*  $A^{\gamma(\sigma)}(\sigma\in (0, 1))$ ,  $B^{\gamma(0)}$ ,  $C^{\gamma(1)}$  are Fredholm operators with index 0.

(iii) If in addition  $m=1$  or  $d=0$ , then the operators  $A^{\dagger}$ ,  $B^{\dagger}$ ,  $C^{\dagger}$  of (i), (ii) are *invertible.* 

*Proof of Theorem 2.1.* Let  $m = 1$ . If the Galerkin method converges, then Lemma 2.2 implies that operators  $A^{\tau}(\tau \in \Gamma)$  or  $A^{\gamma(\sigma)}$  ( $\sigma \in (0, 1)$ ),  $B^{\gamma(0)}$ ,  $C^{\gamma(1)}$  are invertible. Now one shows analogously to the proof of Corollary 4.2 in [19] that (1.4) is satisfied. Thus the strong ellipticity follows via Lemma 1.2.

Finally, we consider the collocation method described in [19]. Assume  $\Gamma$ closed and  $A_k := \{0 = \sigma_0, \sigma_1 = 1/k, \sigma_2 = 2/k, ..., \sigma_k = 1\}$  (k=n). We seek a piecewise linear approximation  $x_{4k}$  of the solution x of  $Ax = y$  satisfying

$$
(Ax_{4k})(\gamma(j/k)) = y(\gamma(j/k))(j=0,\ldots,k-1). \tag{2.7}
$$

Now define

$$
\psi_j^{(k)}(\gamma(s)) := \begin{cases}\n\frac{\gamma(s) - \gamma((j-1)/k)}{\gamma(j/k) - \gamma((j-1)/k)} & \text{if } s \in [(j-1)/k, j/k] \\
\frac{\gamma((j+1)/k) - \gamma(s)}{\gamma((j+1)/k) - \gamma(j/k)} & \text{if } s \in [j/k, (j+1)/k] \\
0 & \text{otherwise}\n\end{cases}
$$

 $k-1$ and  $K_k f := \sum f(\gamma(j/k)) \psi_i^{(k)}$ . Obviously,  $\{\psi_i^{(k)}\}_{i=0}^k$  is a base of the space  $j=0$  $\tilde{S}_1(\Lambda_k, \mathbb{C}^m)$  of all piecewise linear (linear in  $t = \gamma(s)$ ) functions subordinate to  $\Lambda_k$ and  $K_k$  is the interpolation projection onto this subspace. Equations (2.7) are equivalent to

$$
K_k A x_{\underline{A}_k} = K_k y.
$$

In [19] we proved that if A is invertible and the operators  $K_k A | \tilde{S}_1(\Lambda_k, \mathbb{C}^m)$ are stable, then  $x_{4}$ ( $k \rightarrow \infty$ ) converges to the solution x of  $Ax = y$  for all continuous right-hand sides y. We also gave a necessary and sufficient condition for stability of the aforementioned convergence scheme (see Theorem 4.1 in [19]). We now state

**Corollary 2.1.** Let  $a, b \in PC(\Gamma, \mathbb{C}^{m \times m})$ . If the invertible operator  $A = aP_r$  $+bO_r\in\mathscr{L}(L^2(\Gamma,\mathbb{C}^m))$  is strongly elliptic, then the approximate operators  $K_k A | \tilde{S}_1(A_k, \mathbb{C}^m)$  of the collocation method (2.7) are stable.

*Proof.* By Theorem 1.1, the Galerkin method with respect to any admissible sequence of equidistant partitions converges for A. Therefore, Lemma 2.2 implies the invertibility of the operators  $A^{\dagger}(d=0)$ . However, operators  $A^{\dagger}(d=0)$ coincide with the associated operators  $A_{\tau}$  of Theorem 4.1 in [19]. Theorem 4.1 of [19] yields Corollary 2.1.

## **3. The Asymptotical Order of Convergence**

For smooth curve *F*, let  $H^s(\Gamma)$  ( $s \in \mathbb{R}$ ,  $s \ge 0$ ) be the usual Sobolev space of order s on *F*. By  $H_{t_0,\,\ldots,\,t_N}^s(F)$  we denote the sum  $H^s[t_0,t_1]\oplus\ldots\oplus H^s[t_{N-1},t_N],$  where  $[t_i, t_{i+1}]$  is the arc on  $\Gamma$  with the end points  $t_i, t_{i+1}$ .

