

## NON-LOCAL ELLIPTIC PROBLEM IN HIGHER DIMENSION

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### Abstract

Non-local elliptic problem,  $-\Delta v = \lambda(e^v / (\int_{\Omega} e^v dx)^p)$  with Dirichlet boundary condition is considered on  $n$ -dimensional bounded domain  $\Omega$  with  $n \geq 3$  for  $p > 0$ . If  $\Omega$  is the unit ball,  $3 \leq n \leq 9$  and  $2/n \leq p \leq 1$ , we have infinitely many bendings in  $\lambda$  of the solution set in  $\lambda - v$  plane. Finally if  $\Omega$  is an annulus domain and  $p \geq 1$ , we show that a solution exists for all  $\lambda > 0$ .

### 1. Introduction

In this paper we consider the following elliptic equation with non-local term:

$$(1) \quad \begin{cases} -\Delta v = \lambda \frac{e^v}{(\int_{\Omega} e^v dx)^p} & x \in \Omega, \\ v = 0 & x \in \partial\Omega, \end{cases}$$

where  $\lambda, p$  are positive constants and  $\Omega$  is a bounded domain in  $\mathbf{R}^n$  with smooth boundary  $\partial\Omega$ . Actually, usual Gel'fand problem, in the theories of thermonic emission ([5]), isothermal gas sphere ([4]), and gas combustion ([1]), is formulated as the nonlinear eigenvalue problem

$$(2) \quad \begin{cases} -\Delta v = \sigma e^v & x \in \Omega, \\ v = 0 & x \in \partial\Omega, \end{cases}$$

with a constant  $\sigma > 0$ . Problems (1) and (2) are equivalent through the relation

$$\sigma = \frac{\lambda}{(\int_{\Omega} e^v dx)^p},$$

and hence some features of the solution set

$$\mathcal{C} = \{(\lambda, v) \mid v = v(x) \text{ is a classical solution of (1) for } \lambda > 0\}$$

resemble those of the solution set for (2), denoted by  $\mathcal{S}$ . (1) is the non-local stationary problem of

$$\begin{cases} v_t = \Delta v + \lambda \frac{e^v}{\left(\int_{\Omega} e^v dx\right)^p} & x \in \Omega, t > 0, \\ v = 0 & x \in \partial\Omega, t > 0, \\ v|_{t=0} = v_0(x) & x \in \Omega. \end{cases}$$

Such problems are studied in ([2]). They arise in the study of phenomena associated with the occurrence of shear bands in metals being deformed under high strain rates ([3]) and Ohmic heating ([10], [11]). We note that if  $p = 1$ , the motivation to study (1) is the Keller-Segel system ([9]) which describes the chemotactic aggregation of cellular slime molds given by

$$(3) \quad \begin{cases} \varepsilon u_t = \nabla \cdot (\nabla u - u \nabla v) & x \in \Omega, t \in (0, T), \\ \tau v_t = \Delta v + u & x \in \Omega, t \in (0, T), \\ \frac{\partial u}{\partial \nu} - u \frac{\partial v}{\partial \nu} = v = 0 & x \in \partial\Omega, t \in (0, T), \\ u|_{t=0} = u_0(x) \geq 0 & x \in \Omega, \\ v|_{t=0} = v_0(x) & x \in \Omega, \end{cases}$$

where  $\Omega$  is a bounded domain with smooth boundary  $\partial\Omega$ ,  $\tau, \varepsilon$  are positive constants, and  $\nu$  is the outer unit normal vector, respectively. In the stationary state of (3), it is reduced to (1) ([20]). In fact, since  $u = u(x, t) \geq 0$  and

$$\frac{d}{dt} \int_{\Omega} u dx = \frac{1}{\varepsilon} \int_{\Omega} \nabla \cdot (\nabla u - u \nabla v) dx = 0,$$

the total mass is conserved, that is,  $\|u(\cdot, t)\|_1 = \|u_0\|_1$ . Here and henceforth,  $\|\cdot\|_p$  denotes the standard  $L^p$  norm. The system of equations (3) has a Lyapunov function

$$J(u, v) = \int_{\Omega} \left( u(\log u - 1) - uv + \frac{1}{2} |\nabla v|^2 \right) dx$$

and it holds that

$$\varepsilon \frac{d}{dt} J(u, v) + \tau \varepsilon \|v_t\|_2^2 + \int_{\Omega} u |\nabla(\log u - v)|^2 dx = 0.$$

It implies that in the stationary state

$$\log u - v = \log \sigma$$

holds for some constant  $\sigma > 0$ . In other words

$$u = \sigma e^v$$

holds. Putting  $\lambda = \|u_0\|_1$  we have

$$\sigma = \frac{\lambda}{\int_{\Omega} e^v dx}$$

by mass conservation. Thus the second equation in (3) implies that

$$0 = \Delta v + \lambda \frac{e^v}{\int_{\Omega} e^v dx}$$

with Dirichlet boundary condition.

We have known the result on  $\mathcal{C}$  when  $n = 1, 2$  and  $\Omega$  is the unit ball  $B = \{x \in \mathbf{R}^n \mid |x| < 1\}$ . For  $n = 1$  with  $\Omega = B$ , if  $p \geq 1$ , then (1) has a unique solution for all  $\lambda > 0$ . On the contrary if  $0 < p < 1$ , there exists  $\bar{\lambda} > 0$  such that (1) has two solutions for  $\lambda < \bar{\lambda}$ , (1) has one solution for  $\lambda = \bar{\lambda}$  and (1) has no solution for  $\lambda > \bar{\lambda}$ . For  $n = 2$  with  $\Omega = B$ , if  $p < 1$ , then (1) has a unique solution for all  $\lambda > 0$ . If  $p = 1$ , then (1) has a unique solution for all  $0 < \lambda < 8\pi$  but no solution for  $\lambda \geq 8\pi$ . On the contrary if  $0 < p < 1$ , there exists  $\bar{\lambda} > 0$  such that (1) has two solutions for  $\lambda < \bar{\lambda}$ , (1) has one solution for  $\lambda = \bar{\lambda}$  and (1) has no solution for  $\lambda > \bar{\lambda}$ . These facts are proven in [2].

