

## A note on algebraic aspects of boundary feedback control systems of parabolic type

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**1. Introduction.** In the study of stabilization of boundary control systems, most fundamental is the static feedback control scheme: Based on a finite number of the observed data (outputs), it is the scheme to feed them back *directly* into the system through the boundary. Let  $\Omega$  denote a bounded domain of  $\mathbb{R}^m$  with the boundary  $\Gamma$  which consists of a finite number of smooth components of  $(m - 1)$ -dimension. The control system studied here is the following initial-boundary value problem:

$$(1) \quad \begin{aligned} \frac{\partial u}{\partial t} + \mathcal{L}u &= 0 \quad \text{in } (0, \infty) \times \Omega, \\ \tau u &= \sum_{k=1}^N \langle u, w_k \rangle_{\Omega} h_k \quad \text{on } (0, \infty) \times \Gamma, \\ u(0, \cdot) &= u_0(\cdot) \quad \text{in } \Omega. \end{aligned}$$

Here,  $\mathcal{L}$  denotes a uniformly elliptic differential operator of order 2 in  $\Omega$  defined by

$$\begin{aligned} \mathcal{L}u &= - \sum_{i,j=1}^m \frac{\partial}{\partial x_i} \left( a_{ij}(x) \frac{\partial u}{\partial x_j} \right) \\ &\quad + \sum_{i=1}^m b_i(x) \frac{\partial u}{\partial x_i} + c(x)u, \end{aligned}$$

and  $a_{ij}(x) = a_{ji}(x)$  for  $1 \leq i, j \leq m$ ,  $x \in \bar{\Omega}$ . The boundary operator  $\tau$  associated with  $\mathcal{L}$  is either  $\tau_1$  of the Dirichlet type or  $\tau_2$  of the Robin type:

$$\begin{aligned} \tau_1 u &= u|_{\Gamma}, \\ \tau_2 u &= \frac{\partial u}{\partial \nu} + \sigma(\xi)u \\ &= \sum_{i,j=1}^m a_{ij}(\xi) \nu_i(\xi) \frac{\partial u}{\partial x_j} \Big|_{\Gamma} + \sigma(\xi)u|_{\Gamma}, \end{aligned}$$

where  $(\nu_1(\xi), \dots, \nu_m(\xi))$  denotes the unit outer normal at  $\xi \in \Gamma$ . Necessary regularity on  $\bar{\Omega}$  and on  $\Gamma$  of coefficients of  $\mathcal{L}$  and  $\tau$  is assumed tacitly. The inner product and the norm in  $L^2(\Omega)$  are denoted by  $\langle \cdot, \cdot \rangle_{\Omega}$  and  $\| \cdot \|$ , respectively. The symbol  $\| \cdot \|$  is also used for the  $\mathcal{L}(L^2(\Omega))$ -norm. In eq. (1),  $\langle u, w_k \rangle_{\Omega}$

denote the outputs, where  $w_k \in L^2(\Omega)$ , and  $h_k$  the actuators belonging to  $H^{3/2}(\Gamma)$  in the case of the Dirichlet boundary condition, or  $H^{1/2}(\Gamma)$  in the Robin boundary condition.

Let us define the linear operators  $L_i$  and  $M_i$ ,  $i = 1, 2$  in  $L^2(\Omega)$  by

$$\begin{aligned} L_i u &= \mathcal{L}u, \quad u \in \mathcal{D}(L_i), \\ \mathcal{D}(L_i) &= \{u \in H^2(\Omega); \tau_i u = 0 \text{ on } \Gamma\} \end{aligned}$$

and

$$\begin{aligned} M_i u &= \mathcal{L}u, \quad u \in \mathcal{D}(M_i), \\ \mathcal{D}(M_i) &= \left\{ u \in H^2(\Omega); \right. \\ &\quad \left. \tau_i u = \sum_{k=1}^N \langle u, w_k \rangle_{\Omega} h_k \text{ on } \Gamma \right\}, \end{aligned}$$

respectively. Henceforth  $L$  stands for either  $L_1$  or  $L_2$  when it is distinguished from the context. The same symbolic convention applies to  $M_i$  as well as other operators. Eq. (1) is then simply rewritten as the equation in  $L^2(\Omega)$ :