For  $\varepsilon$  an arbitrary positive number satisfying  $\varepsilon < \max(s_i - s_{i-1})$   $(i = 1, ..., N)$ , we choose the partitions  $A_k$  such that

$$
\Delta_k \cap (s_i - \varepsilon, s_i + \varepsilon) = \left\{ s_i \pm \left( \frac{j}{N_k} \right)^{\beta_i} \middle| j \in \mathbb{N} \right\} \cap (s_i - \varepsilon, s_i + \varepsilon), \quad i = 1, ..., N - 1
$$
\n
$$
\Delta_k \cap [0, \varepsilon) = \left\{ \left( \frac{j}{N_k} \right)^{\beta_0} \middle| j \in \mathbb{N} \right\} \cap [0, \varepsilon)
$$
\n
$$
\Delta_k \cap (1 - \varepsilon, 1] = \left\{ 1 - \left( \frac{j}{N_k} \right)^{\beta_N} \middle| j \in \mathbb{N} \right\} \cap (1 - \varepsilon, 1], \tag{3.1}
$$

where  $\beta_i(\beta_i \ge 1)$  is a given real number and  $N_k \in \mathbb{N}$ . In addition assume  $N_k \to \infty$  $(k \rightarrow \infty)$  and the existence of a positive constant C with

$$
h_{A_k} \le C 1/N_k, \qquad C^{-1} N_k \le n_k \le C N_k. \tag{3.2}
$$

If  $\beta_i = 1$  (i=0, ..., N), then equidistant partitions satisfy (3.1) and (3.2).

For functions  $a, b \in PC(\Gamma, \mathbb{C})$  continuous on  $\Gamma \setminus \{t_0, \ldots, t_N\}$ , we set

$$
\kappa_j := \frac{1}{2\pi i} \ln \left( \frac{a(t_j + 0) b(t_j - 0)}{a(t_j - 0) b(t_j + 0)} \right),
$$

where In denotes the continuous branch of the logarithm in  $\mathbb{C} \setminus (-\infty, 0]$  which takes real values on the positive real axis. Further, if  $\Gamma$  is open, set  $a(t_0-0)$  $= b(t_0 - 0) = a(t<sub>N</sub> + 0) = b(t<sub>N</sub> + 0) = 1.$ 

**Theorem 3.1.** *Let*  $\{A_k\}$  *satisfy* (3.1) *and* (3.2), *a, b*∈ $H_{t_0, ..., t_N}^{r+2}(F)$  (r∈**N**),  $1/2+r$  $+ \text{Re } \kappa_i < s < 3/2 + r + \text{Re } \kappa_i$  (i=0, ..., N),  $A = aP_f + bQ_f$  be invertible in  $L^2(\Gamma, \mathbb{C})$ *and strongly elliptic. If*  $y \in H_{to, ..., tx}^s(F)$ *,*  $x = A^{-1}y$  *and*  $x_{A_k}$  *is the solution of (0.3),* 

*then* 

$$
||x_{\Delta_k} - x||_{L^2(\Gamma, \mathbb{C})} \leq C_0 (1/n_k)^{\mu} (\ln n_k)^{\nu},
$$

where  $C_0$  is a constant and

$$
\mu := \min \{s, d+1, \beta_i (\text{Re } \kappa_i + 1/2) | i = 0, ..., N\},
$$
  

$$
v := \begin{cases} 1/2 & \text{if } s \ge d+1 \\ 0 & \text{otherwise.} \end{cases} \text{ and } \min \{\beta_i (\text{Re } \kappa_i + 1/2), i = 0, ..., N\} = d+1
$$

For simplicity sake, we shall prove Theorem 3.1 for the following special case only (see Theorem 3.3). Let  $\Gamma = (0, 1)$ ,  $a, b \in C[0, 1]$ . Thus instead of the "critical" points  $\{t_i | i = 0, ..., N\}$  we have merely the "critical" points  $\{t_0 = 0, t_N = 1\}$ . We now follow the ideas developed in [9]. First we determine the asymptotic behavior of  $x = A^{-1}y$  at 0 and 1. From the proofs in [9], Sect. 2, we conclude

