(29) GCOE Event The 2ndGCOE International Symposium on "Weaving Science Web beyond Particle-Matter Hierarchy" 9:00 - 18:3, February 19, 2010 Date & Time: February 18, 2010 9:30-18:00 -Place: Faculty of Science Aobayama Campus of Tohoku University Invited Speakers: Dr. Indranil Mazumdar (Tata Institute of Fundamental Research, India) **Prof. Philippe Marcq** (PhysicoChimie Curie Institut Curie, France) **Prof. Young-Woo Son** (Korea Institute for Advanced Study, Korea) **Dr.Jean Coupon** (Institut d'Astrophysique de Paris, France) Prof. Ryushi Goto (Osaka University, Japan) Prof. Jenann Ismael (University of Arisona, USA and University of Sydney, Australia) **Prof. Jie Meng** (Beihang University, China) **Prof. Huaizhe Xu** (Beihang University, China) Assistant Prof. Hojun IM (Hirosaki University) M.S. Masahiro FUTAKI (The University of Tokyo) Dr. Masahiko IGASHIRA (Osaka University)

Organizing Committee:

Prof. Hideo Kozono (Mathematics, Tohoku Univ.): Chairman Assistant Prof. Katsuhiko Sato (Condensed Matter physics, Tohoku Univ.) Assistant Prof. Yousuke Itoh (Astrophysics, Tohoku Univ.) Assistant Prof. Yoshihiro Ueda (Mathematics, Tohoku Univ.) Assistant Prof. Yuichi Nohara (Mathematics, Tohoku Univ.) Assistant Prof. Masaki Asano (Particle Theory, Tohoku Univ.) Assistant Prof. Yoshiyuki Onuki (Particle Experiment, Tohoku Univ.) Assistant Prof. Takeshi Koike (Nuclear Experiment, Tohoku Univ.) Assistant Prof. Hidekatsu Nemura (Nuclear Theory, Tohoku Univ.) Assistant Prof. Jin Sung Park (Condensed Matter Physics, Tohoku Univ.) Assistant Prof. Tatsuro Yuge (Condensed Matter Physics, Tohoku Univ.) Assistant Prof. Tomohiro Yoshikawa (Astrophysics, Tohoku Univ.) Assistant Prof. Shota Sato (Mathematics, Tohoku Univ.) Assistant Prof. Yohei Matsuda (Nuclear Experiment, Tohoku Univ.) Assistant Prof. Masafumi Kurachi (Particle Theory, Tohoku Univ.) Dr. Masaru Yonehara (Philosophy, Tohoku Univ.)

Tohoku University GCOE program "Weaving Science Web beyond Particle-Matter Hierarchy"

Outline:

The aim of this Symposium is to explore new science frontiers through "Weaving Science Web beyond Particle-Matter hierarchy" among physics - astrophysics - mathematics - philosophy.

Number of participants:

Symposium:

Banquet:

ł	Tobeku Univer The 2 Inter "Weaving S February " Main Lecture Aobayama Ca	And GCOE to be a consecutive active to the second s
	Brogram & Invited	eneskare
	Program & Invited	speakers
	Plenary session 0930 Indranil Mazumdar Philippe Marcq Young-Woo Son Jean Coupon Ryushi Goto	1220 Feb. 18 9 Main Lecture Hall (Tata Institute of Fundamental Research, India) (PhysicoChimile Ourie Institut Curle, France) (Korea Institute for Advanced Study, Korea) (Institut d'Astrophysique de Paris, France) (Osaka University, Japan)
	Jenann Ismael	(University of Ansona, USA and University of Sydney, Australia)
-	Special Session 1690- Jie Meng Huaizhe Xu	1718 Feb.18 R203 Sogatoh Bidg. (Beihang University, China) (Belhang University, China)
	Parallel Session 1350- Hojun IM Masahiro FUTAKI Masahiko IGASHIRA	1936 Feb. 18, 1350-1730 Feb. 19 R203-R205 Segetoh Bidg. (Hirosaki University, Japan) (The University of Tokyo, Japan) (Osaka University, Japan)
	Collaborative Research	Section 1730-1908 Erb 19 8283-8305 Eventsh Bide
	Poster Session 1710-	1830 Feb.18 Entrance-Hall Sugotoh Bidg.
	Banquet 1820-	2000 Feb.19 Campus Cafeteria
	Organizing Commit	tee
	Hideo Kozono (Chairman Katsuhiko Sato, Yousuke J Yoshiyuld Onuki, Takeshi K Tomohiro Yoshikawa, Shot	1) toh, Yoshihiro Ueda, Yuichi Nohara, Masaki Asano, Kolka, Hidekatsu Nemura, Jin Sung Park, Tatsuro Yuge, Ja Sato, Yohoi Matsuda, Masafumi Kurachi
623	Inguirles: GCOE Science Succord Division Reasonatic Genera for New no Buienne Tel E-mel: gcost (Jacienceset), toheku, au jp	ndau U-incraity, 8-3 Aran avi-Ace-Aube, Aube-ku: Scytter 950-5578, Japa-

Program:

Februar	y 18 Thursday
09:30 -	"Opening address" Kunio INOUE (GCOE program leader)
Plenary S	Session Chairmen H. Tamura / T. Kawakatsu / Y. Hirayama
09:30-10:20	"Halo World: The Story according to Faddeev, Efimov and Fano" I. MAZUMDAR
	(Chairman H. Tamura)
10:30-11:20	"Mechanics of Stress Fibers" P. MARCQ (<i>Chairman T. Kawakatsu</i>)
11:30-12:20	"Relativistic Dirac Electrons in Condensed Matters – Graphene and others" Y. W. SON (Chairman Y. Hirayama)
12:30-13:50	Lunch
Parallel S	Session A - Particle/ Nuclear/ Astro-physics Group Chairman K. Hikasa
13:50-14:20	"General WIMP search at International Linear Collider" M. ASANO
14:30-15:00	"Primordial nucleosynthesis and recent topics in particle cosmology" K. KOHRI
15:10-15:40	"Higgsless Models for the Electroweak Symmetry Breaking" M. KURACHI
15:50-16:00	Coffee Break
Parallel S	Session B – Condensed Matter Physics Group / Nuclear Physics Chairman R. Saito
13:50-14:20	"Exciton effect of Raman resonance window of single wall carbon nanotubes" J. PARK
14:30-15:00	"Properties of response function of nonequilibrium steady state" T. YUGE
15:10-15:40	"Hyperon-nucleon interactions from lattice QCD" H.NEMURA
15:50-16:00	Coffee Break
Parallel S	Session C – Mathematics and Philosophy GroupChairmanH. Kozono
13:50-14:20	"Application of the anti-derivative method to the half space problem for damped
	wave equation with non-convexity" Y.UEDA
14:30-15:00	"Toric degeneration of Gelfand-Cetlin systems" Y.NOHARA
15:10-15:40	"Singular Backward Self-Similar Solutions of a Semilinear Parabolic Equation" S. SATO
15:50-16:00	Coffee Break
Special Second	ession Chairman O. Hashimoto
16:00-16:30	"Covariant density functional theory for Nuclear structure and application in astrophysics"
	J.Meng
16:35-17:05	"Control of Magnetic Property and Magnetic Coupling Mechanism in ZnO based DMS"
	H.Xu
Poster Se	ssion
17:10-18:30	80 Research assistants present their research.

February	19 Friday
Plenary Se	ssion Chairmen K. Naoe / R. Miyaoka / T. Futamase
09:30-10:20	"A Fresh Direction for Quantum Research: What Might Entanglement Be Telling Us?" J.ISMAEL (<i>Chairman K. Naoe</i>)
10:30-11:20	"Calabi-Yau structures and Einstein-Sasaki structures" R. GOTO(Chairman R. Miyaoka)
11:30-12:20	"Probing the large scale structure in the Universe with CFHTLS"
	J. COUPON (Chairman T. Futamase)
12:30-13:50	Lunch
Parallel Se	ssion A – Particle/ Nuclear/ Astro-physics Group <i>Chairmen H. Tamura / T. Futamase</i>
13:50-14:20	"Measurements of Cabibbo-Kobayashi-Maskawa unitary triangle angle phi3 at Belle experiment"
	Y. ONUKI (Chairman H. Tamura)
14:30-15:00	"Hypernuclear γ-ray spectroscopy at J-PARC" T. KOIKE (<i>Chairman H. Tamura</i>)
15:10-15:40	"Search for the alpha cluster condensed state in ¹⁶ 0" M. ITOH (<i>Chairman H. Tamura</i>)
15:50-16:10	Coffee Break
16:10-16:40	"A dipole anisotropy of galaxy distribution: Does the CMB rest frame exist in the
	local universe?" Y. ITOH (Chairman T. Futamase)
16:50-17:20	"MOIRCS Deep Survey: Near-Infrared Observations of Galaxies at the Star-Forming
	Epoch of the Universe" T. YOSHIKAWA (<i>Chairman T. Futamase</i>)
Parallel Se	ssion B – Condensed Matter Physics Group Chairmen Y. Kuramoto / T. Takahashi
13:50-14:20	"From graphene to Z2 topological insulator" K. IMURA (<i>Chairman Y. Kuramoto</i>)
14:30-15:00	"Shear-induced phase separation of complex fluids" K. SATO (<i>Chairman Y. Kuramoto</i>)
15:10-15:40	"Ce 4f electronic structure of heavy-fermion systems across quantum critical point:
1.5.50.1.5.1.5	a resonant angle-resolved photoemission study" H. IM (<i>Chairman Y. Kuramoto</i>)
15:50-16:10	Coffee Break
16:10-16:40	"Spin-resolved ultrahigh-resolution ARPES study of Rashba effect on semi-metal surface" S. SOUMA (<i>Chairman T. Takahashi</i>)
16:50-17:20	"Metal-contact effect on graphene" R. NOUCHI (<i>Chairman T. Takahashi</i>)
Parallel Se	ssion C – Mathematics and Philosophy Group Chairmen T. Shiova / K. Naoe
13:50-14:20	"Around Homological Mirror Symmetry" M. FUTAKI (<i>Chairman T. Shioya</i>)
14:30-15:00	"On the heat equation in a half space with a nonlinear boundary condition"
•	T. KAWAKAMI (Chairman T. Shioya)
15:10-15:40	"An Approach to Experimental Philosophy of Mind" M. IGASHIRA (Chairman K. Naoe)
15:50-16:10	Coffee Break
16:10-16:40	"Utilitarianism and Rawls" Masaru YONEHARA (<i>Chairman K. Naoe</i>)

16:50-17:20	"Husserl and Disjunctivism: On the Theory of Perceptual Experience in
	Transcendental Phenomenology"S. SATO (Chairman K. Naoe)
Collaborat	ive Research Session A Chairman K. Inoue
17:30-17:40	"Research Center for Electron Photon Science (Laboratory of Nuclear Science)
	Development of an electro-magnetic calorimter made up of ceramic Pr:LuAG
	scintillator" T. ISHIKAWA
17:40-17:50	"Study of geo neutrinos with KamLAND" Y. SHIMIZU
17:50-18:00	"Geometric Analysis on Einstein Equation" S. YAMADA
Collaborat	ive Research Session B Chairman H. Yamamoto
17:30-17:40	"Study of elementary particles, nuclei, and high pressure condensed matter in
	gravitational wave astronomy" Y. ITOH
17:45-17:55	"Development of Readout Board for FPCCD Detector" Y. TAKUBO
Collaborat	ive Research Session C Chairman T. Ogawa
17:30-18:40	"Theoretical formulation of morphology dynamics of membranes based on a
	combination of mathematical models and differential geometry" K. SATO
17:45-17:55	"The search for safety and certainty in foundations of mathematics from the
	logical and philosophical point of view" K. YOKOYAMA

Poster Presentations:

РО	Title / Name
no.	
1.	Coherent control of excited states by multipulse photo excitation
	Kenta Abe (Physics, D2, Tohoku Univ.)
2.	de Haas –van Alphen effect studies in the Antiferro-Quadrupolar ordering systems:
	$Pr_{1-x}La_xPb_3$ and $U_{1-x}Th_xPd_3$
	Toshiyuki Isshiki (Physics, D3, Tohoku Univ.)
3.	Measurement of the Superparticle Mass Spectrum in the Long-Lived Stau Scenario at
	the LHC
	Takumi Ito (Physics, D1, Tohoku University)
4.	In situ NMR imaging of lithium- ion batteries during charge/discharge cycle
	Yoshiki Iwai (Physics, D3, Tohoku University)
5.	Crack propagation in largely deformed rubber sheet
	Daiki Endo (Physics, D1, Tohoku University)
6.	Optical response of photonic crystal with multi-layered structure
	Rihei Endo (Physics, D2, Tohoku University)
7.	Pseudo-spin Kondo effect in a capacitively-coupled parallel double quantum dot
	Yuma OKAZAKI (Physics, D1, Tohoku University)
8.	Spin and charge dynamics in a photo-excited double exchange system
	Yu Kanamori (Physics, D2, Tohoku University)
9.	Observation of ⁸ B Solar Neutrinos with KamLAND
	Yoshiaki Kibe (Physics, D3, Tohoku University)
10.	z~3 LYMAN BREAK GALAXY CLUSTER SURVEY IN SSA22 WITH VIMOS
	Katsuki Kousai (Physics, D2, Tohoku University)
11.	EELS and SXES studies of electronic structures of Al-TM alloys
	Shogo Koshiya (Physics, D1, Tohoku University)
12.	Comparison of membrane physical property changes between DMPC membrane and
	DMPE membrane induced by melittin
	Atsuji Kodama (Physics, D3, Tohoku University)
13.	Resonant x-ray scattering experiment on Pr(Ru _{1-x} Rh _x) ₄ P ₁₂
	Kotaro Saito (Physics, D1, Tohoku University)
14.	Study of geo neutrinos with KamLAND
	Yuri Shimizu (Physics, D3, Tohoku University)