$$(2) \quad \frac{du}{dt} + Mu = 0, \quad u(0) = u_0.$$

Given a  $\mu > 0$ , the problem is to find  $w_k$ 's and  $h_k$ 's such that the semigroup  $\exp(-tM)$  satisfies the decay estimate

$$(3) \quad \|e^{-tM}\| \leq \text{const } e^{-\mu t}, \quad t \geq 0.$$

In [4], this estimate was established via the fractional powers  $L_c^{\omega}$ ,  $L_c = L + c$ ,  $c > 0$  and the related fractional calculus. In the case of the Robin boundary condition, for example, we set

$$x(t) = L_{2c}^{-\omega} u(t), \quad \frac{1}{4} < \omega < \frac{1}{2},$$

and, noticing the relation:  $\mathcal{D}(L_{2c}^{\omega}) = H^{2\omega}(\Omega)$  for

$0 \leq \omega < 3/4$  [2], turn eq. (2) into

$$\frac{dx}{dt} + L_2x = \sum_{k=1}^N \langle L_{2c}^\omega x, w_k \rangle_\Omega L_{2c}^{1-\omega} \psi_k,$$

$$x(0) = L_{2c}^{-\omega} u_0,$$

where  $\psi_k \in H^2(\Omega)$  satisfy  $(\mathcal{L} + c)\psi_k = 0$ ,  $\tau_2\psi_k = h_k$ ,  $1 \leq k \leq N$ . The problem is then reduced to that of finding the estimate

$$\left\| \exp \left\{ -t \left( L_2 - \sum_{k=1}^N \langle L_{2c}^\omega \cdot, w_k \rangle_\Omega L_{2c}^{1-\omega} \psi_k \right) \right\} \right\|$$

$$\leq \text{const } e^{-\mu t}, \quad t \geq 0.$$

We propose in this note an alternative algebraic approach to the stabilization which requires no fractional powers of  $L_c$ . The common idea is, however, to turn the problem into another with no feedback term on  $\Gamma$ . A merit of the present approach is that the idea is equally applied to a variety of boundary control systems. In fact, the approach via fractional powers requires exact characterization of  $\mathcal{D}(L_c^\omega)$ . This seems in general a difficult (but challenging) problem when general elliptic operators with more complicated boundary conditions are studied.

The spectrum  $\sigma(L)$  consists only of eigenvalues  $\lambda_i$ ,  $i \geq 1$ , lying symmetrically in the interior of a parabola:  $\{\lambda = (a\tau^2 - b) + \sqrt{-1}\tau; \tau \in \mathbb{R}^1\}$ ,  $a > 0$  [1]. They are labelled according to increasing  $\text{Re } \lambda_i$ . As usual,  $P_{\lambda_i} = 1/(2\pi\sqrt{-1}) \int_{|\lambda - \lambda_i| = \varepsilon} (\lambda - L)^{-1} d\lambda$  is a projection which maps  $L^2(\Omega)$  onto the generalized eigenspace for  $\lambda_i$ , where  $\varepsilon > 0$  is small enough. Set  $\dim P_{\lambda_i} L^2(\Omega) = m_i (< \infty)$ , and let  $\varphi_{i1}, \dots, \varphi_{im_i}$  be the basis for  $P_{\lambda_i} L^2(\Omega)$ . As is well known [1],  $P_{\lambda_i}^*$  maps  $L^2(\Omega)$  onto the generalized eigenspace for  $\bar{\lambda}_i$  of  $L^*$ , and  $\dim P_{\lambda_i}^* L^2(\Omega) = m_i$ . The basis for  $P_{\lambda_i}^* L^2(\Omega)$  is denoted by  $\psi_{i1}, \dots, \psi_{im_i}$ .

For a given  $\mu > 0$ , suppose that

$$\text{Re } \lambda_1 \leq \dots \leq \text{Re } \lambda_K \leq \mu < \text{Re } \lambda_{K+1}.$$