We consider the  $p = 1$  case in [13] and expand their results to the general  $p > 0$  in this paper.

We also mention the structure of  $\mathcal{S}$ . For  $n \geq 3$  with  $\Omega = B$ ,  $\mathcal{S}$  is a one-dimensional open manifold with the end points in  $(\sigma, v) = (0, 0)$  and  $(\sigma, v) = (2(n-2), 2 \log(1/|x|))$ , respectively, the latter being a weak solution of (2). Moreover for  $3 \leq n \leq 9$ ,  $\mathcal{S}$  bends infinitely many times with respect to  $\sigma$  around  $\sigma = 2(n-2)$ . Morse indices increases by one whenever it bends. On the other hand if  $n \geq 10$ , no bending occurs. They are shown in [14], [15].

We define the section of  $\mathcal{C}$  cut by  $\lambda > 0$  as follows:

$$\mathcal{C}^\lambda = \{v \in C^2(\Omega) \cap C(\bar{\Omega}) \mid v = v(x) \text{ solves (1)}\}.$$

The first theorem is concerned with the star-shaped domain, so that  $x \cdot \nu > 0$  holds for each  $x \in \partial\Omega$ .

**Theorem 1.** *If  $\Omega$  is star-shaped with respect to the origin with  $n \geq 3$  and  $p \leq 1$ , then there is  $\bar{\lambda} \in (0, +\infty)$  such that (1) has no solution for  $\lambda > \bar{\lambda}$ . Moreover,  $\mathcal{C}_0$  is unbounded in  $\lambda - v$  plane, where  $\mathcal{C}_0$  stands for the connected component of  $\mathcal{C}$  satisfying  $(0, 0) \in \bar{\mathcal{C}}_0$ .*

The second theorem is concerned with the ball case.

**Theorem 2.** *If  $\Omega$  is the unit ball  $B = \{x \in \mathbf{R}^n \mid |x| < 1\}$  with  $n \geq 3$ , then  $\mathcal{C}$  is a one-dimensional open manifold and can be parametrized as*

$$\mathcal{C} = \{(\lambda(s), v(\cdot, s)) \mid 0 < s < +\infty\}$$

with the end points in  $(\lambda, v) = (0, 0)$  and the weak solution

$$(\lambda, v) = \left( 2\omega_n^p (n-2)^{1-p}, 2 \log\left(\frac{1}{|x|}\right) \right),$$

respectively, where  $\omega_n$  denotes the area of the unit sphere in  $\mathbf{R}^n$ . Moreover for  $3 \leq n \leq 9$  and  $2/n \leq p \leq 1$ ,  $\mathcal{C}$  bends infinitely many times with respect to  $\lambda$  around  $\lambda = 2\omega_n^p (n-2)^{1-p}$ . On the other hand if  $n \geq 10$  and  $p \leq 1$ , no bending occurs.

The third theorem is on the spectral property of the linearized operator. To state the result, we define Morse index as follows. For given  $(\lambda, v) \in \mathcal{C}$ , the linearized eigenvalue problem is given by

$$(4) \quad \begin{cases} \Delta \phi + \lambda \frac{e^v}{\left(\int_{\Omega} e^v dx\right)^p} \phi - p\lambda \frac{\int_{\Omega} e^v \phi dx}{\left(\int_{\Omega} e^v dx\right)^{p+1}} e^v = -\mu \phi & x \in \Omega, \\ \phi = 0 & x \in \partial\Omega. \end{cases}$$

Then, the Morse index  $i = i(\lambda, v)$  and the radial Morse index  $i_R = i_R(\lambda, v)$  denote the number of negative eigenvalues and that of radially symmetric eigenfunctions, respectively.

**Theorem 3.** *If  $\Omega$  is the unit ball with  $3 \leq n \leq 9$ ,  $2/n \leq p \leq 1$  and  $n \geq 10$ ,  $p \leq 1$ , respectively, then  $i = i_R$  holds and  $i = i(\lambda, v)$  increases by one at each bending point.*

The last theorem is on the annulus domain  $A_a = \{x \in \mathbf{R}^n \mid a < |x| < 1\}$  with  $a \in (0, 1)$ . We deal with only radial solutions. Then we define the solution set by

$$\mathcal{C}_a = \{(\lambda, v) \mid v = v(|x|) \text{ is a classical solution of (1) for } \lambda > 0\}.$$

**Theorem 4.** *If  $\Omega$  is the annulus domain  $A_a$  with  $n \geq 3$  and  $p \geq 1$ , then  $\mathcal{C}_a$  is a one-dimensional open manifold and can be parametrized as*

$$\mathcal{C}_a = \{(\lambda(s), v(\cdot, s)) \mid 0 < s < +\infty\}$$

with the end points in  $(\lambda, v) = (0, 0)$  and  $(\lambda, v)$  satisfying

$$\lim_{s \uparrow +\infty} \lambda(s) = +\infty \quad \text{and} \quad \lim_{s \uparrow +\infty} \sup_{a < x < 1} |v(x, s)| = +\infty.$$

This paper is composed of four sections. In §2, we treat a star-shaped domain and prove Theorem 1. Next in §3, we study the ball case and prove Theorems 2 and 3. Finally, §4 is on the annulus domain case and we prove Theorem 4.