15.	Pressure dependence of <i>B-T</i> phase diagram in heavy-fermion superconductor CeRhSi ₃
	Tetsuya Sugawara (Physics, D2 Tohoku University)
16.	Neutral meson photoproduction with electromagnetic calorimeter complex FOREST
	Koutaku Suzuki (Physics, D3, Tohoku University)
17.	Micro-photoluminescence around spin phase transition of v=2/3 fractional quantum
	Hall regime
	Jun-ichiro Hayakawa (Physics, D1, Tohoku University)
18.	Ionic conductivity of Trehalose-Water-Lithium iodide mixture in glass and supercooled
	liquid state
	Reiji Takekawa (Physics, D2, Tohoku University)
19.	Progress of the new DAQ system for wide-band solar neutrino observation with
	KamLAND
	Takemoto Yasuhiro (Physics, D1, Tohoku University)
20.	Double Chooz PMT preparation
	TABATA Hiroshi (Physics, D2, Tohoku University)
21.	Effects of lanthanoid ion on phosphate head groups in DMPC membrane and
	DMPC/DHPC mixtures.
	Kouya Tamatsukuri (Physics, D3, Tohoku University)
22.	Status of double beta decay experiment with KamLAND
	Azusa Terashima (Physics, D1, Tohoku University)
23.	Simulations on Dynamics of Wormlike Micellar System Using Particle-Field Hybrid
	Models
	Masatoshi Toda (Physics, D3, Tohoku University)
24.	Magneto-dielectric phenomena in charge ordered system with frustrated geometrical
	lattice
	Makoto Naka (Physics, D2, Tohoku University)
25.	Pi-mesonic decays of Λ hypernuclei
	Yoji Nakagawa (Physics, D1, Tohoku University)
26.	Axionic Mirage Mediation
	Shuntaro Nakamura (Physics, D3, Tohoku University)
27.	The Study of the Origin of Lyman alpha Emitters with Large Equivalent Widths
	Yuki Nakamura (Physics, D3, Tohoku University)
28.	Ultrafast broadband THz response of photo-induced metallic state in charge ordered

	insulator <i>a</i> -(BEDT-TTF) ₂ I ₃
	Hideki Nakaya (Physics, D3, Tohoku University)
29.	Ring-Exchange interaction in Orbital Degenerate System
	Joji Nasu (Physics, D2, Tohoku University)
30.	Shape Memory Effect Induced by Magnetic-Field Rotation in MnV ₂ O ₄ Spinel
	Compound
	Yoichi Nii (Physics, D1, Tohoku University)
31.	The ω meson photoproduction on the nucleon in the threshold region
	Ryo HASHIMOTO (Physics, D3, Tohoku University)
32.	Prototype fast imaging detector for YN-scattering
22	Ryotaro Honda (Physics, M1, Tohoku University) Ground state phase diagram of graphene in a high Landau level: A density matrix
33.	renormalization group study
	Tatsuya Higashi (Physics, D2, Tohoku University)
34.	Investigation of the $n(\gamma, K^{\theta})$ A reaction near the threshold
	K.Futatsukawa (Physics, D3, Tohoku University)
35.	Traffic flow of two lanes with a bottleneck
	Sho Furuhashi (Physics, D2, Tohoku University)
36.	Staggered Order with Kondo and Crystalline Field Singlets in f ² System
	Shintaro Hoshino (Physics, D1, Tohoku University)
37.	Proton elastic scattering of ⁹ C at 290 MeV
	Yohei Matsuda (Physics, Assistant professor, Tohoku University)
38.	High Resolution and High Statistics Λ Hypernuclear Spectroscopy by the (e,e'K ⁺)
	Reaction
	Akihiko Matsumura (Physics, D3, Tohoku University)
39.	External Gamma-ray Backgrounds of the KamLAND Detector
	Yukie Minekawa (Physics, D2, Tohoku University)
40.	Higgs Triplet Model
	Yusuke Motoki (Physics, D1, Tohoku University)
41.	Electron diffraction study of the low temperature phase of Hollandite-type oxide
	K ₂ Cr ₈ O ₁₆
	Daisuke Morikawa (Physics, D1, Tohoku University)
42.	Light scattering by collective excitation of phonon in quantum paraelectrics

	Ryuta Takano (Physics, D1, Tohoku University)
43.	Formation of the Milky Way Based on the Analysis of Kinematics and Chemical
	Abundance of the Outer Halo Stars
	Miho Ishigaki (Astronomy, D3, Tohoku University)
44.	Searching for Luminous Core-Collapsed Supernovae in a High-z Proto-Cluster
	Nana Morimoto (Astronomy, D1, Tohoku University)
45.	A Study of Light Curves from Rapidly Rotating Neutron Stars
	Kazutoshi Numata (Astronomy, D3, Tohoku University)
46.	CMB Bispectrum from the Second-Order Cosmological Perturbations
	Daisuke Nitta (Astronomy, D3, Tohoku University)
47.	General Properties of Non-Radial Pulsations
	Aprilia (Astronomy, D3, Tohoku University)
48.	A Study on Mathematical Fuzzy Logic
	Ahmad Termimi Bin Ab Ghani (Mathematics, D1, Tohoku University)
49.	Some space-time integrability estimates of the solution for heat equations in two
	dimension
	Norisuke Ioku, (Mathematics, D2, Tohoku University)
50.	Well-posedness for Navier-Stokes equations in modulation spaces with negative
	derivative indices
	Tsukasa Iwabuchi (Mathematics, D2, Tohoku University)
51.	The soul conjecture for Riemannian orbifolds
	Naoki Oishi (Mathematics, D3, Tohoku University)
52.	A new approach to the existence of harmonic maps
	Toshiaki Omori (Mathematics, D2, Tohoku University)
53.	Stability of the interface of a Hele-Shaw flow with two injection points
	Michiaki Onodera (Mathematics, D3, Tohoku University)
54.	The Isomorphism between Motivic Cohomology and K-groups for Equi-Characteristic
	Regular Local Rings
	Yuki Kato (Mathematics, D3, Tohoku University)
55.	Bifurcations in semilinear elliptic equations on thin domains
	Toru Kan (Mathematics, D1, Tohoku University)
56.	On the curvature of the pseudo-volume form defining the carath_eodory measure
	hyperbolicity
	Shin Kikuta (Mathematics, D1, Tohoku University)

57.	Gödel's Incompleteness Theorem, Recursively Axiomatizable Theories, and Medvedev
	Degrees of Unsolvability
	Takayuki Kihara (Mathematics, D1, Tohoku University)
58.	Hypergeometric series on a <i>p</i> -adic field
	Kensaku Kinjo (Mathematics, D2, Tohoku Univercity)
59.	On Hasse principle of purely transcendental extension field in one variable
	Makoto Sakagaito (Mathematics, D3, Tohoku University)
60.	Large time behavior of solutions for system of nonlinear damped wave equations
	Hiroshi Takeda (Mathematics, D3, Tohoku University)
61.	Non-abelian generalization of Iwasawa theory
	Kazuaki Tajima (Mathematics, D1, Tohoku University)
62.	The quadratic subextension of the class field of a real quadratic field
	Toshihide Doi (Mathematics, D2, Tohoku University)
63.	Pattern formation by receptor-based models for regeneration experiments on Hydra
	Madoka Nakayama (Mathematics, D2, Tohoku University)
64.	The computational methods of canonical heights on elliptic curve
	Tadahisa Nara (Mathematics, D3, Tohoku University)
65.	Spatial branching process in random environment
	Nishimori Yasuhito (Mathematics, D1, Tohoku University)
66.	On a periodic decomposition of meromorphic functions
	Takanao Negishi (Mathematics, D2, Tohoku University)
67.	Difficulties of solving problems
	Kojiro Higuchi (Mathematics, D1, Tohoku University)
68.	Undecidability and weak theory of concatenation
	Yoshihiro Horihata (Mathematics, D2, Tohoku University)
69.	Davies' Conjecture for Pseudo-Schrödinger Operators and its Applications to
	Penalization Problem
	Masakuni Matsuura (Mathematics, D1, Tohoku University)
70.	Torsion points of Abelian varieties with values in infinite extension fields
	Yuken Miyasaka (Mathematics, D1, Tohoku University)
71.	Maximum principle for a biological model related to the motion of amoebae
	Harunori Monobe (Mathematics, D2, Tohoku University)
72.	Asymptotic behavior of solutions to the drift-diffusion equation in the whole spaces
	Masakazu Yamamoto (Mathematics, D3, Tohoku University)

73.	A study of the idea of systematic knowledge:
	On the relation between nature and spirit in the organizational view of nature
	Fukuko Abe (Philosophy, D3, Tohoku University)
74.	Medical technology and surrogate decision-making
	Haruka Hikasa (Philosophy, D3, Tohoku University)
75.	Hume's empiricism and the experimental method of reasoning
	Hiromichi Sugawara (Philosophy, D1, Tohoku University)
76.	Epistemic deference and transmission of knowledge:
	How should we (non-scientists) acquire warranted beliefs about scientific propositions?
	Mariko Nihei (Philosophy, D3, Tohoku University)
77.	The Finiteness of Human Beings and the Role of Technology
	Ryozo Suzuki (Philosophy, D3, Tohoku University)
78.	What is ethically problematic in Biogenetics?
	Takuma Obara (Philosophy, D3, Tohoku University)
79.	The ontological genesis of the theoretical attitude
	Tetsurou Yamashita (Philosophy, D3, Tohoku University)
80.	Risk, uncertainty and the precautionary principle: How to deal with scientific
	uncertainty?
	Yasuhiko Fujio (Philosophy, D3, Tohoku University)
81.	The mechanism of suppressed dynamical friction in a constant density core of dwarf
	galaxies
	Shigeki Inoue (Astronomy, D2, Tohoku University)

.....

Coherent control of excited states by multipulse photo excitation

Kenta Abe (Physics, D2, Tohoku Univ.)



Forbidden S₁ State of Carotenoids.



Measurement of the Superparticle Mass Spectrum in the Long-Lived Stau Scenario at the LHC

Takumi Ito (Physics, D1) Particle Theory and Cosmology Group Tohoku University

Reference : TI. R. Kitano and T. Moroi, arXiv:0910.5853[hep-ph] The 2rd GCOE International Symposium 2010.2.18.-19.

Long-lived stau scenario

- * It is naturally realized who two higgs' VEVs). en a model has a large tanß (the ra * Cosmological problems, which charged particle, can be avoid LSP (i.e. gravitino). which may be caused by the long-lived e avoided if there is weekly interacting
- A slow-moving charged track (= stau) informs us a production of superparticles in high energy experiments.
- Further more, we can measure stau's momentum, velocity and energy as well as ones of SM particles.
- -> All final state particles in the event are visible (mont v)
- We have a opportunity to probe the SUSY in detail!



		determined underlyin $M_{W^{\prime\prime\prime}}$ = 430.0 \leftrightarrow 425.9
Endpoint	E-A	Uncertainties
	THE NEW	± 5 GeV
The second second	1	* Stau mass measurement ± 100 MeV



* Loop effects (muon Yukawa) $\sim O(100~{\rm MeV})$ for the large tanß * SUGRA effects $\sim m_{3/2}^2/m_\ell^- ~~ \epsilon m_{3/2} \sim O(1~{\rm GeV})~~ {\rm case}?$

We discuss mass measurements of superpartic	cles
in the long-lived stau scenario at the LHC expe	eriments.
* Neutralino Masses	
by endpoint analysis : $\delta m_{\hat{B}} \sim 1~{ m GeV},\delta m$	$_{W^0} \sim 5 ~{ m GeV}$
charge subtraction method is useful.	
* Selectron & Smuon Masses	
particularly in the case $m_{\hat{H}} > m_{\hat{t}_{H}} > m_{\hat{t}_{H}}$	
as a sharp peak : $\delta m_e \sim \delta m_\mu \sim 1~{ m GeV}$	
• mass difference : $\delta(M_{\tilde{v}} - M_{\tilde{\mu}}) \sim 100$) MeV
* Squark Masses	
as a clear peak : $\delta m_a \sim 10 \text{ GeV}$	16

- Supersymmetry : A famous extension of the Sta odard Model There are several reasons to consider SUSY seriously. * Gauge coupling unification at the very high energy scale (GUT) * Candidates of Dark Matter of the Universe * Predicts a light higgs (prefered by EW preci Naturally solves gauge hierarchy problem → It will appear around TeV scale
 - \leftrightarrow in the energy range of the LHC experiments!
- So, Question is

What is a signature of SUSY at the LHC?