**Theorem 3.2.** Assume  $a, b \in H^{r+2}(0, 1)$  (re**N**),  $A = aP_{(0, 1)} + bQ_{(0, 1)}$  invertible in  $L^2((0,1),\mathbb{C})$  and  $y \in H^s(0,1)$ , where  $\text{Re } \kappa_i + 1/2 + r < s < \text{Re } \kappa_i + 3/2 + r$  (i=0, 1). Then the function  $x = A^{-1}y$  has the representation

$$
x(t) = x_0(t) + t^{\kappa_0} \sum_{m=0}^{r} \sum_{s=0}^{m} \alpha_{m,s} t^m (\ln t)^s
$$
  
+ 
$$
(1-t)^{\kappa_1} \sum_{m=0}^{r} \sum_{s=0}^{m} \beta_{m,s} (1-t)^m (\ln (1-t))^s,
$$

*where*  $x_0 \in H^s(0, 1)$  *and*  $\alpha_{m, s}$ ,  $\beta_{m, s} \in \mathbb{C}$ .

If  $x_{4}$  is the solution of (0.3) and  $x = A^{-1}y$ , then

$$
||x - x_{A_k}|| \le C ||x - P_{A_k} x||. \tag{3.3}
$$

(See, e.g., [10].) Thus the proof of Theorem 3.1 for the immediately preceeding special case requires only estimating  $||(I-P_{A})x||$ .

**Lemma 3.1** (see, e.g., [9]). *If*  $x_0 \in H^s(0, 1)$ , *then there exists a constant C such that* 

$$
||(I - P_{A_k}) \times_0 ||_{L^2(\Gamma, \mathbb{C})} \leq C(1/n_k)^{\min(s, d+1)}.
$$

For further estimates, we need the quasi-interpolant  $Q_{\mu k}$  introduced by de Boor [7],

$$
Q_{A_k}g:=\sum_{i=-d}^{n_k-1}\lambda_i(g)\,\psi_i\,\sqrt{\frac{\sigma_{i+d+1}-\sigma_i}{d+1}},
$$

where  $\psi_i$  is defined as in Sect. 1 and

$$
\lambda_i(g) := \sum_{j=0}^d (-1)^{d-j} \omega_i^{(d-j)}(\tau_i) g^{(j)}(\tau_i) \qquad (i = -d, \dots, n_k - 1),
$$
  

$$
\omega_i(t) := (t - \sigma_{i+1}) \dots (t - \sigma_{i+d}) \cdot 1/d!,
$$
  

$$
\tau_i \in (\sigma_i, \sigma_{i+d+1}) \cap [1/2\sigma_1, 1].
$$
 (3.4)

The operator  $Q_{4k}$  is a projection onto  $S_d(\Lambda_k, \mathbb{C})$ .

**Lemma 3.2.** For any  $(d+1)$ -times continuously differentiable function f on  $(0, 1)$ , *the estimate* 

$$
\int_{\sigma_i}^{\sigma_{i+1}} |f - Q_{A_k} f|^2 \le C(\sigma_{i+d+1} - \sigma_{i-d-1})^{2(d+1)+1} \max_{t \in [\sigma_{i-d-1}, \sigma_{i+d+1}]} |f^{(d+1)}(t)|^2
$$
\n
$$
(i = 0, ..., n_k - 1)
$$
\n(3.5)

*holds, where C is a constant independent of f and*  $n_k$ *.* 

*Proof.* For  $t \in [\sigma_i, \sigma_{i+1}]$ , we obtain

$$
\left| f(t) - \sum_{j=0}^{d} f^{(j)}(\sigma_i) \frac{(t-\sigma_i)^j}{j!} \right| \le C(\sigma_{i+1} - \sigma_i)^{(d+1)} \max_{t \in [\sigma_i, \sigma_{i+1}]} |f^{(d+1)}(t)|
$$

in particular,

$$
\int_{\sigma_i}^{\sigma_{i+1}} \left| f(t) - \sum_{j=0}^d f^{(j)}(\sigma_i) \frac{(t-\sigma_i)^j}{j!} \right|^2 dt \le C(\sigma_{i+1} - \sigma_i)^{2(d+1)+1} \max_{t \in [\sigma_i, \sigma_{i+1}]} |f^{(d+1)}(t)|^2.
$$
\n(3.6)