**2. Star-shaped domain**

In this section, we assume that  $\Omega$  is a star-shaped bounded domain with respect to the origin in  $\mathbf{R}^n$  with  $n \geq 3$  with the smooth boundary  $\partial\Omega$  and that  $\nu$  is the outer unit normal vector.

Proof of Theorem 1. The first part of Theorem has already proven in [2], but we provide the proof for completeness. In fact we apply the Pohozaev identity ([16]) to (1).

(5)

$$\begin{aligned} \frac{1}{2} \int_{\partial\Omega} (x \cdot \nu) \left( \frac{\partial v}{\partial \nu} \right)^2 ds &= \frac{n\lambda}{\left( \int_{\Omega} e^v dx \right)^p} \int_{\Omega} (e^v - 1) dx + \frac{2-n}{2} \frac{\lambda}{\left( \int_{\Omega} e^v dx \right)^p} \int_{\Omega} e^v v dx \\ &\leq \frac{n\lambda}{\left( \int_{\Omega} e^v dx \right)^{p-1}}, \end{aligned}$$

where  $ds$  is the area element of  $\partial B$  with standard metric.

On the other hand it follows from (1) that

$$\frac{\lambda}{\left( \int_{\Omega} e^v dx \right)^{p-1}} = \int_{\Omega} (-\Delta v) dx = \int_{\partial\Omega} \left( -\frac{\partial v}{\partial \nu} \right) ds,$$

and therefore we have

$$\frac{\lambda^2}{\left( \int_{\Omega} e^v dx \right)^{2(p-1)}} \leq \int_{\partial\Omega} (x \cdot \nu) \left( \frac{\partial v}{\partial \nu} \right)^2 ds \int_{\partial\Omega} \frac{1}{(x \cdot \nu)} ds.$$

Combining this inequality and (5), we have

$$\frac{\lambda^2}{\left( \int_{\Omega} e^v dx \right)^{2(p-1)}} \leq \frac{2n\lambda}{\left( \int_{\Omega} e^v dx \right)^{p-1}} \int_{\partial\Omega} \frac{1}{(x \cdot \nu)} ds.$$

Hence since  $p \leq 1$  and  $v > 0$  in  $\Omega$ , we have

$$\lambda \leq \frac{2n}{\left( \int_{\Omega} e^v dx \right)^{1-p}} \int_{\partial\Omega} \frac{1}{(x \cdot \nu)} ds \leq \frac{2n}{|\Omega|^{1-p}} \int_{\partial\Omega} \frac{1}{(x \cdot \nu)} ds,$$

where  $|\Omega|$  is the measure of  $\Omega$ . This gives  $\bar{\lambda}$  in the statement.

We denote by  $\mathcal{C}_1$  the branch of solutions of (1) starting from  $(\lambda, v) = (0, 0)$ . Supposing that  $\mathcal{C}_1$  is bounded, we prove unboundedness of the component  $\mathcal{C}_1$  by contradiction by making use of the standard degree argument similar to [18]. We provide the proof for completeness and proceed it in the same method as in [19]. Putting

$$F(\lambda, v) = \Delta v + \lambda \frac{e^v}{\left(\int_{\Omega} e^v dx\right)^p},$$

we apply the implicit function theorem. Then there exists a solution  $(\lambda, v)$  for  $0 < \lambda \ll 1$ .  $\mathcal{C}_1$  is compact from the assumption and the existence of the upper bound  $\bar{\lambda} < +\infty$ . Because

$$\mathcal{C} \cap \{\lambda = 0\} = \{(0, 0)\},$$

we can take open set  $\mathcal{U}$  containing  $\mathcal{C}_1$  with the properties

$$\partial\mathcal{U}_\lambda \cap \mathcal{C} = \emptyset \quad \text{and} \quad \mathcal{U}_\lambda = \emptyset \quad \text{for } \lambda \gg 1,$$

and  $\mathcal{U}_\lambda \cap \mathcal{C}$  is composed of the solution of (1) with  $\mu_1(\lambda, v) > 0$  for  $0 < \lambda \ll 1$ , where

$$\mathcal{U}_\lambda = \{v \in C(\overline{\Omega}) \mid (\lambda, v) \in \mathcal{U}\}.$$

In Banach space  $C(\overline{\Omega})$ , Leray-Schauder degree  $d(\Psi_\lambda, 0, \mathcal{U}_\lambda)$  is taken for any  $\lambda > 0$ , where  $\Psi_\lambda = I_{C(\overline{\Omega})} - \Phi_\lambda$  with

$$\Phi_\lambda(v) = (-\Delta)^{-1} \lambda \frac{e^v}{\left(\int_{\Omega} e^v dx\right)^p}.$$

From the homotopy invariance ([18]),  $d(\Psi_\lambda, 0, \mathcal{U}_\lambda)$  is independent of  $\lambda > 0$ . However by existence and nonexistence of the solution of (1), we have

$$\begin{cases} d(\Psi_\lambda, 0, \mathcal{U}_\lambda) = 0 & \text{for } \lambda \gg 1, \\ d(\Psi_\lambda, 0, \mathcal{U}_\lambda) = 1 & \text{for } 0 < \lambda \ll 1, \end{cases}$$

which is a contradiction. □

### 3. Ball case

In this section, we assume that  $\Omega = B$ , where  $B = \{x \in \mathbf{R}^n \mid |x| < 1\}$ .