We discuss a determination of the superparticle masses at the LHC experiments in the model with the long-lived stau. * Neutralinos * Sleptons

* Squarks

We perform a Monte Carlo analysis of sparticle production eve the LHC. Our basic ideas and the simulated accuracies of the ma determinations are discussed in the following slides.







(4) No isolated leptons with $p_{\gamma} > 15 \text{ GeV}$



(MSSM-LSP)

Signature of SUSY at the LHC... -> It strongly depends on what the LSP is.

Popular candidates of LSP are: *The lightest neutralino $\tilde{\chi}_1^0$ _____ Large missing \mathbf{p}_{T}
 Time lymma
 Gravitino
 ψµ

 Gravitino
 ψµ
 Useokly interacting (Gravity)

 =
 Gravitino

 =

Monte Carlo Study; Model & Assumptions GMSB (Dire Nelson Shim Monte Carlo Analysis 1 - 10 IeV, $M_{mem} - 900$ IeV V₁ = 3, tan $\beta = 35, ...$ Event Gen. by HERWIG6.510 Fast Detector Sim. by PGS4 1309 67k events are generated (100 fb1) 1230 1180 Asummptions $\frac{\bar{q}_R}{\bar{W}^0}$ 426 * Stau is stable in the detector. 240 194 193 149 \tilde{B} $\begin{array}{c|c} \mbox{Reduced} & \mu \\ \mbox{I GeV} & \\ \mbox{by hands} & \\ \mbox{$\dot{\mu}_R$} \end{array}$









Jet algorism? Cone size? Event selection? Or



In situ NMR imaging of lithium- ion batteries during charge/discharge cycle Institute for Multidisciplinary Research of Advanced Materials, Tohoku University Yoshiki Iwai, Daiki Ohno, Junichi Kawamura



Previous study of quasistatic crack propagation

- An experiment conducted by Yuse and Sano (Nature (London) 362, 329 (1993)) motivates researchers to study about crack pattern especially in brittle materials such as glasses Glasses can be regarded as linear elastic body.
- · As we change temperature difference ΔT or descent speed V. crack pattern changes.



Experimental procedure

- 1. Prepare a cylindrical rod and cover it with very viscous silicon oil.
- 2. As shown in figure, cover the rod with a stretched cylindrical rubber sheet (an inflated balloon) on oil and wait until the rubber sheet is relaxed to be under uniform strain.
- 3. Create a small initial crack at one end of the sheet to create a spontaneous crack. Then observe the crack propagation.

Five typical crack patterns

[Straight crack] Straight crack propagates in a longitudinal direction of the rod

[Oscillating crack] Oscillatory crack

propagates in a longitudinal direction of the rod [Helical crack] Crack keeps a constant angle

to the longitudinal direction of the rod.

[Irregular crack] Crack exhibits several irregular turns and stops after that.

[Stagnating crack] Crack propagates in a very short distance and stops after that



Fracture in large deformation

- · Unlike glasses, many materials exhibit nonlinear elasticity and large dissipation when they fracture under large deformation.
- $dE_{\text{elastic}} = dE_{\text{surface}} + dE_{\text{dissipation}}$
- *dE*_{dissipation} depends on viscosity, viscoelasticity, and plasticity.



Pattern diagram



Pattern diagram 1 (computer simulation)

This diagram is result in nonviscoelastic condition. In large deformation regime, we can see oscillating pattern and straight pattern which may correspond to oscillating crack and straight crack in the experiment.

Conclusion

- We observe five typical crack patterns in the experiment of rupture of rubber sheets.
- · From computer simulation, oscillating pattern and straight pattern appear in large deformation regime in same order of the experiment.
- We conclude that in large deformation regime, very large dissipation arises and causes shortening crack path length per unit length of the rod.

· If we change the amount of silicon oil covering the rod, we can change characteristic relaxation time of the sheet by changing viscous resistance. And if we change diameter of the rod, we can change strain of the sheet.

Note

Crack propagation in largely

deformed rubber sheet

Daiki Endo

Physics D1 Tohoku University

An experiment of rupture of

rubber sheets · Rubber is a typical material which can

deform largely. It shows nonlinear

2.6

2.3

Strain 2.4 +++

Л

 F_c

= 1.0 $\overline{ka_0}$

elasticity and creep.





- characteristic length a₀
- characteristic time $\frac{\gamma_{\text{in}}}{k_{\text{in}}} = T_{in}$
- characteristic force ka₀

Crack patterns (computer simulation)

•As we increase strain or decrease F_c, wave length and amplitude decrease.



Comparing stress-strain

•The model shows nonlinear elasticity and viscoelasticity like a rubber sheet.



Pattern diagram 2 (computer simulation)

This diagram is result of computer simulation which corresponds to pattern diagram in the experiment. Unlike the experiment, we can see no change of the crack pattern by changing characteristic time.

curve





Optical response of photonic crystal with multi-layered structure

Rihei Endo Physics, D2, Tohoku University



Tunneling photons in frustrated total internal reflection





We want to calculate the light pulse propagation of superluminal barrier traversal.

[1] J. Ch. Bose, in Collected Physical Papers of Sir Jagadis Chunder Bose (Longmans, London, 1927), p. 42. [2] F. Goos and H. Ha "nchen, Ann. Phys. (Leipzig) 1, 333 (1947).
 [3] Enders and G.Nimtz, J. Phys. I France 2, 1693-1698 (1992).

II - ii . Refrecrance of Sandwiched layer



II. Theory

II - i . Boundary condition



II - iii . Light tunneling effect of Sandwiched layer

The amount of penetration of light is decided according to the incident angle and the width of the barrier. We calculated the reflectance of light vs barrier width in the case of critical angle. We got the result that reflectance of P-polarization light is smaller than that of S-polarization.



III. Model

We consider about one dimension photonic crystal with the periodic multilayer structure as a model. To enhance "light tunneling effect", we can adjust the refractive indices (n_1, n_2) , the number of layers (N), and the width of the layer (d_1, d_2) . When the angle of incidence is the critical angle, tunneling and propagating are occurred alternately. To examine the behavior of the propagating of the light pulse in this system, we calculate the Maxwell equation.



fig.3: comparison with light propagation of 1D-photonic crystal and that of vacuum

Results

In 1D-PC

In vacuum

At first, we examined the resonance condition of the "light tunneling effect" (Fig1). Next, we examined the relation between reflectivity and the circular frequency when the resonance condition was filled (Fig.2). Finally, we calculated the time evolution of the Gauss pulse that centered on the number of resonance vibrations (Fig.3). In comparison with light propagation of 1D-photonic crystal and that of vacuum, the front velocities quicken more than the light that propagates in the vacuum.

1=-50



What is photonic crystals ?

In generally, photonic crystal (PC) is comprised identical dielectric structures replicated in one-, two-, three-dimensional periodic arrangement as illustrated in Fig.





IV. Conclusion

We have understood the following from calculating the Maxwell equation in one dimension photonic crystal with "light tunneling effect".

1. We got the condition of enhancing "light tunneling effect" in the 1D-photonic crystal.

2. In using "Light tunneling effect", the front velocities quicken more than the light that propagates in the vacuum.

Pseudo-spin Kondo effect in a capacitively-coupled parallel double quantum dot

February 18th by Yuma OKAZAKI





A honeycomb structure appears due to interdot coupling V_{at}= 2.5 μV (77.7Hz) NL NL+1 NL+2 rt-0 N_R+2 N_R+1 NR +V. V. $(N_{L}, N_{R}+1)$ (N, +1,N_a+1) -lock-in-Ø Uine л QPC can deter differential conductance the number of electrons in the both QDs. $(N_{L}+1,N_{R})$ $(N_L N_R)$ $\label{eq:2DEG:n=2.2e11/cm^2, \mu=2e6 cm^2/Vs} \\ \mbox{Single Hetero junction 90nm below the surface.} \\ \mbox{Measurement temperature : 40 mK} \\ \mbox{}$ intedot Coulomb blockade









 $\begin{array}{ll} U_{\rm inter} = 270 \; \mu {\rm eV} \\ \varGamma = 260 \; \mu {\rm eV} \\ t = 20 \; \mu {\rm eV} \end{array} \quad \mbox{interdot tunneling} \\ \label{eq:linearized}$

ZBP was observed in both the on-site and intedot blockade. on-site : FWHM = 70μeV interdot : FWHM ~ 84μeV



Peak splitting in a finite inplane magnetic field



Why B = 1.25T ? A storg magnetic field affects the orbital state, due to poor alignment of the 2DEG. This magnetic field did not change the orbitals: Onsite Coulomb blockade Double peak splitting was observed, indicating a formation of the spin Kondoo singlet: hterdot Coulomb blockade ZBP were slightly troaden, might be coming from weak contribution of spin? No peak splitting was observed, indication for the Kondo correlation.

Why pseudo	-spin dor	ninant?	,	
The Kondo singlet is originating	initial	virtual	final	Effective coupling
from an effective coupling which resulting from second-order spin or pseudo-spin flip process.	Í ∗ ⊧ ⊧	┋┥╢┥┇	┇┼╢┼┢	$J = \frac{V^2}{U - U_{\underline{\bullet}} \cdot / 2}$
-+\++\#+		┪╣╪┢	┋╪╢╪╞	$J = \frac{2V^2}{U_{\rm max}}$
				$J = \frac{2 \mathcal{V}^2}{U_{inter}}$
пр				$J = \frac{2V^2}{U_{inter}}$

The right lower state has not spin. Thus spin and pseudo-spin can not simultaneously flip. When $U > U_{\rm state}$. The effective coupling of pseudo-spin is stronger than spin

Results

Double dot device was fabricated and measured.

Kondo effect was determined from ZBP and Kondo scaling at both on-site and intedot Coulomb blockade.

No peak splitting was observed at interdot Coulomb blockade in a in-plane magnetic field, indicating that the pseudo-spin dominantly contributes to the Kondo correlation.

I showed you why pseudo-spin domonant.

The Physical Society of Japan spring meeting 21/3(21aHV-2)
 Other Society of Japan spring meeting 21/3(21aHV-2)
 Other Society Soci





Observation of 8B Solar Neutrinos with KamLAND

Yoshiaki Kibe

Department of Physics, Tohoku University, Research Center for Neutrino Science



(1), History



The Standard Solar Model (SSM) attempts to describe the solar processes. Various solar neutrino experiments have yielded results in conflict with SSM predictions.

The most plausible solution for this anomaly is neutrino oscillation, which describes neutrino flavor changing while the neutrino propagates. An additional effect happens in the Sun, the neutrino weakly interacts with other particles and the behavior of the oscillation changes. This effect is called the MSW effect.

KamLAND revealed a significant deficit of electron anti-neutrinos from distant power reactors and suggested the LMA-MSW solution.



(4), Backgrounds

The dominant background in ⁸B solar neutrino region(> 5MeV) are from the β -decays of light isotopes produced by muon spallation. The following veto is applied for the light isotopes produced by muon spallation.

- 1 sec whole volume veto after muons
- 5 sec whole volume veto
- (residual charge > 10⁶ p.e.)
- 5 sec veto of 3 m cylindrical volume along

muon track (residual charge < 10⁶ p.e.) These cuts are effective for isotopes with half lives less than ~1sec, so ⁸B, ⁸Li or ¹¹Be are still significant.



The production rate of ⁸Li and ⁸B are estimated with simultaneously fitting the time difference from preceding muon because the lifetime of ⁸Li (1.21 sec) is close to that of ⁸B (1.11 sec). The energy spectrum of βdecays of ⁸Li and ⁸B is also fitted simultaneously. The estimated event rates are ⁸Li = 0.52 ± 0.15, ⁸B = 0.21 ± 0.04 ev/d/kt.



The production rate of ¹¹Be is estimated from the fitting of time difference from muons with tight track correlation. The estimated event rate of ¹¹Be is $0.91 \pm 0.30 \text{ ev/d/kt}$.



The external γ -rays come from (α , n) and (n, γ) reaction produced in mainly surrounding rock and stainless. The estimation for external γ -rays make use of MC simulation, providing 0.24 ± 0.04 and 0.38 ± 0.14 ev/d/kt for before/after purification.

(2), Motivation



KamLAND has the possibility to observe ⁸B solar neutrinos with an energy thresold of 3.5 MeV. This provides the transition of vacuum to matter oscillation, which has not been observed yet. The measurement of solar neutrinos provides the test for the theories of steller evolution and structure. Furthermore, The direct measurement of the transition of vacuum to matter oscillation will be good test for the oscillation parameter space.



KamLAND (Kamioka Liquid scintillator Anti-Neutrino Detector) 17inch PMTs; 1325 tubes 20inch PMT; 554 tubes Outer Detector; 229 20inch PMTs 1000 tons Liquid scintillator; Dodecane(80%) + Psuedocumene(20%) + PPO(1.36g/l)

Number o	t target;	3.429 x	1032	electror

(6), Uncertainty		100
	Before Pur.	After Pur.
Fiducial volume	3 m cylinder	3.5 m spherical
Livetime	952 days	63 days
	Uncertainty(%)	
Trigger Eff.		
Fiducial volume	2.17	11.23
Livetime Calc.	0.03	0.03
Energy Threshold	2.44	3.00
χ2TQ Cut	0.30	0.30
Number of Target	< 0.1	< 0.1
Cross Section	0.50	0.50

Fiducial volume before purification is defined as a cylindrical due to the effect of external γ -rays. On the other hand, fiducial volume after purification is defined as a spherical since ²⁰⁸TI is a dominant background with an energy threshold of 3.5 MeV, then the signal-to-noise ratio gets better.