On the other hand,

$$
\left| \sum_{j=0}^{d} f^{(j)}(\sigma_i) \frac{(t-\sigma_i)^j}{j!} - Q_{d_k} f(t) \right| = \left| Q_{d_k} \left( \sum_{j=0}^{d} f^{(j)}(\sigma_i) \frac{(\cdot-\sigma_i)^j}{j!} - f \right) (t) \right|
$$

$$
= \left| \sum_{s=i-d}^{i+1} \lambda_s \left( \sum_{j=0}^{d} f^{(j)}(\sigma_i) \frac{(\cdot-\sigma_i)^j}{j!} - f \right) \right| \sqrt{\frac{(\sigma_{s+d+1} - \sigma_s)}{d+1}} \psi_s(t)
$$

and (1.1) yields

$$
\int_{\sigma_1}^{\sigma_{i+1}} \left| \sum_{j=0}^d f^{(j)}(\sigma_i) \frac{(t-\sigma_i)^j}{j!} - Q_{A_k} f(t) \right|^2 dt
$$
\n
$$
\leq C \sum_{s=i-d}^{i+1} (\sigma_{i+d+1} - \sigma_{i-d-1}) \left| \lambda_s \left( \sum_{j=0}^d f^{(j)}(\sigma_i) \frac{(\cdot - \sigma_i)^j}{j!} - f \right) \right|^2.
$$
\n(3.7)

Setting  $g(t)=f(t)-\sum_{i=0}^{u} f^{(i)}(\sigma_i) \frac{(t-\sigma_i)^i}{i!}$ , we obtain  $j=0$ 

$$
|g^{(j)}(\tau_s)| \leqq C(\sigma_{i+d+1} - \sigma_{i-d-1})^{d+1-j} \max_{t \in [\sigma_{i-d-1}, \sigma_{i+d+1}]} |f^{(d+1)}(t)|
$$

and

$$
|\omega_s^{(d-j)}(\tau_s)| \leq C(\sigma_{i+d+1} - \sigma_{i-d-1})^j
$$
  $(s = i - d, ..., i + 1).$ 

Therefore, by means of (3.4) we have

$$
|\lambda_s(g)| \leq C(\sigma_{i+d+1} - \sigma_{i-d-1})^{d+1} \max_{t \in [\sigma_{i-d-1}, \sigma_{i+d+1}]} |f^{(d+1)}(t)|.
$$

From the last inequality and (3.7) we obtain

$$
\int_{\sigma_1}^{\sigma_{i+1}} \left| \sum_{j=0}^d f^{(j)}(\sigma_j) \frac{(t-\sigma_i)^j}{j!} - Q_{A_k} f(t) \right|^2 dt
$$
  
\n
$$
\leq C(\sigma_{i+d+1} - \sigma_{i-d-1})^{2(d+1)+1} \max_{t \in [\sigma_{i-d-1}, \sigma_{i+d+1}]} |f^{(d+1)}(t)|^2
$$

which together with  $(3.6)$  implies  $(3.5)$ .

**Lemma 3.3.** *Let*  $\{A_k\}$  *satisfy* (3.1) *and* (3.2) *and*  $f(t) := t^{x+i\beta}(\ln t)^r$  ( $\alpha, \beta \in \mathbb{R}, \alpha > 0$  $-1/2$ ,  $r \in \mathbb{N}$ ). Then there exists a constant C such that

$$
\| (I - P_{A_k}) f \|_{L^2((0, 1), \mathbb{C})} \le \| (I - Q_{A_k}) f \|_{L^2((0, 1), \mathbb{C})}
$$
  