Proof of Theorem 2. According to [6], any solution of (1) is radially symmetric. Hence we have

$$\begin{cases} (r^{n-1}v')' + \lambda r^{n-1} \frac{e^v}{\left(\int_{\Omega} e^v dx\right)^p} = 0 & \text{for } r > 0, \\ v(1) = 0, \quad v'(0) = 0, \end{cases}$$

where

$$v = v(r) \quad \text{for } r = |x|.$$

We begin with the parametrization of the solution set  $\mathcal{C}$ , following [7], [13], [14] and [15]. In fact, any solution is obtained as a solution of the initial value problem

$$(6) \quad \begin{cases} (r^{n-1}v')' + \sigma r^{n-1}e^v = 0 & \text{for } r > 0, \\ v(0) = A, \quad v'(0) = 0, \end{cases}$$

with a certain positive constant  $A$ . Through the Emden transformation

$$(7) \quad v(r) = w(t) - 2t + A, \quad r = \left\{ \frac{2(n-2)}{\sigma e^A} \right\}^{1/2} e^t,$$

(6) is reduced to the autonomous ordinary differential equation

$$(8) \quad \begin{cases} \ddot{w} + (n-2)\dot{w} + 2(n-2)(e^w - 1) = 0 \\ \lim_{t \rightarrow -\infty} (w(t) - 2t) = \lim_{t \rightarrow -\infty} e^{-t}(\dot{w}(t) - 2) = 0. \end{cases}$$

Then there exists a unique global solution  $w = w(t)$  of (8) by [14]. The orbit  $\mathcal{O} = \{(w(t), z(t)) = (w(t), \dot{w}(t)) \mid t \in \mathbf{R}\}$  starts at  $t = -\infty$  along and below the line  $z = 2$  with  $w = -\infty$ , and approaches the origin  $(0, 0)$  as  $t \rightarrow +\infty$ . If  $3 \leq n \leq 9$ , it proceeds clockwise in  $\{(w, z) \mid z < 2\}$ , crosses infinitely many times  $z$ - and  $w$ -axes alternately, while it keeps to stay in  $\{(w, z) \mid w < 0, 0 < z < 2\}$  in the case of  $n \geq 10$ . Through the Emden transformation (8), the boundary condition in (1) is converted to

$$w(\tau) - 2\tau + A = 0$$

with

$$\left\{ \frac{2(n-2)}{\sigma e^A} \right\}^{1/2} e^\tau = 1.$$

Therefore for any  $\tau \in \mathbf{R}$ ,  $(\sigma_\tau, v_\tau)$  defined by

$$(9) \quad v_\tau(r) = w(t) - 2t - \{w(\tau) - 2\tau\} = w(\log r + \tau) - w(\tau) - 2 \log r$$

with  $r = e^{t-\tau}$  and

$$(10) \quad A_\tau = 2\tau - w(\tau), \quad \sigma_\tau = 2(n-2)e^{2\tau-A} = 2(n-2)e^{w(\tau)}$$

satisfies (1). Conversely every solution of (1) can be expressed in the form of (9) and (10). According to [7], [13], [14] and [19], the total set  $\mathcal{S}$  of the solution  $(\sigma, v)$

of (2) is homeomorphic to  $\mathcal{O}$  through the relation (7) with the constants  $A, \sigma$  determined by (10). This means that  $\mathcal{C}$  is homeomorphic to  $\mathcal{O}$ . Each point of  $\mathcal{O}$  is given as  $(w(\tau), z(\tau))$  and hence  $\mathcal{C}$  is parametrized by  $\tau \in \mathbf{R}$ . Then we have

$$(11) \quad \lambda(\tau) = \omega_n^p 2^{1-p} (n-2)^{1-p} (2-z(\tau))^p e^{(1-p)w(\tau)}.$$

In fact, putting  $K = \int_B e^v dx$  we have

$$\lambda^{1/p-1} = K^{1-p} \sigma^{1/p-1}$$

because of  $K^p = \lambda/\sigma$ . By (10), we have

$$\lambda^{1/p} = \lambda K^{1-p} \sigma^{1/p-1} = \lambda K^{1-p} (2(n-2)e^{w(\tau)})^{1/p-1}.$$

Integrating (1) over  $B$ , we have

$$-\lambda K^{1-p} = \omega_n v'(1) = \omega_n (z(\tau) - 2)$$

by (7). Finally combining two equations, we have the desired one. Hence the behaviour of  $\tau = +\infty$  follows at once. On the other hand  $\lim_{\tau \rightarrow -\infty} v_\tau = 0$  and  $\lim_{\tau \rightarrow -\infty} \sigma_\tau = 0$  imply that  $\lim_{\tau \rightarrow -\infty} \lambda_\tau = 0$ . The rest follows.

In the use of (11), we have

$$\dot{\lambda}(\tau) = \omega_n^p 2^{1-p} (n-2)^{1-p} (2-z(\tau))^{p-1} e^{(1-p)w(\tau)} \{(1-p)z(\tau)(2-z(\tau)) - pz(\tau)\}.$$

If  $n \geq 10$  and  $p \leq 1$ ,  $\dot{\lambda}(\tau) > 0$  for all  $\tau \in \mathbf{R}$  which proves the statement. Next we concentrate on the case of  $3 \leq n \leq 9$  and  $2/n \leq p \leq 1$ . To do so, we put  $g(\tau) = (1-p)z(\tau)(2-z(\tau)) - pz(\tau)$ . Then we have

$$(12) \quad \dot{g}(\tau) = 2(1-p)(2-z(\tau))\dot{z}(\tau) + (pn-2)\dot{z}(\tau) + 2p(n-2)e^{w(\tau)}z(\tau).$$