Set  $P = P_{\lambda_1} + \dots + P_{\lambda_K}$ . In view of the expression:  $L\varphi_{ij} = \lambda_i\varphi_{ij} + \sum_{k < j} \alpha_{jk}^i \varphi_{ik}$ ,  $1 \leq j \leq m_i$ , the restriction of  $L$  onto the invariant subspace  $PL^2(\Omega)$  is bounded and similar to the upper triangular matrix  $A$ , the diagonal elements of which are  $\underbrace{\lambda_1, \dots, \lambda_1}_{m_1}, \dots, \underbrace{\lambda_K, \dots, \lambda_K}_{m_K}$ . If  $\lambda$  is in  $\rho(L_i)$ , the boundary value

problem:

$$(\lambda - \mathcal{L})\psi_k = 0, \quad \tau_i\psi_k = h_k,$$

$$1 \leq k \leq N, \quad i = 1, 2$$

admits a unique solution  $\psi_k$  [3] which is denoted by  $N_i(\lambda)h_k$ , where

$$N_1(\lambda) \in \mathcal{L}(H^{3/2}(\Gamma); H^2(\Omega)),$$

$$N_2(\lambda) \in \mathcal{L}(H^{1/2}(\Gamma); H^2(\Omega)).$$

The operators  $N_i(\lambda)$  are simply rewritten as  $N(\lambda)$ .

**2. Main result.** Our first result is

**Theorem 2.1.**

- (i) *The operator  $M$  is densely defined. The problem (2) is well posed, and the semigroup  $e^{-tM}$  is analytic in  $t > 0$ .*
- (ii) *The adjoint  $M^*$  is given by*

$$M_1^*u = \mathcal{L}^*u + \sum_{k=1}^N \left\langle \frac{\partial u}{\partial \nu}, h_k \right\rangle_\Gamma w_k,$$

$$u \in \mathcal{D}(M_1^*) = H^2(\Omega) \cap H_0^1(\Omega),$$

$$M_2^*u = \mathcal{L}^*u - \sum_{k=1}^N \langle u, h_k \rangle_\Gamma w_k,$$

$$u \in \mathcal{D}(M_2^*) = \{u \in H^2(\Omega); \tau^*u = 0 \text{ on } \Gamma\},$$

where  $(\mathcal{L}^*, \tau^*)$  denotes the formal adjoint of  $(\mathcal{L}, \tau)$ .

For notational convenience, let us introduce the symbol  $[u]$  as

$$[u] = \begin{cases} \frac{\partial u}{\partial \nu}, & \text{in the case of the Dirichlet boundary} \\ & \text{condition,} \\ u, & \text{in the case of the Robin boundary} \\ & \text{condition.} \end{cases}$$

Then  $M_i^*$  are simply rewritten as

$$M_i^*u = \mathcal{L}^*u - (-1)^i \sum_{k=1}^N \langle [u], h_k \rangle_\Gamma w_k, \quad i = 1, 2.$$

For a large  $c > 0$  with  $-c \in \rho(L)$ , set  $PN(-c)h_k = \sum_{i \leq K, j} \zeta_{ij}^k \varphi_{ij}$ . It is well known -via Green's formula- that there is an  $S \times S$  nonsingular matrix  $A$  such that ( $S = m_1 + \dots + m_K$ )

$$\begin{pmatrix} \zeta_{11}^k \\ \vdots \\ \zeta_{Km_K}^k \end{pmatrix} = A \begin{pmatrix} \langle h_k, [\psi_{11}] \rangle_\Gamma \\ \vdots \\ \langle h_k, [\psi_{Km_K}] \rangle_\Gamma \end{pmatrix}.$$

We define the  $S \times S$  matrix  $\tilde{A}$  and the  $S \times N$  matrix  $H$  as

$$(4) \quad \tilde{A} = A^{-1}AA,$$

and

$$(5) \quad H = \left( \langle h_k, [\psi_{ij}]_T; \begin{matrix} k \rightarrow 1, \dots, N \\ (i, j) \downarrow (1, 1), \dots, (K, m_K) \end{matrix} \right),$$

respectively.

Based on Theorem 2.1, our main result is stated as follows:

**Theorem 2.2.** *Suppose that  $(\tilde{A}, H)$  is a controllable pair, i.e.,*

$$(6) \quad \text{rank}(H \tilde{A}H \tilde{A}^2H \dots \tilde{A}^{S-1}H) = S.$$

Then there is a set of  $w_k$ 's  $\in P^*L^2(\Omega)$  such that the estimate (3) holds.