\n
$$
\le C \begin{cases} \left(\frac{1}{n_k}\right)^{\beta_0(\alpha + \frac{1}{2})} (\ln n_k)^r & \text{if } \beta_0(\alpha + \frac{1}{2}) < d + 1 \\ \left(\frac{1}{n_k}\right)^{d+1} (\ln n_k)^{2r+\frac{1}{2}} & \text{if } \beta_0(\alpha + \frac{1}{2}) = d + 1 \\ \left(\frac{1}{n_k}\right)^{d+1} (\ln n_k)^{2r} & \text{if } \beta_0(\alpha + \frac{1}{2}) > d + 1 \end{cases}
$$
(3.8)

*Proof.* Considering Lemma 3.2,  $\sup_{t \in [a, 1]} |f^{(d+1)}(t)| < \infty$  and  $|\sigma_{i+d+1}| < \sigma$ .  $-\sigma_{i-d-1} \leq C \frac{1}{n_k}$  we can easily see that

$$
\int_{\varepsilon}^{1} |f - Q_{A_k} f|^2 \leq C \, 1/n_k^{2(d+1)}.
$$

Consequently, we may assume without loss of generality  $\varepsilon = 1$  (see (3.1)) - i.e.,  $\sigma_i = \sigma_i^k = (i/n_k)^{\beta_0}$   $(i = 0, 1, ..., n_k)$ .

Let  $i \geq 2(d+1)$ . This implies

$$
(\sigma_{i+d+1} - \sigma_{i-d-1}) \le C(i/n_k)^{\beta_0 - 1} 1/n_k,
$$
  
\n
$$
\max_{t \in [\sigma_{i-d-1}, \sigma_{i+d+1}]} |f^{(d+1)}(t)| \le C(i/n_k)^{\beta_0(\alpha - d - 1)} (\ln i)^r (\ln n_k)^r + C
$$

and Lemma 3.2 yields

$$
\int_{\sigma_{2(d+1)}}^1 |f - Q_{A_k} f|^2 \leq C \sum_{i=2(d+1)}^{n_k} (1/n_k)^{2(d+1)+1} (i/n_k)^{(\beta_0 - 1)(2(d+1)+1)}
$$

$$
\cdot \{(i/n_k)^{2\beta_0(\alpha - d - 1)} (\ln i)^{2r} (\ln n_k)^{2r} + 1 \},\
$$

$$
\left(\int_{\sigma_{2(d+1)}}^1 |f - Q_{A_k} f|^2\right)^{1/2} \leq C \begin{cases} (1/n_k)^{\beta_0(\alpha + 1/2)} (\ln n_k)^r & \text{if } \beta_0(\alpha + 1/2) < d + 1\\ (1/n_k)^{d+1} (\ln n_k)^{2r+1/2} & \text{if } \beta_0(\alpha + 1/2) = d + 1\\ (1/n_k)^{d+1} (\ln n_k)^{2r} & \text{if } \beta_0(\alpha + 1/2) > d + 1 \end{cases}
$$
(3.9)

For  $i \leq 2(d+1)$ , we have

$$
|\omega_s^{(d-j)}(\tau_s)| \le C(1/n_k)^{\beta_0+j}, \qquad |f^{(j)}(\tau_s)| \le C(1/n_k)^{\beta_0(\alpha-j)}(\ln n_k)^r.
$$

Now, (3.4) and (1.1) together yield

$$
\int_{0}^{\sigma_{2(d+1)}} |Q_{\Delta_{k}}f|^{2} \leq (\sigma_{3(d+1)} - \sigma_{0}) \sum_{s=0}^{2(d+1)} |\lambda_{s}(f)|^{2} \leq C(1/n_{k})^{\beta_{0}(2\alpha+1)} (\ln n_{k})^{2r},
$$

$$
\int_{0}^{\sigma_{2(d+1)}} |Q_{\Delta_{k}}f|^{2} \Big)^{1/2} \leq C(1/n_{k})^{\beta_{0}(\alpha+1/2)} (\ln n_{k})^{r}, \tag{3.10}
$$

whereas

$$
\left(\int_{0}^{\sigma_{2(d+1)}} |f|^{2}\right)^{1/2} \leq C(1/n_{k})^{\beta_{0}(\alpha+1/2)}(\ln n_{k})^{r}.
$$
\n(3.11)

Finally, (3.8) follows from (3.9), (3.10) and (3.11).

An immediate consequence of Theorem 3.2, (3.3), Lemma 3.1 and Lemma 3.3 is the following special case of Theorem 3.1.