Let  $\mathcal{O}_k$  ( $k \geq 2$ ) denote the successive points  $w$ -axis and  $z = -2(e^w - 1)$  crossed by the orbit  $\mathcal{O}$  in  $w-z$  plane in order. Moreover we set  $\mathcal{O}_1 = (-\infty, 2)$ . Then  $\dot{\lambda}(\tau) > 0$  and  $\dot{\lambda}(\tau) < 0$  on the arc  $\mathcal{O}_{4k-3}\mathcal{O}_{4k-2}$  and  $\mathcal{O}_{4k-1}\mathcal{O}_{4k}$ , respectively for  $k \geq 1$ . On the other hand  $\dot{g}(\tau) < 0$  and  $\dot{g}(\tau) > 0$  on the arc  $\mathcal{O}_{4k-2}\mathcal{O}_{4k-1}$  and  $\mathcal{O}_{4k}\mathcal{O}_{4k+1}$ , respectively for  $k \geq 1$ . Hence there exists a unique point  $\tilde{\mathcal{O}} = (w(\tau), z(\tau))$  on the every arc  $\mathcal{O}_{4k-2}\mathcal{O}_{4k-1}$  and  $\mathcal{O}_{4k}\mathcal{O}_{4k+1}$  respectively for  $k \geq 1$  such that  $\dot{\lambda}(\tau) = 0$ . The proof is complete.  $\square$

We proceed the proof of Theorem 3 in the same argument and computation as in [15].

**Proof of Theorem 3.** As far as we consider negative eigenvalues in (4), the corresponding eigenfunctions are radially symmetric ([12], [13]). Hence  $i(\lambda, v) = i_R(\lambda, v)$  for

$(\lambda, v) \in \mathcal{C}$ . We denote by  $\mu_\tau^l$  the  $l$ -th eigenvalue of (4) in  $(\lambda(\tau), v(\tau)) \in \mathcal{C}$  corresponding to radially symmetric eigenfunctions. Any of them is simple. If  $(\lambda(\tau), v(\tau)) \in \mathcal{C}$  is the turning point of  $\mathcal{C}$ , then there exists  $l \geq 1$  such that  $\mu_\tau^l = 0$  by the implicit function theorem. On the contrary,  $\mu_\tau^l = 0$  for some  $l \geq 1$  at  $(\lambda(\tau), v(\tau)) \in \mathcal{C}$  implies that it is a turning point by the bifurcation theorem from the critical point of odd multiplicity ([17], [18]). Since  $\mu_{-\tau}^1 > 0$  for sufficiently large  $\tau > 0$ , we have  $i(\lambda, v) = 0$  for  $(\lambda(\tau), v(\tau)) \in \mathcal{C}$  for  $n \geq 10$  and  $p \leq 1$ .

For  $3 \leq n \leq 9$  and  $2/n \leq p \leq 1$ , let  $T_k = (\lambda(\tau_k), v(\tau_k))$  for  $\tau_1 < \tau_2 < \dots$  denote the turning point of  $\mathcal{C}$ . Then we have  $\mu_{\tau_k}^l = 0$  for some  $l \geq 1$  and we have only to show that  $\dot{\mu}_{\tau=\tau_k}^l < 0$  for all  $k \geq 1$ .

Differentiating (1) with respect to  $\tau$ , we have

$$(13) \quad \begin{cases} \Delta \dot{v} + \dot{\lambda} \frac{e^v}{(\int_B e^v dx)^p} + \lambda \frac{e^v \dot{v}}{(\int_B e^v dx)^p} - \lambda p \frac{\int_B e^v \dot{v} dx}{(\int_B e^v dx)^{p+1}} e^v = 0 & x \in B, \\ \dot{v} = 0 & x \in \partial B, \end{cases}$$

and hence

$$\begin{cases} \Delta \dot{v}_k + \dot{\lambda}_k \frac{e^{v_k} \dot{v}_k}{(\int_B e^{v_k} dx)^p} - \lambda_k p \frac{\int_B e^{v_k} \dot{v}_k dx}{(\int_B e^{v_k} dx)^{p+1}} e^{v_k} = 0 & x \in B, \\ \dot{v}_k = 0 & x \in \partial B, \end{cases}$$

for  $v_k = v(\cdot, \tau_k)$ . Then we have

$$\dot{v}(r, \tau) = \dot{w}(\log r + \tau) - \dot{w}(\tau) \neq 0,$$

and therefore,  $\dot{v}_k$  is an eigenfunction of (4) corresponding to  $\mu = \mu_{\tau_k}^l = 0$ . Then, the standard perturbation theory ([8]) guarantees the existence of  $\phi = \phi(\cdot, \tau)$  and  $\mu = \mu(\tau)$  satisfying (4),  $\phi(\cdot, \tau_k) = \dot{v}_k$ , and  $\mu(\tau_k) = \mu_{\tau_k}^k = 0$ . Differentiating (4) and (13) with respect to  $\tau$ , subtracting each other with  $\tau = \tau_k$ , multiplying by  $\dot{v}$  and integrating it over  $B$ , we have

$$(14) \quad \ddot{\lambda} \frac{\int_B e^v \dot{v} dx}{(\int_B e^v dx)^p} = \dot{\mu} \int_B \dot{v}^2 dx,$$

where  $\dot{\mu} = \dot{\mu}(\tau_k)$ ,  $\ddot{\lambda} = \ddot{\lambda}(\tau_k)$ ,  $v = v(\cdot, \tau_k)$  and  $\dot{v} = \dot{v}(\cdot, \tau_k)$ . As is stated in the proof of Theorem 2,

$$(15) \quad \dot{\lambda}(\tau) = \omega_n^p 2^{1-p} (2 - z(\tau))^{p-1} e^{(1-p)w(\tau)} g(\tau) = \omega_n^p \sigma(\tau)^{1-p} (2 - z(\tau))^{p-1} g(\tau),$$

where  $g(\tau) = (1 - p)z(\tau)(2 - z(\tau)) - pz(\tau)$ . Differentiating (15), we have

$$\begin{aligned} \ddot{\lambda}(\tau) &= \omega_n^p \sigma(\tau)^{1-p} (2 - z(\tau))^{p-2} \\ &\quad \times \{(1 - p)(-z(\tau)^2 + 2z(\tau) + \dot{z}(\tau))g(\tau) + (2 - z(\tau))\dot{g}(\tau)\}. \end{aligned}$$