*Outline of the proof.* The proof of Theorem 2.1, (i) is almost the same as in [5, Theorem 2.3]: There exists a sector  $\bar{\Sigma}_\alpha = \{\lambda - \alpha \in \mathbb{C}; \theta_0 \leq |\arg \lambda| \leq \pi\}$ ,  $0 < \theta_0 < \pi/2$ ,  $\alpha \in \mathbb{R}^1$ , such that

$$(7) \quad (\lambda - M)^{-1}f = (\lambda - L)^{-1}f + [N(\lambda)h_1 \dots N(\lambda)h_N](1 - \Phi(\lambda))^{-1} \cdot \langle (\lambda - L)^{-1}f, \mathbf{w} \rangle_\Omega, \quad \lambda \in \bar{\Sigma}_\alpha,$$

where  $\langle \cdot, \mathbf{w} \rangle_\Omega$  denotes the transpose of a vector:

$$(\langle \cdot, w_1 \rangle_\Omega \dots \langle \cdot, w_N \rangle_\Omega),$$

and

$$\Phi(\lambda) = \left( \langle N(\lambda)h_k, w_j \rangle_\Omega; \begin{matrix} k \rightarrow 1, \dots, N \\ j \downarrow 1, \dots, N \end{matrix} \right) \rightarrow 0, \quad |\lambda| \rightarrow \infty, \quad \lambda \in \bar{\Sigma}_\alpha,$$

uniformly. Thus the estimate:

$$\|(\lambda - M)^{-1}\| \leq \frac{\text{const}}{1 + |\lambda|}, \quad \lambda \in \bar{\Sigma}_\alpha$$

holds, and  $e^{-tM}$  is analytic in  $t > 0$ .

The expression of the adjoint  $M_1^*$  is found in [5, Proposition 2.4], and  $M_2^*$  is obtained in almost the same manner as  $M_1^*$ .

As to the proof of Theorem 2.2, the main feature is to propose an approach entirely different from and simpler than in [4]. Let us define the operator  $T$  by

$$(8) \quad v = Tu = u - \sum_{k=1}^N \langle u, w_k \rangle_\Omega N(-c)h_k.$$

Here the vectors  $w_j$ 's are to be determined later in relation to  $c > 0$  and the associated finite-dimensional

stabilization problem (12a). The operator  $T$  belongs to  $\mathcal{L}(L^2(\Omega)) \cap \mathcal{L}(\mathcal{D}(M); \mathcal{D}(L))$ . The bounded inverse  $T^{-1}$  exists, and is given by

$$u = T^{-1}v = v + [N(-c)h_1 \dots N(-c)h_N] \cdot (1 - \Phi(-c))^{-1} \langle v, \mathbf{w} \rangle_\Omega.$$

Here we have assumed with no loss of generality that  $(1 - \Phi(-c))^{-1}$  exists. In fact, consider the case where  $\det(1 - \Phi(-c)) = 0$ . We then replace  $w_j$ 's by  $(1 + \varepsilon)w_j$ 's for a sufficiently small  $\varepsilon$ . The function  $\det(1 - (1 + \varepsilon)\Phi(-c))$  in  $\varepsilon$  is a polynomial of degree at most  $N$ ; not a constant; and analytic. Thus  $\det(1 - (1 + \varepsilon)\Phi(-c)) \neq 0$  for some small  $\varepsilon \neq 0$ . As far as  $\varepsilon$  is small enough, this does not affect the stabilization problem under consideration. The other properties of  $T$  are easily examined.

For a solution  $u \in \mathcal{D}(M)$  to the problem (2), set

$$(9) \quad v(t) = Tu(t), \quad t \geq 0.$$

Then  $v(t) \in \mathcal{D}(L)$  satisfies the equation

$$\frac{dv}{dt} + TM_cT^{-1}v = cv, \quad t > 0, \quad v(0) = Tu_0,$$

where  $M_c = M + c$ . We calculate as

$$\begin{aligned} TM_cT^{-1}v &= T\mathcal{L}_c(v + [N(-c)h_1 \dots N(-c)h_N] \cdot (1 - \Phi(-c))^{-1} \langle v, \mathbf{w} \rangle_\Omega) \\ &= T\mathcal{L}_cv = TL_cv \\ &= L_cv - \sum_{k=1}^N \langle L_cv, w_k \rangle_\Omega N(-c)h_k. \end{aligned}$$

We assume that  $w_k$ 's belong to  $P^*L^2(\Omega) \subset \mathcal{D}(L^*)$ . Then the equation for  $v$  is rewritten as

$$(10) \quad \frac{dv}{dt} + Lv - \sum_{k=1}^N \langle v, L_c^*w_k \rangle_\Omega N(-c)h_k = 0, \quad t \geq 0, \quad v(0) = Tu_0.$$