The v-e⁻ elastic scattering is used to detect the solar neutrino events. In anti-neutrino detection like reactor neutrinos, the inverse β decay reaction, v_e+p -> e⁺+n and the neutron is thermalized in the liquid scintillator. Prompt signal from e⁺ gives information on the incident v_e energy. Neutron capture on hydrogen emits delayed signal with ~200µsec. This delayed coincidence method is effective to the background rejection. The way to detect the solar neutrino is single event, so it is necessary to estimate the background in the ⁸B solar neutrino energy region in details.

(5), Purification

The newly developed purification system reduces the radioactive impurities in a liquid scintillator. This purification system is mainly composed of the distillation tower and nitrogen purge tower. The purification provides a lower energy threshold.



Analytical reduction of ²⁰⁸Tl makes use of the delayed coincidence between prompt ²¹²Bi αdecay and delayed ²⁰⁸Tl β-decay. The lifetime of ²⁰⁸Tl is 3.053 min. Due to this analytical reduction, we can explore 8B solar neutrinos with an energy threshold of 3.5 MeV.

(7), Conclusion

The study of ⁸B solar neutrinos with KamLAND at energy threshold of 3.5 MeV will provide the observation of the transition of vacuum to matter enhanced neutrino oscillation in addition to the result of ⁷Be solar neutrino observation (energy is 0.875 MeV). For instance, 2-year livetime provides statistical uncertainty of ~10%.



The 2nd Scienceweb GCOE International Symposium on "Weaving Science Web beyond Particle-Matter Hierarchy/"

z~3 LYMAN BREAK GALAXY CLUSTER SURVEY IN SSA22 WITH VIMOS

Katsuki Kousai

Physics, D2, Tohoku University

We report our redshift survey of z⁻³ Lyman break galaxies (LBGs). We have obtained spectral redshifts (spec-z) of 94 LBGs in SSA22 (2:17:34, +00:15:04) 912 arcsec² field to study three dimensional distribution of LBGs. SSA22 is a field in which Steidel et al discovered high density region of LBGs at z⁻³.09 (Steidel et al 1998). They determined spec-z of 99 LBGs at z⁻³ in 162 arcsec² field. They found density peak of LBGs from their redshift distribution. We have detected a lot of Lyman alpha emitters (LAEs) at z⁻³.09 in 912 arcsec² field which contains the high density region discovered by Steidel et al, in our narrowband survey for LAEs with Suprime-Cam to find a large scale structure of LAEs (Hayashino et al 2004). Also, we are carrying out LBG redshift survey in our LAE survey area. We obtained spectral redshifts of 94 LBGs in 2006 and 2008 VIMOS observation. We have found spike not only at z⁻³.1, z⁻³.3 and z⁻³.7 in our redshift distribution of LBGs.

1. High redshift galaxy

There is two major types of high redshift galaxy Lyman α emitter (LAE) and Lyman break galaxy (LBG).The LAEs has a strong Lyman α emission line (λ_{rest} =1216 Å). The emission line are observed by narrow-band imaging technique. The LBGs are observed with strong break at λ_{rest} = 912 Å ; Lyman break. These UV photons of z > 3 galaxies are shifted to visible wave length by redshift. We obtained these candidate by observation of visible range.



2. SSA22

elative Dec [an

10

0

10

VIMOS2006 Steidel et al 2003

VIMOS2008

20

ve RA [arcmin]

z=3.1LAE high density region VIMOS2008 field of view

Steidel et al discovered high density region of Lyman Break galaxies (LBGs) in SSA22 at z=3.09 (Steidel et al 1998). They observed 81 arcsec2 × 2 in their spectroscopic LBG survey. We have detected Lyman alpha emitters (LAEs) at z=3.06~3.12 in 912 arcsec² field that contain the high density region which was discovered by Steidel et al (Hayashino et al 2004).



5. Three dimensional Spatial distribution

elative Dec

10

10

relative RA [arcmin]

20

 VIMOS2008 z=3.275-3.325
 VIMOS2008 z=3.325-3.375

 VIMOS2006 z=3.275-3.325
 VIMOS2006 z=3.325-3.375

 Steidel et al 2003 z=3.275-3.325
 Steidel et al 2003 z=3.325-3.375

VIMOS2008 field of view



4. Redshift distribution

We obtained spec-z of $34\ \text{LBGs}$ by the data that observed at $2006\ \text{and}\ 60\ \text{LBGs}$ by the data that observed at 2008.

We have found spike at z=3.1 and z=3.3 in the VIMOS2006 redshift distribution. We found another spike at z=3.7 in the VIMOS2008 redshift distribution. The magenta histogram represent the over all redshift selection function , normalized to the observed number of galaxies . We derived the selection function by Monte Carlo simulation.



Property of spike

- [·	VIMOS2006			· [\	/IMOS2008		
L	redshift	LBG / (selection function)	significance		redshift	LBG / (selection function)	significance
	3.1	5.1	5.1 σ		3.1	4.0	5.0 σ
	3.3	3.2	2.4σ		3.35	2.3	2.1 σ
					3.75	2.7	2.4σ



10

relative RA [arcmin]

VIMOS2008

20

VIMOS2008 field of view

arci

ive Dec

10

- We have found spikes not only at z=3.1 discovered by Steidel et al but also at z=3.3 and 3.7 in our redshift distribution of LBGs.
- Most of the LBGs at z=3.075~3.125 distribute inside the high density region of LAE skymap at z=3.06~3.12.
- It seems that LBGs in other spike z=3.3 and 3.7 are also clustered in 2D sky maps.

ational Symposium on "Weaving Science Web beyond Particle-Matter Hierarchy" The 2ndGCOE Inte PO no. 11.





(a) Al₇₅Cu₁₅V₁

0.3e

 L_3

72 74 74 Energy loss (eV) M. Terauchi, *et al.*, Phil. Mag., 87, Nos.18-21, 2947, (2007).

EF

Abstract After the discovery of an icosahedral symmetry material (quasicrystals) in a melt-quenched Al_6Mn alloy [1], great effort has been spent for understanding the presence of quasiperiodic structured materials. In recent years, Hume-Rothery mechanism, which predicts an existence of a pseudogap around Fermi level ($E_{\rm F}$), is accepted as a major reason for the stabilization of quasicrystals. level (E_p), is accepted as a major reason for the stabilization of quasicrystals. The presences of pseudogap structures in quasicrystals were experimentally confirmed by X-ray photoemission spectroscopy and EELS. EELS experiments also pointed out characteristic chemical shifts in al L-shell excitation spectra of AL-based quasicrystals [2], which suggested a decrease of valence electron charge at Al sites. Recently, a covalent bonding nature in quasicrystals was reported by MEM/Rietveld analysis of 1/0-Al₁₂Re and 1/1-Ar₃Re₅Si₁₀ [3]. Thus, it is interesting to investigate the relation between chemical shift and bonding nature of quasicrystals. On the other hand, the extrameting measurement, ed electronic dructores of the hemeel [4] transition systematic measurements of electronic structures of the 'normal' Al-transition metal (Al-TM) alloys help the understanding of the bonding nature of the Albased quasicrystal. D. Schechtman, et al., Phys. Rev. Lett., 53, 1951, (1984). [2] M. Terauchi, et al., Phil. Mag, 87, 2947, (2007).
 K. Kirihara, et al., Phys. Rev. B, 64, 212201, (2001).

Background 3: MEM/Rietveld analysis

Covalent bond

→ Decrease of valence electron charge ??

High resolution EELS-TEM instrument

Chemical shift (EELS)

Covalent bonds between Al atoms

1/1-Al72.5Re17.5Si10

 $1/0-Al_{12}Re$

Interatomic bonds: Metallic



ground 1: Quasicrystals and Approxim

and crystalline order by using EELS and SXES Approximants vs Chemical shifts I.

0.2eV →

Ouasicrystal shows chemical shift

Presence of covalency

Decrease of valence electron charge

- $1/0-Al_{12}Re$, $1/1-Al_{73}Re_{17}Si_{10}$
- II. Crystalline order vs Chemical shifts Am, QC, Cryst phases of Al53Si27Mn20
- III. Normal Al-TM alloy vs Chemical shifts B2 structure $(Pm\overline{3}m)$ Al-TM (TM=Fe, Co, Ni, Pd, Pt)

Purpose of this study

Investigate the relation between chemical shifts

I-a. Approximants vs Chemical shifts: EELS



II-a. Crystalline order vs Cher Al L-shell $(2p) \rightarrow C.B.(3s, 3d)$ Cryst 0C arb units. +0.4eV Am ~0eV

Energy Loss (eV) Chemical shift: QC phase



AIPt: chemical shift to smaller binding energy side (increase the valence electron) nount of shift: AIFe, AIPt>AIPd≒AICo≒AINi≒pure Al AlFe, AlPt: ch









nits)

arb.

AlFe: shift to small binding energy side nt with EELS re

Comparison with Pure Al Pure Al 1/0, 1/1 1/0-Al₁₂Re +0.4eV +6eV ~0eV 1/1-Al₇₃Si₁₇Re₁₀ +0.5eV +6eV ~0eV This result suggests that the covalent

Results I: Approximants vs Chemical shifts

Al-Ka emission

Re-Mß emission

Grat

Analyzed area: ~100nm¢ Energy dispersion for X-ray

0.5eV/pixel (640eV X-ray) 1.7eV/pixel (1.5keV X-ray) 2.2eV/pixel (1.7keV X-ray)

Al L-shell excitation

electrons between Al sites was provided from Si sites in 1/1 approximants.



Si-Ka emission

Pure Si

+3eV

Re	sults Ⅱ: Crys	talline order v	s Chemical sh	ifts	
	Al L-shell excitation	Al-Kα emission	Si-Ka emission	Mn-La emission	
Comparison with	Pure Al	Pure Al	Pure Si	Am, QC, Cryst	
Amorphous	~0eV	~0eV	~0eV	~0eV	
Quasicrystal	+0.4eV	+4eV	+6eV	~0eV	
Crystal	+0.1~0.2eV	~0.5eV	~1eV	~0eV	
Chemical shift is characteristic for QC phase					

 \rightarrow Al and Si atoms change those valence electron states Mn site does not change its valence electron charge



The relations between chemical shifts and crystalline order have been investigated by using EELS and SXES I. Approximants vs Chemical shifts:

- Al, Re: Chemical shift in 1/0, 1/1 but no difference Chemical shift in 1/1
- Si: Added Si supply the charge for the covalency
- II. Crystalline order vs Chemical shifts:
 - Al, Si: Chemical shift is characteristic for QC phase
 - No difference \rightarrow Because of localization of 3d electron ?? Chemical shift

 \rightarrow Decrease of valence electron charge

- Covalent bonds between Al or Si atoms in QC
- III. Normal Al-TM alloy vs Chemical shifts: AlFe, AlPt: Chemical shift to smaller binding energy side Charge transfer, different from quasicrystal

Larger atomic radius of TM \rightarrow Larger shift

Comparison of membrane physical property changes between DMPC membrane and DMPE membrane induced by melittin

Atsuji Kodama

Introduction

Melittin, an amphipathic peptide composed of 26 amino acid residues from ho has been widely used in the study of lipid-peptide interaction in biomembrane



orted that the strongest est interaction of melittin observed in 1,2-dimyristoyl-sn-) membrane among PC membranes which have different of matching of hydrophobic regions between melittin and we have investigated the interaction of melittin with PC as ength of fatty acyl chains with respect to effect of melittin It has h glycero length c DMPC PE men investigated the interaction fatty acyl chains with resp s. In this study, we h

Differential Scanning Calorimetry



Methods

Synthetic phospholipids were purchased from Avanti Polar Lipids (Alabaster, AL). Melittin purchased from Sigma Chemical Co. (St Louis, MO). Laurdan and Prodan were purchased I Molecular Probes (Eugene, OR).

Vesicle Pre ration

<u>Testice representation</u>. The desired phospholipid was taken from the chloroform solution into a test tube. In the study of GP measurement, luardan or prodan dissolved in DMSO was also added (0.5mol%). The solvent was evaporated first by a nitrogen stream, then under reduced pressure overnight. The lipid was dispersed into buffer (HEPES S0mM, EDTA, SMM, PH 7.3) containing desired concentration of meliftin in a buffer/HEPES S0mM, EDTA, SMM, PH 7.3) containing desired concentration of The total lipid concentration was 1mM for DSC and 0.2mM for fluorescene spectrometry.