Since  $\dot{\lambda}(\tau) = 0$  at  $\tau = \tau_k$ , it holds that  $g(\tau_k) = 0$ , namely,

$$(1-p)z(\tau_k)(2-z(\tau_k)) - p\dot{z}(\tau_k) = 0$$

from (15). We have  $\lambda = \omega_n^p(2-z)^p\sigma^{1-p}$  from (10) and (11),

$$\dot{\lambda} = \omega_n^p(2-z)^{p-1}\sigma^{1-p}\{(2-z)\sigma^{-1}\dot{\sigma}(1-p) - pz\}$$

by differentiating, and

$$\sigma\dot{z} = \frac{(2-z)(1-p)}{p}\dot{\sigma}$$

at  $\tau = \tau_k$ . Hence we have

$$(16) \quad \begin{aligned} \ddot{\lambda}(\tau) &= \omega_n^p\sigma(\tau)^{-p}(2-z(\tau))^{p-1} \\ &\times \left\{ \frac{2(p-1)^2(2-z(\tau))^2}{p} + \frac{(pn-2)(1-p)(2-z(\tau))}{p} + p\sigma(\tau) \right\} \dot{\sigma}(\tau) \end{aligned}$$

at  $\tau = \tau_k$ . Since  $\sigma = \lambda / (\int_B e^v dx)^p$ , it holds that

$$(17) \quad \dot{\sigma} = \frac{\dot{\lambda}}{(\int_B e^v dx)^p} - \lambda p \frac{\int_B e^v \dot{v} dx}{(\int_B e^v dx)^{p+1}}.$$

Hence we have

$$\begin{aligned} \dot{\mu} \int_B \dot{v}^2 dx &= -\lambda p \frac{(\int_B e^v \dot{v} dx)^2}{(\int_B e^v dx)^{2p+1}} \omega_n^p \sigma(\tau)^{-p} (2-z(\tau))^{p-1} \\ &\times \left\{ \frac{2(p-1)^2(2-z(\tau))^2}{p} + \frac{(pn-2)(1-p)(2-z(\tau))}{p} + p\sigma(\tau) \right\} \end{aligned}$$

at  $\tau = \tau_k$  from (14) and (16). Let  $\int_B e^v \dot{v} dx = 0$ . Then (17) means that  $\dot{\lambda} = 0$  and  $\dot{\sigma} = 0$  vanish at  $\tau = \tau_k$  simultaneously. However it is impossible from (10) and (11). Finally we have  $\dot{\mu} < 0$  at  $\tau = \tau_k$ .  $\square$

#### 4. Annulus domain

In this section, we assume that  $\Omega = A_a = \{x \in \mathbf{R}^n \mid a < |x| < 1\}$  with  $a \in (0, 1)$  and consider radially symmetric solutions of (1).

We cite known results in [14] for the case of  $3 \leq n \leq 9$ . Radial solutions of (2) satisfy

$$(18) \quad \begin{cases} (r^{n-1}v')' + \sigma r^{n-1}e^v = 0 & \text{for } a < r < 1, \\ v(a) = v(1) = 0. \end{cases}$$

We can continue the solution of (18) up to  $r = +0$  satisfying

$$\lim_{r \downarrow 0} \left( v'(r) - \frac{L}{r^{n-1}} \right) = 0$$

and

$$\lim_{r \downarrow 0} \left( v(r) + \frac{L}{(n-2)r^{n-2}} \right) = M$$

for some  $L, M > 0$ . Through the modified Emden transformation

$$v(r) = w(t) - 2t + M, \quad r = Be^t$$

with

$$B = \left\{ \frac{2(n-2)}{\sigma e^M} \right\}^{1/2},$$

(18) is reduced to the autonomous ordinary differential equation

$$(19) \quad \begin{cases} \ddot{w} + (n-2)\dot{w} + 2(n-2)(e^w - 1) = 0, \\ \lim_{t \rightarrow -\infty} (w(t) - 2t + \alpha e^{-(n-2)t}) = 0, \\ \lim_{t \rightarrow -\infty} e^{-t} (\dot{w}(t) - 2 - \alpha(n-2)e^{-(n-2)t}) = 0, \end{cases}$$

where  $\alpha = LB^{-(n-2)}/(n-2)$ .

Then there exists a unique global solution  $w = w_\alpha(t)$  of (19) for every  $\alpha > 0$ . The orbit  $\mathcal{O}_\alpha = \{(w(t), z(t)) = (w_\alpha(t), \dot{w}_\alpha(t)) \mid t \in \mathbf{R}\}$  starts at  $t = -\infty$  above the line  $z = 2$  with  $w = -\infty$ , and approaches the origin  $(0, 0)$  as  $t \rightarrow +\infty$ . Then the family of orbits  $\{\mathcal{O}_\alpha\}_{\alpha \geq 0}$  forms a foliation, that is,  $\mathcal{O}_\alpha \cap \mathcal{O}_\beta = \emptyset$  if  $\alpha \neq \beta$ . Every orbit  $\mathcal{O}_\alpha$  with  $\alpha > 0$  crosses the line  $z = 2$  just once. Moreover if  $\alpha > \beta > 0$ ,  $\mathcal{O}_\alpha$  lies in the left side of  $\mathcal{O}_\beta$  as  $t$  increases. For every point  $(\eta_0, \zeta_0) \in \Sigma \equiv \{(w, z) \mid w > 0, z > 2\}$ , there is a unique nonnegative  $\alpha$  such that  $(\eta_0, \zeta_0) \in \mathcal{O}_\alpha$ .