The problem (10) generates an analytic semigroup. Thus the problem (2) also generates an analytic semigroup  $\exp(-tM)$ , and

$$(11) \quad \exp(-tM) = T^{-1} \cdot \exp \left\{ -t \left( L - \sum_{k=1}^N \langle \cdot, L_c^*w_k \rangle_\Omega N(-c)h_k \right) \right\} \cdot T, \quad t \geq 0.$$

In view of the relation (11), we have to establish a stabilization result for the problem (10). At this stage, the problem is simple since  $w_k$ 's belong

to  $P^*L^2(\Omega)$ . The restrictions of  $L$  onto the invariant subspaces  $PL^2(\Omega)$  and  $(1-P)L^2(\Omega) \cap \mathcal{D}(L)$  are denoted by  $L^1$  and  $L^2$  respectively. Then, by setting

$$v_1 = Pv, \quad v_2 = (1-P)v,$$

eq. (10) is decomposed into

$$(12a) \quad \frac{dv_1}{dt} + L^1v_1 - \sum_{k=1}^N \langle v_1, L_c^*w_k \rangle_{\Omega} PN(-c)h_k = 0,$$

$$(12b) \quad \frac{dv_2}{dt} + L^2v_2 - \sum_{k=1}^N \langle v_1, L_c^*w_k \rangle_{\Omega} (1-P)N(-c)h_k = 0.$$

In (12a), replace  $L_c^*w_k$  by  $y_k = \sum_{i,j (i \leq K)} y_{ij}^k \psi_{ij}$ . Then (12a) is equivalent to the equation in  $\mathbb{C}^S$ :

$$\frac{dv}{dt} + (\Lambda - Z\bar{Y}\Pi)v = 0,$$

where

$$Z = \left( \zeta_{ij}^k; \begin{matrix} k \rightarrow 1, \dots, N \\ (i, j) \downarrow (1, 1), \dots, (K, m_K) \end{matrix} \right) = AH,$$

$$Y = \left( y_{ij}^k; \begin{matrix} k \downarrow 1, \dots, N \\ (i, j) \rightarrow (1, 1), \dots, (K, m_K) \end{matrix} \right), \quad \text{and}$$

$$\Pi = \left( \langle \varphi_{ij}, \psi_{pq} \rangle_{\Omega}; \begin{matrix} (i, j) \rightarrow (1, 1), \dots, (K, m_K) \\ (p, q) \downarrow (1, 1), \dots, (K, m_K) \end{matrix} \right).$$

Note that  $\Pi$  is nonsingular and  $\langle \varphi_{ij}, \psi_{pq} \rangle_{\Omega} = 0$  when  $i \neq p$ . According to the assumption (6),  $(\Lambda, Z)$  is a controllable pair, i.e.,

$$\text{rank}(Z \Lambda Z \Lambda^2 Z \dots \Lambda^{S-1} Z) = S.$$

Thus the well known pole assignment argument of finite dimension [6] implies that there exists an  $N \times S$

matrix  $Y$  or  $w_k$ 's in  $P^*L^2(\Omega)$  such that

$$\|e^{-t(\Lambda - Z\bar{Y}\Pi)}\| \leq \text{const } e^{-\mu t}, \quad t \geq 0.$$

By recalling that  $\|e^{-tL^2}\| \leq \text{const } e^{-\mu' t}$ ,  $t \geq 0$ ,  $\mu < \mu' < \text{Re } \lambda_{K+1}$ , (12b) immediately gives the desired estimate for  $v$ . Note that  $\mu'$  cannot be generally replaced by  $\text{Re } \lambda_{K+1}$ , due to the algebraic multiplicities of the eigenvalues on the vertical line:  $\text{Re } \lambda = \text{Re } \lambda_{K+1}$ .  $\square$

As a concluding remark, another algebraic approach to Theorem 2.2 is possible via Theorem 2.1, (ii). In view of the relation

$$\|e^{-tM}\| = \|(e^{-tM})^*\| = \|e^{-tM^*}\|,$$

the problem is reduced to the one with  $M^*$ , and the assumption (6) ensures suitable vectors  $w_k$ 's in  $P^*L^2(\Omega)$ .

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