Differential Scanning Calorimetry

Calorimetric scans were performed with a differential adiabatic microcalorimeter (Privalov calorimeter DASM-4, Sinku Riko, Yokohama, Japan) at a heating scan rate of 0.5K/min in 0.5ml cells under a pressure of 2.0 atm to prevent bubble formation. Fluo

scence spectrometry were performed with an AmincoBowman Series 2 luminescence meter (SLM Instruments, Urbama, IL) and rf-5000 (shimazu). The temperature of the was controlled by a water-circulating bath. Fluores

DSC thermograms of DMPC/Melittin and DMPE/Melittin



DSC thermograms of membranes with additives (2)



Temperature dependence of GP value for DMPC/Melittin and DMPE/Melittin



Changing of GP value for DMPC in phase transition became moderate with in of melitin, whereas GP value for DMPE changed drastically in the presence or compared with in the absence of melitin at the phase transition temperature. of melitti

Model for interaction of melittin with PC and PE membrane

atura did



The reason why transition ten

Melittin inserted into membranes and the system achieved equilibrium

not change with increasing of concentration of melittin has to be revealed.

Melittin RANN

PE membrane

Melittin oriented pallarel to the membrane at

Melittin oriented pallarel to the membranes disturbed penetration of water into membranes

Heating - cooling cycles induced insertion of melittin into membranes

Melittin · European honey bee venom

26 amino acid residues
 α-helix structure (interact with membrane)



Phospholipid $DMPC(diC_{140}) T_m = 23^{\circ}C$

Materials

-Y \bigcirc

. V



The scan times dependence of DSC thermograms









Laurdan spectra in DMPC membrane



Generalized Polarization (GP)

Temperature dependence of λ_{max} for tryptophan residue in DMPC and DMPE



In DMPC/Melittin systems, λ max at all mole fractions of melittin were between 330nm and 335nm through all temperature ranges. In DMPE/Melittin systems, λ max were at about 345nm.







Summary of DSC measurment

The dependence of concentration of melittin on DSC thermograms

DMPC/Melittin

DMPC/Mentum The behaviors of DSC thermograms could not be explained qualitatively by the equation because the peak position of phase transition did not change even though it became broadened with increasing of concentration of melittin.

DMPE/Melittin The behaviors of DSC thermograms could be explained qualitatively by the equation at K < 1

The dependence of the scan times on DSC thermograms

DMPC/Melittin

There was no dependence of scan times on DSC thermograms

DMPE/Melittin

The behaviors of DSC thermograms could be explained qualitatively by the equation at K < 1.

It indicates melittin inserted into with increasing of DSC scan times.

Tryptophan fluorescence

Hydrophobic environment

Shift to shorter wavelength



300 300 400 420 3cWavelength(nm) λ max was obtained from tryptophan spectru fitted by log-normal distribution. [Alexey S. Ladokhin et al (2000)]

DMPE (diC14:0) Tm = 49°C

DMPC has larger headgroup than DMPE

J.M.Sturtevant(1982) $\frac{T_0}{T} = 1 + RT_0 \left\{ \frac{1}{\Delta H_{vH}} ln \left(\frac{1 - \alpha}{\alpha} \right) - \frac{ln X_2}{\Delta H_{cal}} \frac{1}{\frac{K}{1 - K} + \alpha} \right\}$ Red part of the equation Blue part of the equation van't Hoff equation for two state transition van't Hoff enthalpy Mole fraction of LC phase Midle temperature of phase transiti

DSC thermograms of membranes with additives (1)



• The partition of additives for water phase is not taken into account. partition coeeficient of melittin (membrane / water) --- 10000



Reference

[1] H. Kusunose, JPSJ Online - News and Comments (Sep. 10, 2007)
 [2] Y. Kuramoto, PTPS 176 77
 [3] D.Mannix et al., PRL 95 117206

Acknowledgement

Author thank to Dr. Hironori Nakao for helping us at BL-4C, Photon Factory.

Study of geo neutrinos with KamLAND Yuri Shimizu, D3, RCNS

Geo-Neutrino

Radioactive sources in Earth emit electron antineutrinos (geo-neturinos) via beta-decays.

 $\begin{array}{rcl} ^{238}{\rm U} & \rightarrow & ^{206}{\rm Pb} + 8^{4}{\rm He} + 6{\rm e}^{-} + 6\bar{\nu_{e}} + 51.7\,{\rm MeV} \\ ^{232}{\rm Th} & \rightarrow & ^{208}{\rm Pb} + 6^{4}{\rm He} + 4{\rm e}^{-} + 4\bar{\nu_{e}} + 42.7\,{\rm MeV} \\ ^{40}{\rm K} & \rightarrow & ^{40}{\rm Ca} + {\rm e}^{-} + \bar{\nu_{e}} \end{array}$

 ${}^{40}K + e \rightarrow {}^{40}Ar + \nu_e$

Geo-neutrino measurement is the only direct way to investigate chemical composition of the Earth's interior.

KamLAND detects geo-neutrinos with inverse betadecay because it has high sensitivity for low energy neutrinos.

KamLAND detector



Earth's internal structure

Seismic wave measurements have revealed the detailed structure of Earth.

Core : mainly Fe and Ni. No U and Th exist

<u>Mantle</u>: Unknown composition by direct measurement → estimated using CI chondrites



+ Seismic wave measurement + Lithology Reference Earth model Calculate geo-neutrino flux

Meaning of geo-neutrino observation

The geo-neutrino observation can directly measure the Earth's internal structure.

U and Th mass in the mantle <u>Radiogenic heat</u> from U and Th in Earth

Total heat flow 44 or 31 TW BSE model prediction 19 TW (U 8TW, Th 8TW, ⁴⁰K 3TW)

The contribution of radiogenic heat is about half of total heat flow. \rightarrow U and Th are main heat sources in Earth.

Problems and solutions for geo-neutrino observation

✓ Statistics

 Analysis using 2076.6-day data set (KL08+586.0days)

✓ Backgrounds

 LS purification reduced the (α, n) background by 1/16.

✓ Uncertainties

 LS purification caused a decrease of light yield. Before the analysis, reconstructed energy is calibrated for positions and time.

LS purification



Background Reduction



Energy correction

After the purification, light yield decreased by 20% and reconstructed energy has large time variation and an ununiformity.

That significantly affects an uncertainty of energy scale. \rightarrow The post-purification period is dead time...

Energy is calibrated for z-axis, off-axis and time using ⁶⁰Co gamma and n-capture gamma.

That suppresses the increase of the uncertainty of energy scale. Time variation and an uniformity are the same as before the purification.

		Be	fore	the	corre	ection	۱			
Z dependence	(Mean energy (MeV)	2.1 2.1 1.9 1.4 1.7	<u></u>		÷	Befo	ore the	e purif	ication 5%	C. Contaction (1)
Off-axis dependence	Mean anargy (MeV)	22 22 22 22 22 22 22 22 22 22 22 22 22	J.	6 10		- ;		 1	z (m) 4%	Contraction (1)
Time variation	Mean Energy [MeV]	24			مد ید ی			± Purifi	15% cation	C Designation (12)
			Dec.31 2002	Jan.01	Dec.31	Dec.31	Dec.31	Jan.01	Dec.31	

Z dependence
Off-axis dependence
2.00 2.00 2.00 2.00 2.00 2.00 2.00 2.00

Expected spectra







Accuracy is significantly improved due to the reduction of the $(\boldsymbol{\alpha},\,\boldsymbol{n})$ background.

The hypothesis of zero signal is excluded at over 3σ C.L.. (investigation \rightarrow evidence)

Restrict for radiogenic heat



Models of radiogenic heat BSE model 16TW

Fully-radiogenic model I
 37TW (H_U+H_{Th}+H_K=H_{Earth})

Fully-radiogenic model II
 44TW (Hu+HTh=HEarth)

Fully-radiogenic models are excluded at over 2σ C.L.

Conclusion

Geo-neutrino analysis was performed with 2076.6day data set (KL08+586.0days).

 $(\alpha,\,n)$ background was reduced by 1/16 resulted in the significant improvement of accuracy.

The energy shift due to the decrease of light yield was carefully calibrated.

The analysis provided the evidence of geoneutrinos and excluded fully-radiogenic models at the first time.

Pressure dependence of B-T phase diagram in heavy-fermion superconductor CeRhSi3

T. Sugawara¹, N. Kimura¹, H. Aoki¹, T. Terashima², F. Lévy³ and I. Sheikin⁴

- ¹Graduate School of Science, Tohoku University ²National Institute for Materials Science ³ DPMC, Université de Genève, Geneva, Switzerland
- ⁴Grenoble High Magnetic Field Laboratory, CNRS

There is a bend point on the superconducting transition line $(B_{c2}(T))$

Below the bend point field (B_{BP}), the form of $B_{c2}(T)$ s are quite similar.

Introduction

CeRhSi2 is a pressure-induced heavy-fermion superconductor discovered in 2005 by N. Kimura et al. It crystallizes in the BaNiSn3-type tetragonal structure which lacks an inversion center. It exhibits the antiferromagnetic ordering below the Néel temperature $T_{\rm N}$ = 1.6 K at ambient pressure. The superconductivity emerges above 0.2 GPa and its transition temperature T_c becomes comparable to T_N at P_c = 2.4 GPa.

We have measured the resistivity of CeRhSi₃ for magnetic field along the tetragonal c-axis under several pressures. We obtained a novel superconducting B-T phase diagram, suggesting the existence of multiple superconducting phases





Sample

Starting materials: Ce(4N), Rh(2N8), Si(5N) Czochralski pulling method in a tetra-arc furnace Anneal: 900°C, 2×10⁻⁶ Torr, 1 week RRR > 100 (l > 2000Å: from dHvA)

Experimental method

Pressure: Ni-Cr-Al / Cu-Be piston cylinder cell i- and n-propanol 0~2.85GPa

Resistivity: 4 wires AC method 0.02~1.6K, 0~16T(~28T at 2.85GPa)

> 4 wires DC method 1.5~300K, 0T (for obtaining the absolute values)

B-T phase diagram

above P_c.





 $B_{\rm BP}$ and the bend point temperature $(T_{\rm BP})$ have P-linear dependence

Considering a linear extrapolation, the bend point vanishes at 3.20 GPa.

The initial slope of $B_{c2}(T)$ at T_c , $B_{c2}' = -dB_{c2}(T)/dT |_{T=Tc}$ shows abrupt change at Po.

This and the bend on the $B_{c2}(T)$ -line suggest the existence of three (at least two) distinct superconducting phases and the superconducting state changes at Pc and the bend point





$\rho(T)$ under several pressures



ρ (μΩcm)

As increasing magnetic field, T_N becomes ambiguous and $\rho(T)$ turns to have a positive curvature.



Above B_{BP} , the slope of $\rho(T)$ changes at T_{a} . This behaviour is different from that of T_N.



2.85GPa



The change of the slope of $\rho(T)$ becomes obvious with increasing pressure.



Neutral meson photoproduction experiments with electromagnetic calorimeter complex FOREST

Reserch Center for Electron Photon Science, Tohoku University K. Suzuki

1. Motivation

Recently, a new baryon resonance was observed at W=1670 MeV in the $\gamma d \rightarrow \eta np$ reaction, as shown in Fig. 1. The resonance width is narrow compared with that of a typical baryon resonance. No indication appears in the $\gamma p \rightarrow \eta p$ reaction in the same energy region.



2. Accelerator and beam

For photoproduction experiments, we use a bremsstrahlung photon beam. We have an electron synchrotron, called STB, accelerating electrons up to 1.2 GeV. The bremsstrahlung photon beam is generated by employing a radiator placed just upstream from a bending magnet of the STB ring. The radiator made of a carbon fiber, 11 µm in diamitor, is inserted into the circulating electrons. Recoiled electrons are analyzed by the bending magnet and detected with the STB-Tagger II system which consists of 116 telescopes of two-layer scintillating fibers (Fig. 3). The energy $E\gamma$ of a generated photon is determined by the energy of the recoiled electron Ee and that of the circulating electron E_0 as $E\gamma = E_0$ - Ee.



4. Experimental data and analysis

Photo-production experiments have been carried out with a hydrogen/ deuterium target. Up to now, more than 2×10^9 events were obtained for each target. A typical trigger rate and DAQ efficiency are 1.8 kHz and 80%, respectively.



The π^0 and η mesons decay into 2γ -rays as soon as they are generated. The decay branching ratios of π^0 and η mesons are 99% and 39%, respectively. FOREST measures energies and positions of γ -rays comming from these mesons (Fig. 5). The $\gamma\gamma$ invariant mass can be calcurated by using the energy and position information for 2γ -rays as,





5. Future plan

The goal of this analysis is to determine the spin and parity of the new baryon resonance N*(1670). It is necessary to simulate the acceptances of FOREST to obtain the cross sections for the $\gamma p \rightarrow \eta p$ and $\gamma d \rightarrow \eta n p$ reactions. And a partial wave analysis will be made to determine the spin and parity of N*(1670).



respectively. They can identify charged particles, and measure the position of the incident particles. FOREST covers about 90% of 4 π [str], thus it has large geometrical acceptances for the $\gamma p \rightarrow \pi^0 p$ and $\gamma p \rightarrow \eta p$ reactions. FOREST is powerful in suppressing background events (missing γ -ray events from $2\pi^0$ or $3\pi^0$...) since the missing γ -ray happens with a low probability under the large geometrical acceptance.



To select $\gamma p \rightarrow \eta p$ and $\gamma d \rightarrow \eta np$ events, we require the following conditions: 1. 2γ -rays are deteced by Backward Gamma, and

2. 1 charge or 1 neutral is detected by SCISSORS III.