Let  $w_\alpha = w_\alpha(t)$  be the solution of (19). Then every point  $(w, z)$  on  $\mathcal{O}_\alpha$  determines the time  $t$  so that the boundary condition

$$v(a) = v(1) = 0$$

is converted into

$$(20) \quad w(t^-) - 2t^- + M = 0 \quad \text{and} \quad w(t^+) - 2t^+ + M = 0$$

for

$$(21) \quad t^- = \log \frac{a}{B} \quad \text{and} \quad t^+ = \log \frac{1}{B}.$$

Henceforward we omit the subscripts  $\alpha$  and so on unless there is any confusion. Further  $w^\pm$  and  $z^\pm$  denote  $w_\alpha(t^\pm)$  and  $z_\alpha(t^\pm)$ , respectively. From (20) and (21), it holds that

$$(22) \quad w^+ - w^- = -2 \log a \quad \text{and} \quad t^+ - t^- = -\log a.$$

Conversely, if there exists a pair of points  $P^\pm(w^\pm, z^\pm)$  on  $O_\alpha$  satisfying (22), we have a radial solution  $v = v(r)$  for (18) with some positive constant  $\sigma$ . In fact, we define  $B, M, L, \sigma$  and  $v(r)$  as

$$\begin{aligned} B &= ae^{-t^-} = e^{-t^+}, \\ M &= 2t^- - w^- = 2t^+ - w^+, \\ L &= \alpha(n-2)a^{n-2}e^{-(n-2)t^-} = \alpha(n-2)e^{-(n-2)t^+}, \\ \sigma &= 2(n-2)e^{w^+} = 2(n-2)a^{-2}e^{w^-}, \\ v(r) &= w(t) - 2t + M. \end{aligned}$$

Therefore the structure of the solution of (18) is reduced to that of pairs of  $P^\pm$  on  $O_\alpha$ . We call  $\{P_\alpha^\pm\}$  the boundary pair on  $O_\alpha$  associated with the annulus  $A_\alpha$  in  $\mathbf{R}^n$ .

For every  $a \in (0, 1)$ , there exists a unique pair of points  $P_\alpha^\pm = P^\pm(w_\alpha^\pm, z_\alpha^\pm)$  satisfying (22) on each orbit  $O_\alpha$  ( $\alpha > 0$ ). The points  $P_\alpha^+$  and  $P_\alpha^-$  lie below and above the line  $z = 2$  respectively, and further, these points  $P_\alpha^\pm$  depend on  $\alpha$  continuously. Conversely for each point  $(w_\alpha(t), z_\alpha(t))$  on  $O_\alpha$  with  $z_\alpha(t) < 2$  ( $z_\alpha(t) > 2$ ), there exists a unique  $a^* = a^*(t) \in (0, 1)$  ( $a_* = a_*(t) \in (0, 1)$ ) such that

$$\begin{aligned} w_\alpha(t) - w_\alpha(t + \log a^*) &= -2 \log a^*, \\ (w_\alpha(t - \log a_*) - w_\alpha(t)) &= -2 \log a_*. \end{aligned}$$

Hence we have only to study  $K_\alpha = \{P_\alpha^-(w_\alpha^-, z_\alpha^-) \mid \alpha > 0\}$  to get the structure of solution of (18). The set  $K_\alpha$  forms a continuous curve in  $\mathbf{R}^2$ , which is homeomorphic to  $\mathbf{R}$ . Now we have two lemmas.

**Lemma 1** ([14], Lemmas 4.8 and 4.9). *For any fixed  $a \in (0, 1)$ , we have*

$$\begin{cases} \lim_{\alpha \rightarrow 0} w_\alpha^+ = \lim_{\alpha \rightarrow 0} w_\alpha^- = -\infty, \\ \lim_{\alpha \rightarrow 0} z_\alpha^+ = \lim_{\alpha \rightarrow 0} z_\alpha^- = 2, \\ \lim_{\alpha \rightarrow 0} \|v_\alpha\|_{C(\overline{\Omega})} = 0, \\ \lim_{\alpha \rightarrow 0} \sigma_\alpha = 0, \end{cases}$$

where  $P_\alpha^\pm(w_\alpha^\pm, z_\alpha^\pm)$  are the pair of boundary points on  $O_\alpha$  associated with  $A_\alpha$ .

**Lemma 2** ([14], Remark 6.1). *For any fixed  $a \in (0, 1)$ , we have*

$$\begin{cases} \lim_{\alpha \rightarrow +\infty} w_\alpha^+ = \lim_{\alpha \rightarrow +\infty} w_\alpha^- = -\infty, \\ \lim_{\alpha \rightarrow +\infty} \|v_\alpha\|_{C(\overline{A_a})} = +\infty, \\ \lim_{\alpha \rightarrow +\infty} \sigma_\alpha = 0, \end{cases}$$

where  $P_\alpha^\pm(w_\alpha^\pm, z_\alpha^\pm)$  are the pair of boundary points on  $O_\alpha$  associated with  $A_a$ .