**Theorem 3.3.** Let a,  $b \in H^{r+2}(0, 1)$  ( $r \in \mathbb{N}$ ),  $A = a P_{(0, 1)} + b Q_{(0, 1)}$  be invertible in  $L^2((0, 1), \mathbb{C})$  and strongly elliptic,  $y \in H^s(0, 1)$ , where  $\text{Re } \kappa_i + 1/2 + r < s < \text{Re } \kappa_i + 3/2$  $+r$  (*i*=0, 1), *and*  $\{\Delta_k\}$  *satisfy* (3.2) *and the second and third formula of* (3.1). *If* x  $A^{-1}$  *y* and  $x_{4k}$  denotes the solution of (0.3), then

$$
||x_{\Delta_k} - x||_{L^2((0,1),\mathbb{C})} \leq C(1/n_k)^{\mu} (\ln n_k)^{\nu},
$$

*where C is a constant and* 

$$
\mu := \min \{s, d+1, \beta_0 (\text{Re } \kappa_0 + 1/2), \beta_1 (\text{Re } \kappa_1 + 1/2) \}
$$
  

$$
v := \begin{cases} 1/2 & \text{if } s \ge d+1 \text{ and } \min \{ \beta_0 (\text{Re } \kappa_0 + 1/2), \beta_1 (\text{Re } \kappa_1 + 1/2) \} = d+1 \\ 0 & \text{otherwise.} \end{cases}
$$

*Remark.* Theorems 3.2 and 3.3 remain true for operators  $A = P_{(0, 1)}a + Q_{(0, 1)}b$ and functions  $y \in H^{s}(0, 1)$  satisfying

$$
y^{(k)}(0) = y^{(k)}(1) = 0, \quad k = 0, 1, ..., [s-1/2].
$$

Similarly, Theorem 3.1 can be proved for operators  $A = P<sub>r</sub> a + Q<sub>r</sub> b$  and functions  $y \in H^s(\Gamma)$ .

#### **4. A Numerical Example**

Let us consider the equation

$$
x(s) - 1/\pi \int_{-1}^{1} \frac{x(t)}{t - s} dt = 1
$$
 (4.1)

which has the exact solution

$$
x(s) = \frac{1}{\sqrt{2}} \left( \frac{1-s}{1+s} \right)^{1/4}.
$$

Using Galerkin's method we determine a piecewise constant approximate solution for x (i.e.  $d = 0$ ).

If we set  $a=1-i$ ,  $b=1+i$ , then the integral operator A defined by the left hand side of  $(4.1)$  takes the form  $(0.2)$ . Obviously, A is strongly elliptic (see Lemma 1.2) and with the notations of Sect. 3 we have

$$
\operatorname{Re}\kappa_0 = -1/4, \quad \operatorname{Re}\kappa_1 = 1/4.
$$

Thus, according to Sect. 3 we set

$$
\sigma_j := -1 + \left(\frac{j}{n_k/2}\right)^5, \quad \sigma_{n_k-j} := 1 - \left(\frac{j}{n_k/2}\right)^2, \quad j = 0, \dots, n_k/2,
$$

$$
\Delta_k := {\sigma_0, \dots, \sigma_{n_k}}.
$$

By exact computation of system (0.3) we obtain the following results.



If the logarithmic error  $\ln ||x_{A_k}-x||_{L^2}$  is considered as a function of  $\ln n_k$ , then the last table shows that the slope of this function tends to  $-1$ . This fact

confirms Theorem 3.3 which asserts  $||x_{4k}-x||_{L^2} \sim 1/n_k$ . In the next table we will compare the Galerkin approximation with the true solution at interior points.



In his paper [28] Thomas implemented and applied Galerkin's method to the same equation. By augmenting the space of piecewise linear functions on  $[-1+\delta_0, 1-\delta_1]$  with the singular functions  $(1+s)^{\kappa_0}$  on  $[-1,-1+\delta_0]$  and  $(1-s)^{\kappa_1}$  on  $[1-\delta_1, 1]$ , he formed a new trial space and obtained a better

Table 1

approximation. This modification could be useful also in calculation with splines of higher order. The reader can readily verify the same asymptotical order of convergence for the modified method.

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