In the 2-nd requirement, we choose Mx < 1100 MeV events as indicated in Fig. 7 (Mx : missing mass of $\gamma N \rightarrow \eta X$) to ensure that the charged or neutral particle is a proton or neutron. Background events are suppressed through these event selections as shown in Fig. 8.



The η meson yield, $Y\eta$, for every tagged energy is calcurated by using counts of η mesons, $N\eta$, and counts and efficiency of the tagging counters, $N\gamma$ and $\epsilon\gamma$, as $Y\eta=N\eta/(N\gamma\,\epsilon\gamma).$

Figures 9 and 10 show the η meson yields of the $\gamma p{\rightarrow}\eta p$ and $\gamma d{\rightarrow}\eta np$ reactions.





<u>J. Hayakawa</u>^A, K. Muraki^B, G. Yusa^A

^A Department of Physics, Tohoku University ^B NTT Basic Research Laboratories, NTT Corporation







can reveal local configurations of electrons

Experimental setup



Photoluminescence of 2DEG at zero magnetic field





How to assign photoluminescence peaks ?





Longitudinal resistance in the vicinity of v=2/3



Simultaneous measurement







T = 150 mK I = 200 nA Laser power = 0.65 mW/cm²

Plan of experiments ~ micro PL ~ Micro PL Unpolarized Unpolarized Floton energy Ploton energy The peak intensity of charged excitons can be used as a probe to detect local spin polarization.

Local probe techniques for luminescence





Conclusion

We have found the peak intensity of singlet charged exciton tends to increase up to transition point and it tends to decrease beyond this point by decreasing electron density. On the other hand the peak intensity of triplet charged exciton tends to monotonically increase by decreasing electron density. Thus the peak intensity of charged excitons can be used as a probe to detect local spin polarization.

In order to acquire the image of spin domains, we developed a scanning optical microscope in a dilution refrigerator. This experimental set-up will enable us to directly observe the spin domains with micron order spatial resolution.







Yasuhiro Takemoto (Physics, D1, Tohoku Univ.) PO.19

Backgrounds



Electronics



DAQ System Construction



AQ System		
host tria	host hub	host mogu07
KinokoCollector	KinokoTransporter	KinokoTransporter
MoguraAnalyzer MoguraTriggerHorizontalBuilder	KinokoTransporter	MoguraTriggerAnalyzer MoguraMogAnalyzer
host_mog01	KinokoBuffer	KinokoViewer
KinokoCollector MoguraAnalyzer	KinokoVerticalBuilder	
MoguraHorizontalBuilder		
host_mog02	host disk	
host_mog03		
host_mog04	KinokoTransporter	
host_mog05		
host_mog06	KinokoRecorder	
KiNOKO is used for	or this DAQ System.	

Inline event building is the new idea Spreading builders and analyzers lessens loads.

Summary

Physics and Requirements



- before DAQ launches.
- After starting DAQ, tagging method and PSD method have to be established.
- For longer live time, we will start DAQ as soon as possible



Double Chooz PMT preparation

The 2nd GCOE International Symposium @ Tohoku University,18-19/Feb/2010 TABATA Hiroshi

What is Double Chooz?:



New concept of Double Chooz: The Double Chooz experiment is a neutrino oscillation experiment at Chooz nuclear power plant in France to accurately measure the last mixing angle θ_{13} . Double Chooz uses 2 identical detectors at different distances from nuclear reactors to cancel systematic uncertainties on neutrino flux and detector response. Short baseline reactor experiments can provide a clean θ_{13} measurement.

$$P(\overline{v}_e \to \overline{v}_e) \approx 1 - \sin^2 2\theta_{13} \sin^2 \left(\frac{\Delta m_{31}^2 L}{4E}\right)$$

Sensitivity: The current upper limit is $\sin^2 2\theta_{13} < 0.15$. We will measure the $\sin^2 2\theta_{13}$ with 5 times better sensitivity than the current limit measured by the previous CHOOZ experiment.

Double Chooz detector: Neutrino target: 10.3 m³ of liquid



Anti neutrino events: $\overline{v}_e + p \rightarrow e^+ + n$ Prompt: e^+ with $T_e = 1 - 8MeV$ in the target.

1

3

Delayed: radiative capture of *n* on Gd with $\tau \sim 30 \mu s$, r < 1m and $\sum E_{\gamma} \sim 8MeV$. Neutrino signals are estimated as 70/day @far and 500/day @near.

Preparation in the lab:



PMT delivery: PMT group purchased 800 PMTs for both detector. After acceptance test in Japan and Germany, cleaning and magnetic shield assembly were held in Germany. Then, 400 PMTs were delivered to Chooz. Other 400 PMTs for near detector are kept in Germany.

Preparation of the lab: Before starting PMT installation, clean and new scaffolding was installed in buffer tank. Then, clean tents were installed. One tent covered tank. Another tent was for dressing area and preparation area.

Cleanliness organization[1]: To keep lab clean, we cleaned all of the lab at least twice everyday. The tunnel in front of the lab was cleaned also. All people have to wear clean suit, clean shoes and clean glove.

Pretest[2]: PMT was tested just before installation. If some kind of problem is found in this test, the PMT is rejected. PMTs preparation[3]: After testing, legs to fix PMT to wall were mounted on PMT. PMTs were covered by black plastic bag to avoid light and dust. Position survey markers were put on PMT.

Lid PMT installation (Nov-Dec,2009):

Preparation of the buffer lid[1]: After acrylic vessel integration, PMT group worked for lid PMT installation. Since tank was open, it was covered by plastic sheet to avoid any dust. PMT installation[2]: People worked below buffer lid.

scintillator doped with 0.1 g/l of Gd.

Buffer: 110 m³ of non scintillating oil

Inner veto: 90 m³ of liquid scintillator

PMT (Photon Multiplier Tube): Main

(buffer) system. We developed low background PMT with HAMAMATSU.

responsibility of Japanese group is PMT

Gamma catcher: 22.3 m³ of liquid

scintillator

& 390 PMTs (10 inch)

& 78 PMTs (8 inch)

Worker paid attention very much not to touch installed PMTs.

PMT cabling[3]: To be passed through cable exit pipe on the buffer lid, we made bundle. Making good bundle, worker was unrolling PMT cable. Buffer lid closing[4]: After lid PMT installation, buffer lid was closed. Then PMT group routed lid PMT cable to divider module and connected just before Christmas of 2009.



Side and Bottom PMT installation (Jun-Jul, 2009):

PMT transportation[1]: PMTs were transferred to each place. During side and bottom PMT installation, we used a bucket to lowering the PMTs.

PMT installation[2]: As you can see, fixation rails had already been installed on the buffer tank in advance. 2 people mounted PMTs on the rail. We could install 25 PMTs everyday. This was hard work, because they worked in very small space between scaffolding and wall.



PMT cabling[3]: PMTs have a high voltage cable. Cable goes to divider module which installed on wall of the lab. 6 people worked for this cable routing. This task was hardest work because workers went up and down many times. Done[4]: All PMTs were installed until end of July 2009. After precise

until end of July 2009. After precise PMT position survey was held, plastic bags and survey markers were removed.





Future plan:

Electronics and DAQ integration: Electronics are under production and testing. It will take 3 months for electronics and DAQ integration. Detector closing: Inner veto lid is open now. Until detector closing, chimney construction (between buffer and veto tank), veto lid closing and chimney construction (above veto lid) have to be done.

Filling: Liquid and gas filling system is under construction by filling group. Since chimney is not so thick, difference of liquid levels is critical for acrylic vessel. Therefore filling operation is very important. Filling will be finished until July 2010.

Data taking: After filling , we will turn on the high voltage of PMTs to detect first neutrino for Double Chooz experiment. This is scheduled June 2010. Detector commissioning and data analysis will be continued. Near detector construction: Tunnel construction of near detector will start October 2010. Data taking with double detector is scheduled middle of 2012.

Effects of lanthanoid ion on phosphate head groups in DMPC membrane

Kouya Tamatsukuri¹, Tetsuhiko Ohba¹, Gen Sazaki², Kazuo Ohki¹ tamatuku@bio.phys.tohoku.ac.jp ¹Department of Physics, Tohoku University, Sendai, 980-8578 ²Institute for Materials Research, Tohoku University, Sendai, 980-8577







Generalized polarization (G.P.)











DPH



DMPC/Eu=0.1

The orientation of DPH was observed above 40 $^{\rm O}{\rm C}$ on microscopy, and there is no dependence of orientation of DPH on lanthanoid ion .



Conclusion

Addition of Eu³⁺ ORotational diffusion of acyl chains is facilitated polar head group is inhibited OMain transition temperature is shifted to higher →The occupied area of polar head group is relatively small in the liquid crystalline phase. Othere is no dependence of orientation of acyl chain or polar head group on lanthanoid ion. OUnder the static magnetic field, phospholipids became

OUnder the static magnetic field, phospholipids became orientated parallel to the static magnetic field.











Status of double beta decay experiment with KamLAND

Azusa Terashima

Research Center for Neutrino Science, Tohoku University



In detail ⇒ please see Takemoto's poster!



Masatoshi Toda Physics, D3, Tohoku University February 18, 2010

Wormlike Micelle





Linear Rheology -- Wormlike Micellar Case-104

10 10 1 H. Rehage & H. Hoffmann (1988) One Separated Relaxation Mode Exists!! What is it ?

A Proposed Relaxation Mechanism (Phantom Crossing Model) T. Shikata (1987)







Hybrid Model —Purely Dissipative Equations-(Non dimensional) $$\begin{split} \frac{\partial \mathbf{v}}{\partial t} &= -\nabla p + \nabla^2 \mathbf{v} + \sum_{i=1}^{N} \Big[\frac{\mathrm{d} \mathbf{r}_i}{\mathrm{d} t} - \mathbf{v}(\mathbf{r}_i) - \zeta_i \Big] \delta(\mathbf{r} - \mathbf{r}_i) + \boldsymbol{\theta} \\ (\text{at Low Reynolds Number)} \quad & \frac{\mathrm{Body Force}}{T_i} \Big] \frac{\partial \boldsymbol{v}}{\partial t} \\ \end{split}$$



Bending Elasticity of Membranes



.

11

Network Structure

without shear Percolation Prob. = 0.955



Cluster Size Distribution



Conclusion & Future Problems

(1) Modeling of Wormlike Micellar System
Particle-Field Hybrid Model
Micellar Network Formation
(2) As a Result of a Coarse Graining, Discontinuity in Interparticle Interaction Potential
Failure of Canonical Distribution

New Stochastic Methods







Pi-mesonic decays of Λ hypernuclei Yoji Nakagawa (physics, D1)




理学研究科物理学専攻 中村俊太郎

素粒子物理学には標準理論と呼ばれる様々な実験結果を説明する 理論がある

しかし、近年その存在が確定的となった暗黒物質(DM)は説明するこ とができない 標準模型を超える物理を知ることが目標



Rパリティ(標準理論粒子+超対称粒子ー)が保存してれば最も軽 い超対称粒子(LSP)は安定になる



超対称粒子はまだ見つかっていないので超対称性は破れている

超対称粒子は対の粒子よりも重い neutralinoなどが暗黒物質の候補になる

超対称粒子の質量スペクトルは超対称性の破れ(SUSY)によっている どのようにSUSYが伝わっているかを理解することは重要

注目を浴びているのがmirage mediation



これら二つの寄与は常にあるのでmirage mediationはとても 自然なSUSYのメカニズム

しかし、mirage mediationには問題がある

moduli 問題











 $\psi_{3/2}$: gravitino $\tilde{\chi}$: neutralino LSP

neutralinoは大量に存在するので対消滅する



宇宙は膨張しているので対消滅は $n_{\hat{\chi}}(\sigma v) \simeq H(T_{3/2})$ のとき終わる $(T_{3/2}: gravitino崩壊時の熱浴の温度)$ これが最終的な $\tilde{\chi}$ の残存量

T_{3/1}が十分高くない(10MeV)ため消滅効果が一番効くwino LSPでも $m_{\hat{\chi}} \frac{n_{\hat{\chi}}}{s} = \frac{m_{\hat{\chi}}}{\langle \sigma v \rangle M_{\rm Pl} T}$ は暗黒物質量を超えてしまう

我々はmirage mediationを拡張することでmoduli問題を解決する

Axionic mirage mediation = mirage mediation + axion



Axionic mirage mediationにはaxionの超対称粒子であるaxinoも 入ってくる



O(100)GeV - O(1)TeV

axinoの質量は

$$\begin{split} m_{\tilde{a}} &= \frac{5N}{16\pi^2} \lambda^2 A_{S \Phi \Psi} \Big\langle \frac{\hat{S}^{\dagger}}{\hat{S}} \Big\rangle \simeq \mathcal{O}(1) \mathrm{GeV} \\ & \frown \qquad \mathcal{O}(m_{\mathrm{soft}}) \end{split}$$

axinoがLSPになる **DMの**候補

axinoはneutralinoよりも軽いのでmirage mediationのmoduli問題を 解決できそう

moduliの崩壊で生成されたaxinoが現在のDM量を説明できるどうか を見る必要がある

axinoの生成過程



axinoの残存量

moduli問題の解決



質量がO(100)MeV のaxinoで暗黒物質を説明できた

($\Omega_{\rm DM} h^2 = 0.105^{+0.007}_{-0.013}$) Spergel, et. al '06

mirage mediationは有望な超対称性の破れの機構だが moduli問題と呼ばれる宇宙論的な問題があった

我々はmirage mediationにaxionを加える拡張をすることによって neutralinoよりも軽いaxino LSPを実現できた

質量がの(100)MeVのaxinoで暗黒物質を説明でき moduli問題も解決することができた

The Study of the Origin of Lyman Alpha Emitters with Large Wquivalent Widths

SSA22 AT Z=3.1

Our wide-area (2.4deg² in total) and deep narrow-band survey detected ~2000 LAEs and reveal that there are a number of large EW objects in "SSA22" and general fields EW₀>=240Å; 240 LAEs at SSA22, 95 LAEs at general fields). In order to discriminate the origin of large EW objects, first we should make large EW samples which include objects enhanced for EW by mechanism of Lyman alpha scattering and/or galactic superwind. We measured the both Lyman alpha eimission and continuum components by pseudo total magnitude (2.5 × kron radius of SExtractor software) as the Lyman alpha emission of these objects have extended shape. (Noteworthy, we measured Lyman alpha emission and continuum within a given aperture in previous work. When we focus on the EW at the exact position where the star-formation occur to know the stellar age, this measurement is effective.) In results, we found the larger number of high EW objects by this method than that by previous method and the ratio of large EW objects (EW₀=400-700Å) to small EW objects (EWo<100Å) is 1.6 times larger in "SSA22" fields than general fields. Furthermore we investigate the statistic properties of these objects such as the size of Lyman alpha emission, Luminosity Function and colors

1. Observation

We conducted wide-area (2.4deg2 in total) and deep narrow-band survey with Suprime-Cam of Subaru Telescope.