Proof of Theorem 4. We concentrate on the case of  $3 \leq n \leq 9$ . The behaviour of  $\lambda_\alpha$  and  $v_\alpha$  as  $\alpha \rightarrow +0$  follows from Lemma 1, that is,

$$\lim_{\alpha \rightarrow 0} \lambda_\alpha = \lim_{\alpha \rightarrow 0} \sigma_\alpha \left( \int_{A_a} e^{v_\alpha} dx \right)^p = 0 \quad \text{and} \quad \lim_{\alpha \rightarrow 0} \|v_\alpha\|_{C(\overline{A_a})} = 0.$$

Integrating (1) over  $A_a$ , we have

$$(23) \quad \lambda = \omega_n^p (2(n-2)a^{-2}e^{w^-})^{1-p} (a^{n-2}(z^- - 2) - (z^+ - 2))$$

in the same way as we deduce (11). Now it holds that

$$\lim_{\alpha \rightarrow +\infty} z_\alpha^- = +\infty.$$

In fact, we assume that  $z_\alpha^- < M$  for any  $\alpha \in \mathbf{R}$ , where  $M$  is a positive constant owing to  $z_\alpha^- > 2$ . Since  $K_a$  is homeomorphic to  $\mathbf{R}$ , there is a constant  $K > 0$  such that  $w_\alpha^- < K$  for any  $\alpha \in \mathbf{R}$  by Lemma 2. By setting  $Q = \{(w, z) \mid w < K, 2 < z < M\}$ , it holds that  $O_\beta \cap \overline{Q} = \emptyset$  for some  $\beta \in \mathbf{R}$ . Actually, for  $(w, z) = (K, M)$  there is  $\alpha > 0$  such that  $(K, M) \in O_\alpha$  because of  $(K, M) \in \Sigma \subset \bigcup_{\alpha \geq 0} O_\alpha$ . Since  $z$ - and  $w$ -coordinate are decreasing and increasing with respect to  $t$  respectively in  $\{(w, z) \mid w > g(z), z > 0\}$ ,  $O_\alpha \cap \overline{Q} = (K, M)$ . Hence if we put  $\beta = \alpha + \delta$  for any  $\delta > 0$ , we have  $O_\beta \cap \overline{Q} = \emptyset$  because  $O_\beta$  lies in the left side of  $O_\alpha$  as  $t$  increases. The points  $\{(w_\alpha^-, z_\alpha^-) \mid \alpha > 0\}$  on  $O_\beta$  don't satisfy  $w_\alpha^- < K, z_\alpha^- < M$  simultaneously, which is a contradiction. Therefore we have  $\lim_{\alpha \rightarrow +\infty} z_\alpha^- = +\infty$ . Finally from (23),  $z_\alpha^+ < 2$  and  $p \geq 1$ , we have  $\lim_{\alpha \rightarrow +\infty} \lambda_\alpha = +\infty$ . In the case of  $n \geq 10$ , we change  $\Sigma$  and  $K$  by  $\{(w, z) \mid w < 0, z > 2\}$  and a negative constant, respectively.  $\square$

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**References**

[1] J. Bebernes and D. Eberly: *Mathematical Problems from Combustion Theory*, Springer-Verlag, New York, 1989.

- [2] J.W. Bebernes and A.A. Lacey: *Global existence and finite-time blow-up for a class of nonlocal parabolic problems*, Adv. Differential Equations **2** (1997), 927–953.
- [3] J.W. Bebernes and P. Talaga: *Nonlocal problems modelling shear banding*, Comm. Appl. Non-linear Anal. **3** (1996), 836–844.
- [4] S. Chandrasekhar: *An Introduction to the Study of Stellar Structure*, Dover, New York, 1957.
- [5] I.M. Gel'fand: *Some problems in the theory of quasilinear equations*, Amer. Math. Soc. Transl. (2) **29** (1963), 295–381.
- [6] B. Gidas, W.-M. Ni and L. Nirenberg: *Symmetry and related properties via the maximal principle*, Comm. Math. Phys. **68** (1979), 209–243.
- [7] D.D. Joseph and T.S. Lundgren: *Quasilinear Dirichlet problems driven by positive sources*, Arch. Rational Mech. Anal. **49** (1972/73), 241–269.
- [8] T. Kato: *Perturbation Theory for Linear Operators*, Springer, Berlin, 1966.
- [9] E.F. Keller and L.A. Segel: *Initiation of slime mold aggregation viewed as an instability*, J. Theor. Biol. **36** (1970), 399–415.
- [10] A.A. Lacey: *Thermal runaway in a nonlocal problem modelling Ohmic heating: Part I*, European J. Appl. Math. **6** (1995), 127–144.
- [11] A.A. Lacey: *Thermal runaway in a nonlocal problem modelling Ohmic heating: Part II*, European J. Appl. Math. **6** (1995), 201–224.
- [12] C.-S. Lin and W.-M. Ni: *A counterexample to the nodal domain conjecture and a related semilinear equation*, Proc. Amer. Math. Soc. **102** (1988), 271–277.
- [13] T. Miyasita and T. Suzuki: *Non-local Gel'fand problem in higher dimension*; in Nonlocal Elliptic and Parabolic Problems, Banach Center Publ. **66** Polish Acad. Sci., Warsaw, 2004, 221–235.
- [14] K. Nagasaki and T. Suzuki: *Radial solutions for  $\Delta u + \lambda e^u = 0$  on annuli in higher dimensions*, J. Differential Equations **100** (1992), 137–161.
- [15] K. Nagasaki and T. Suzuki: *Spectral and related properties about the Emden-Fowler equation  $-\Delta u = \lambda e^u$  on circular domains*, Math. Ann. **299** (1994), 1–15.
- [16] S.I. Pohozaev: *Eigenfunctions of the equation  $\Delta u + \lambda f(u) = 0$* , Soviet Math. Dokl. **6** (1965), 1408–1411.
- [17] P.H. Rabinowitz: *Some global results for nonlinear eigenvalue problems*, J. Functional Analysis **7** (1971), 487–513.
- [18] P.H. Rabinowitz: *Some aspects of nonlinear eigenvalue problems*, Rocky Mountain J. Math. **3** (1973), 161–202.
- [19] T. Suzuki: *Semilinear Elliptic Equations*, Gakkōtoshō Co., Ltd., Tokyo, 1994.
- [20] G. Wolansky: *A critical parabolic estimate and application to nonlocal equations arising in chemotaxis*, Appl. Anal. **66** (1997), 291–321.

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