Field	Number of LAEs	Volume (Mpc ³)	Density (Mpc ⁻³)	Number density of LAEs
SSA22	1438	9.9 * 10 ⁵	1.4 * 10 ⁻³	in SSAZZ region is 1.5
General Fields	764	7.7 * 10 ⁵	0.99 * 10 ⁻³	general fields.
SDF	196	1.6 * 10 ⁵	1.2 * 10 ^{.3}	↓
GOODS-N	186	1.9 * 10 ⁵	0.98 * 10 ³	SSA22 region is a high-
SXDS	382	4.2 * 10 ⁵	0.90 * 10 ^{.3}	density region of LAEs

2. Calculation of EW

We calculated EW of our detected LAEs by two methods.

- If Ly α photons emitted from star-forming regions are scattered by neutral hydrogen gas and the emission regions are extended, A) focus on the EW of the exact position where star-formation occurs to know the
- stellar age "EWap": EW measured by Lyα emission and continuum fluxes within aperture=2"φ
- (psf=1".0 at SSA22)
- B) include objects enhanced for EW by mechanism of Lya scattering and/or galactic superwind
- 'EWto": EW measured by pseudo total magnitudes of Lyα emission (within 2.5*kron radius defined in NB image) and continuum (within 2.5*kron radius defined in BV image) as the Ly α emission of objects have extended shape (using SExtractor)



Relationship between EWap and EWto



mostly linear correlation EWap ~ 1.5*EWto → most LAEs have extended Lyas emission image as conservative case[EWtoNBdet], objects with extremely extended emission which makes EWto very large.

We consider EWto measured within radius for continuum flux defined in NB Most of objects are measured accurately while some objects are overestimated. It is necessary to check some large EWto objects.

Size of Lyα as a Function of EWto

1.8 1.6 1.4 1.2



Some large EWto objects have compact Ly α emission Others have extended Lya emission.

5. Luminosity Function of Ly α as a Function of EWto



There are no significant differences in the shape of Lya luminosity function between large EWto objects and small EWto objects.

6. EWto Distribution of LAEs



EWto Distribution as a Function of LAE-Density in SSA22



8. Discussion

Characteristics of large EWto objects in SSA22



Future work.

to investigate the $Ly\alpha$ - UV size ratio (it is necessary to increase the S/N value of continuum), the color of large EW objects, the luminosity function of large EW objects

Summary	We made the large sample of \sim 2000 LAEs in SSA22 and general fields. The number density of LAEs in SSA22 region is 1.5 times larger than it in general fields.
	We calculated EW of LAEs by two methods: "EWap"- aperture-photometry and "EWto" - total magnitude. We can newly find a large number of LAEs with high EWto objects which have extremely extended emission. There are the large EWto objects with compact Lyα emission and extended Lyα emission. The large EWto LAEs in SSA22 region have higher ratio than in General fields. → It is the unique characteristic of SSA22-HDR (?)

Ultrafast broadband THz response of photo-induced metallic state in charge ordered insulator α -(BEDT-TTF)₂I₃

H. Nakaya¹, Y. Takahashi¹, K. Itoh¹, S. Iwai^{1,2}, K. Yamamoto³, K. Yakushi³, S. Saito⁴ ¹Department of Physics, Tohoku University, Sendai 980-8578, Japan, ²JST-CREST ³Institute of Molecular Science, Okazaki 444-8585 ⁴Kobe Advanced ICT Research Center, NICT, Kobe 651-2492

Photo-induced insulator to metal (I-M) transition in correlated electron system



Objective Clarifying electronic nature of photo-induced metallic state → Ultrafast broadband terahertz (THz) spectrum (Charge-ordered gap ~ 50-100 meV) c. f. Hilton et al., PRL 99, 226401 (2007). , Averitte et al., PRL87, 017401(2001).







Setup of optical-pump THz-probe measurement



Time evolution of the broadband transient spectra



Spectral weight (low energy ⇒ high energy) = Formation of charge-ordered gap ?



In summary, we investigated the ultrafast dynamics of photoinduced I-M transition by using optical pump-broadband (2 - 36 meV) THz probe spectroscopy in charge ordered insulator α -(BEDT-TTF)₂I₃. At 124 K (τ_{co}), OD spectrum at 0.8 ps has a broad absorption increase, exhibiting generation of the metallic state. Then, the spectral weight shifts to high energy region within several picoseconds. Such spectral behavior in THz region indicates the formation of the charge-ordered gap.







Shape Memory Effect Induced by Magnetic-Field Rotation in MnV₂O₄ Spinel Compound

Yoichi Nii, Takamichi Yagi, Nobuyuki Abe, Kouji Taniguchi¹, Hajime Sagayama¹, Taka-hisa Arima¹ Department of Physics, Tohoku University

¹Institute of Multidisciplinary Research Laboratory, Tohoku University



Magnetic-field-induced shape memory effect

A high speed control of macroscopic shape is possible by using this process where temperature variation are not necessary.



The ω meson photoproduction on the nucleon in the threshold region Rvo HASHIMOTO (Department of Physics, D3) Research Center of Electron Photon Science, Tohoku University

The ω meson is one of vector mesons having a spin parity of 1. Studying near-threshold ω meson photoproduction is interesting theme to search for new baryon resonance states. Recently, the CLAS collaboration has reported that a missing baryon state with $J^p = 5/2^+$ contributes to ω photoproduction around the incident photon energy of 2 GeV. Paying attention to other vector mesons, the anomaly of the cross section near the threshold region is observed for ϕ photoproduction by the CLAS collaboration. A local maximum has observed around Ey = 2 GeV. Is such

a local maximum observed in the ω photoproduction cross section? We plan to measure it from threshold up to 1145 MeV.



vs [GeV] ection Data: •: thi

18]; ×: [5]; \Box : [19]. Partial-wave decomposition [17]; de ; $J^P = \frac{1}{2}^+$; short-dashed line: $J^P \ge \frac{5}{2}^+$ (stemming xchange); dashed line: $J^P = \frac{1}{2}^-$; dashed-dotted line: Jdashed-double-dotted line: $J^P = \frac{3}{2}^-$.

Interesting topics of vector meson photoproduction



A local maximum appeared in the $\gamma p \rightarrow ~\phi~p~cross~section$ (CLAS collaboration)

How about ω meson?

T. Mibe et al.

 ω -photoproduction (a) ELPH

Eγ range 747.6 - 1147.8 MeV Photoproduction threshold 1109 MeV Detector FOREST in the Gev-y experimental hall Target Liquid hydrogen/deuterium Selected events 3γ from $\pi^0\gamma$ (and 1p)

[qi



section for the $\gamma p \rightarrow \omega p$ reaction from threshold up to 2 GeV. The P11 resonance contribution and OPE process are dominant. However, near the threshold, main contribution

G. Penner and U. Mosel calculated the total cross

for the total cross section is the D_{13} resonance. → We expect a local maximum in our data.



60

З

1000 1200

 $M\pi^0\gamma$ (MeV)

1080 1100 1120 1140 1160 Incident energy of γ (MeV)



Prototype fast imaging detector for Hyperon Nucleon (YN) -scattering

Ryotaro Honda Dept. of Physics M1, Tohoku Univ.



If an event was pp scattering, two barrels should be fired. However it's a rare case. Therefore pp event was selected by the number of tracks. A track with hit a barrel was analyzed by the same procedure for pC and **the other** side was not requested a barrel TDC. Virtual hit is created to determine scattering angl

consistent

Low-energy states of electrons in graphene

Motivation

pure Coul

 r/ℓ_B

V(q) : Transform of the Coulomb inter R_i : Guiding center coordinate of the *i*th electron

 $L_N(x)$: Laguerre polynomials



Electronic states of graphene in magnetic fields



Purpose







Comparison with the result of the HF theory



Summary

- By the use of DMRG method, we determined the ground state of electrons at various filling factors of the N=2,3 Landau level of graphene.
- By analyzing the (guiding center) pair correlation function, we obtained the reliable phase diagram of electrons in the N=2,3 Landau level of graphene.

Possibility of realizing the reentrant stripe phase around ν =0.3 of *N*=3 Landau level



Tatsuya Higashi Dept. of Physics, Tohoku U.



 $H_{s} = \sum_{l \leq j} \sum_{\mathbf{q}} e^{-q^{2} \left[L_{N}(q^{2}/2) \right]^{2}} V(q) e^{(\mathbf{q} \cdot |\mathbf{R}_{i} - \mathbf{R}_{j})} V_{eff}(r)$

 $H_{\rm G} = \sum_{i < j} \sum_{\mathbf{q}} e^{-\frac{q}{r_{\rm g}}} [F_n(q)]^2 (q) e^{i\mathbf{q} \cdot (\mathbf{R}_i - \mathbf{R}_j)}$

 $F_{n\neq0}(q) = \frac{1}{2} \left[L_n \left(\frac{q^2}{2} \right) + L_{n-1} \left(\frac{q^2}{2} \right) \right]$

The eigenfunction of nth Landau level of graphene $(n \neq 0)$

-1) (D)m)

 $\Psi_{Kn} = \frac{1}{\sqrt{2}} \left(\frac{ng}{\sqrt{2}} \right)$

The two effective interactions (H_s, H_s) are different.

Effective inter-electron interaction of Nth Landau level

Standard 2D systems

Graphene

CDW ground states of standard 2D electron systems













Comparison with the result of the HF theory DMRG (present study)









Investigation of the $n(\gamma, K^0)\Lambda$ reaction near the threshold

K. Futatsukawa for NKS2 collaboration, Tohoku Univ.



$p\pi^+\pi^-\pi^-$ Event Distribution





Summary

- The $\gamma + n \rightarrow K^{\theta} + \Lambda$ process plays a unique role in the investigation of kaon photoproduction mechanism.
- We have finished the construction of Neutral Kaon Spectrometer (NKS2)
- We performed the physics data taking with the deuteron target.
- The momentum distribution was obtained for two photon energy region, 0.9 to 1.0 GeV and 1.0 to 1.1 GeV.
 It is ready to contact the theoriest and to discuss these
- It is ready to contact the theorists and to discuss these results.

35. Traffic flow of two lanes with a bottleneck

Sho Furuhashi Dept. of Phys. , Tohoku Univ

Traffic flow

Macroscopic viewpoint •Burgers equation •Lighthill-Whitham model

Microscopic viewpoint Cellular automaton model

Cellular automaton (CA) model

models can reproduce the major

•184 model •Asymmetric Simple Exclusion Process (ASEP) Nagel-Schreckenberg (NS) model In one dimensional case, the funda

diagrams (flow-density relationship) of these

characteristic of the ones of real traffic flow.

·Optimal velocity model

Introduction



Road network



Bottleneck is a factor which interrupts traffic flow. At an intersection, cars on one lane interrupt cars on the other lanes. Road network consists of bottle cks linked each other

Objective

Real traffic flow is reproduced by the model (CA model) in the previous studies?

Road network is complex. CA model might be no good for road network.

We check the adequacy of the model.

Model

System

Two lane circuit having one lane in a part

 ρ_1, ρ_2

iental

low

- Total length of a lane L L
- The length of a bottleneck

Initial density

Car dynamics

ASEP with parallel update.

Two lane road which has a one lane road in a part

Simplified model of bottleneck



Previous studies •Y. Ishibashi and M. Fukui J. Phys. Soc. Jpn. 65 (1996) 2793 •M. E. Foulaadvand and M. Neek-Amal EPL. 80 (2007) 60002

Observation

Recorded with video camera



05 0. The cumulative distribution of Index for characterization the number of passing car The number of passing cars $N_i(j)$ Simulation Observation $L - L_b = 500, L_b = 1, p = 0.5$ The number of passing cars on first lane in j-th cycle. $N_1(j) \longrightarrow N_2(j) \longrightarrow N_1(j+1)$ $F(N), F(N_1), F(N_2)$ One cycle £ ρ₁ = ρ₁ = 0.3 Ν ρ₁ = ρ₂ = 0.5 Ν $N_1(k)(j=2k)$ No distinction N(j) $N_2(k)(j=2k+1)$ between lanes $F(N_i) = \int_{-\infty}^{\infty} f(x) dx$

> Numb of car

> > **Exponential distribution**

Cycle correlation



•Cumulative

Result and assignment

Result

The model agrees with real traffic flow in the two indexes. These are Poisson process.

 $(f: N_i \text{ distribution function})$

 $\langle (N_i(j) - \overline{N}_i)(N_i(j + \Delta) - \overline{N}_i) \rangle$

Assignment

We have to check that this result is independent of the density of cars.

When the density is over the critical value, traffic jam occurs. The mental condition of drivers in a traffic jam is different from of drivers in a free flow.

ASEP

·Hop to the next cell with probability *p*.

is occupied.

(Self-driven nature of cars)

•Stop when the next cell

(Excluded volume effect)

o(1/K)

 L_{b}

density [cars per site]











Proton elastic scattering of ⁹C at 290 MeV

Y.Matsuda Department of Physics, Tohoku University

The differential cross-section for proton elastic scattering on %C at 290 MeV/u was measured in inverse kinematics The experimental method, results and discussion are shown below.



3. Experimental results of the $H({}^{9}C,p)$ reaction at 300 MeV/nucleon



4. Comparison of the experimental data and some calculations

Density distribution

(1) Relativistic Hartree Approximation (RMF) [C.J.Horowiz and B.D.Serot, NPA368(1981)503.]

(2) Antisymmetrized Molecular Dynamics (AMD) [Y.Kaneda-En'yo and H.Horiuchi, Prog.Theor.Phys.Suppl.142(2001)647]



Optical potential

Relativistic Impulse Approximation (RIA) [D.P.Murdock and C.J.Horowiz, Phys. Rev.C35(1987)1442.]

(2) Hamburg G-matrix Approximation (G-matrix) [L.Rikus and H.V.von Geramb, NPA426(1984)496.]



5. Summary

The differential cross-section for 9C was measured in the momentum transfer region of 1 to 2 fm⁻¹

This measurement is the first at the energy region at which the proton has the longest mean free path in nuclei. The experiment was performed with very low background using the 5 mm thick solid hydrogen target.

The excitation energy resolution was about 1 MeV.

Compared with ¹²C data, the angular distribution does not show a clear diffraction minimum.

The angular distribution is compared with microscopic model calculations

Optical model description is performed with the RIA and the G-matrix approximation. Description of nuclear structure is done with the RMF and the AMD.

The unclear diffraction minimum is comparatively reproduced by the RMF. It is considerd that this reproduction is derived from that the tail of the proton distribution grow longer than that of the AMD result.

The RMF calculation is performed by adjusting the scalar mass so as to reproduce the matter radius 2.42(3) fm.

The calculation suggests that there might be a large differentce between the proton and neutron radii, which are 2.58 fm and 2.06 fm respectively.

Though more detailed discussion could be done by comparing with more backscattering data, a lack of data and an inaccurateness of nuclear reaction models at the angle make it difficult. In future, it is desired to improve these two points.

High Resolution and High Statistics Λ Hypernuclear Spectroscopy by the (e,e'K+) Reaction



External Gamma-ray Backgrounds of the KamLAND Detector

Yukie Minekawa Physics, D2, Tohoku University









41.

Electron diffraction study of the low temperature phase of Hollandite-type oxide K₂Cr₈O₁₆

Daisuke Morikawa Physics, D1, Tohoku University





Need more high-crystalline single crystal

Introduction

Objective

Parallel _ 100nm

Result.1-2

SAED

Average structure

3008

Diffuse streak in b* direction

acking faul

No super spot

8.

Streak to <110> direction

-Multiple-scattering calculation

etc) -Preprocessing of experimental data -PC clusters

group I4/m

temperature

Streak in SAED pattern

narp diffuse scattering at low

[100]? incidence SAED pattern

No I lattice type area exist

Stacking fault at (010) and (110) plane

Make low symmetry ([100]incidence CBED pattern)?

96 node PC cluster system were installed

[111]

L) (010) plan

Discussion.2-1

Result.4

Experimental.1

Recently, the high quality single crystal of Hollandite-type oxide K₂Cr₈O₁₆ has Recently, the high quality single crystal of Hollandite-type oxide $K_2(r_d)_{24}$ has been successfully made by high-pressure synthesis and the metal-insulator transition keeping ferromagnetism was formed at 95K [1]. This phenomenon is very rare in the strongly-correlated electron materials. The crystal structure of the low temperature phase and the origin of the metal-insulator phase transition have been unclear yet. The purpose of this study is to examine the phase transition at the nano-scale area of K2Cr8O16 using Convergent-beam electron diffraction (CBED) and Selected-area electron diffraction (SAED).

[1] K. Hasegawa et al., Phys. Rev. Lett. 103, 146403 (2009).

Crystal structure of the low temperature phase

Incident

electron beam

Convergent

Rocking curve Space group determination Crystal structure and electro

30K

[100] incidence SAED pattern

Other incidence SAED pattern (30K)

[112]

MBFIT

K. Tsuda et al., Acta Cryst. A (1999,2002)

Summarv

Metal-insulator transition keeping ferromagnetism No structure transition were found CBED pattern symmetry⇒ no contradiction for space

22-2 31-2

sulator transition

CBED

Backgroud.1 Hollandite-type oxide K₂Cr₈O₁₆



i. Visualization of orbital and/or charge ordering state direct determination of 3 dimensional electrostatic potential using CBED method spinel type (FeCr₂O₂), perovskite type (TbMnO₂, Pr_{0,2}Ca_{0,4}MnO₂, Sm_{0,2}Ba_{0,4}MnO₂) ii. Investigate the relation with Physical property and crystal structure RBaMn_{O2} series (MdBaMn_{0,0}etc) iii. Space group determination BaFe_{2x}(r_xAs_{2x})(Pr, La)FeAsO_{1x}, Pr, Pr_{1x}Sr_xFeAsO v. Polarization determination ii. AlGaN, GaN Ca.ca





Quantum Paraelectrics

The polarization arrangement in these materials is disturbed by the zero point vibration at extremely low temperatures. So they do not undergo a ferroelectrics phase transition. SrTiO₃,KTaO₃ etc.

of the Γ point.

Quasi Elastic Light Scattering (QELS)

and Broad Doublet

⇒The density of state of phonon

P.A.Fleury and K.B.Lyons

(u)/2#(GHz)

Rayleigh Scattering

Ani

10

-50 0 ω/2π(GHz) 50

-10

It is necessary to experiment in a wide temperature region and a wide frequency domain

increases at the Γ point.

(TO-phonon + TA-phonon)

KTaO₃ /#

J.D.Axe, J.Harada, and G.Shirane Phys. Rev. B 1, 1227 (1970).

B.Hehlen, A.L.Pérou, E.Courtens, and R.Vacher, PRL 75,2416(1995)

Second Sound?

It is necessary to decide the relaxation

time of τ_{μ}, τ_{N} in the phonon-phonon scattering.

SrTiO₃

50

Shift ω/2π(GHz)

то/

High

Lov

q=(ζ,0,0) 295K 90K 6K





Objective

We measure the spectrum in the Brillouin scattering region in SrTiO3 and KTaO3, and understand a temperature dependence of the spectrum according to the model of collective excitation of phonon.

•We aim to understand a picture of the phonon gas at each temperature, obtaining τ_p and τ_N from the analysis of the light scattering spectrum.



There is a possibility that heat is propagated as a wave in a gas.

 $Q + \frac{\partial Q}{\partial t} = -\kappa \frac{\partial \delta T}{\partial x}$ (Cattaneo's expression)

It becomes a wave equation including decay for the temperature. \Rightarrow A wave of temperature fluctuation (Second Sound).

 $\frac{\partial^2 \delta T}{\partial t^2} + \frac{1}{\partial t^2} \frac{\partial \delta T}{\partial t} = \frac{D_{\text{th}}}{\partial t^2} \frac{\partial^2 \delta T}{\partial t^2}$ 1 :propagation velocity (c : average speed) $\frac{1}{\sqrt{3}}c$ $\tau_{\rm R} \quad \overline{\partial t}$ $au_{
m R} \partial_{x^2}$ ∂_t^2 $D_{\rm th} = \frac{1}{3}c^2 \tau_{\rm R}$: thermal diffusivity coefficient







Temperature Dependence of τ_N and τ_R





Spectrum Expression by Extended Thermodynamics H.Dreyer and H.Struchtrup, Continuum Mech. Thermodyn. 5(1993)3-50

 $1/\tau = 1/\tau_{\rm N} + 1/\tau_{\rm P}$

 $S_{\text{total}}(q, \boldsymbol{\omega}) = P_1 S_1(q, \boldsymbol{\omega}) + P_2 S_2(q, \boldsymbol{\omega}) + \text{const}$

 $\tau_{_R}$ is decided from thermal diffusivity coefficient obtained scattering experiment in the time domain. ($\omega_0\!\ll\!1/\tau_R$) ned by the ligh



Conclusion

- . A narrow quasi-elastic scattering component in the high temperature region is a light scattering spectrum by thermal diffusion. It has been understood that the frequency of the resistive process falls when lowering the temperature in SrTiO3, and this thermal diffusion mode (the second-sound wave that over dumped) shifts to the second-sound wave (propagating mode).
- 2. It cannot be said that the phonon gas is in the state of the collective excitation in the observation scale of the actual experiment. It should be interpreted that the observed BD for KTaO3 only as the second order Raman scattering involving pairs of individual phonons.





$$\Gamma = D_{\rm in} q^2 = \frac{1}{3} c^2 \tau_{\rm in} q^2$$

Propagation as a wave of heat Underdamping Local thermal equilibrium



Light scattering by collective excitation of phonon in quantum paraelectrics Ryuta Takano

Physics, D1, Tohoku University

Energy Transportation in the Phonon Gas

< Thermal Diffusion and Second Sound >

< Diffusion of heat >

< Propagation of heat > The Normal process is predominant

is predominent.

The Umklapp process and the impurity scattering

A compression wave of the phonon gas is propagated.

= Second Sound Change in phonon number = Change in the temperature

 $(Umklapp+impurity+..+etc=Resistive \ process:\tau_R)$

Ţ



Outline

- What can we learn from chemical abundance (+ kinematics) of nearby stars regarding the formation of the Milky Way Galaxy
- The main results of our study on the chemical abundance of the outer stellar halo of the MW Sample selection using kinematics
 - Abundance results of key-elements Implication for the formation of the MW outer halo



Data reduction: Standard IRAF routines



Cumulative distributions of [a/Fe] ratio at -2<[Fe/H]<-1





Mg (SNe II) < Fe (SNe Ia) Epoch of the accretion

Early No time for Fe enrichment Recent

Enough time for Fe enrichment Recent acc

retior [Fe/H] Font et al. 2006

Lower VSF

Lanfranchi & Matteucci 2003

Early acc



- Methods: Subaru/HDS observation + homogeneous
- abundance analysis
- Bank Results:
- At -2<[Fe/H]<-1, the outer halo sample show systematically lower [Mg/Fe]
- The [Mg/Fe] ratios for the outer halo sample spans intermediate range between the inner halo and the MW satellites
- <u>The implication for the halo formation:</u>
 - Certain fraction of the outer halo was formed through merging/ accretion of dwarf galaxies.
 - Later accretions and disruption of stellar systems with lower star formation rates were important.

Searching for Luminous Core-Collapsed Supernovae in a High-z Proto-Cluster

N. Morimoto, T. Yamada, T. Hayashino, Y. Nakamura, K. Kousai (Tohoku Univ.) and Y.Matsuda (Durham Univ.)

It has been revealed that SSA22 region at redshift z=3.1 is an extremely high density region of LAEs (Lyman α emitters). In addition to this, star formation in this region is supposed to be biased to very large mass from the evidence such as a large number of large LyA EW objects. Therefore core- collapsed supernovae (SNe) of massive stars are expected to occur frequently in this region. While SNe has not been detected beyond z>2 to date, it is very important not only to detect core-collapsed SNe at z=3.1 but also to obtain any useful information to understand star formation in early universe. So we estimated the expectation of observing core-collapsed SNe in SSA22 region at z=3.1 based on our own sample of star-forming galaxies, and carried out preliminary observations and data analysis to search for variable objects to investigate the detectability of core-collapsed SNe at high-z by searching for variability of LAEs in this region.



(And there was an unidentified variable object candidate showing SN-like variance.)

Here is possibility of detecting high-z SNe.

+ Multi-epoch photometric data and spectra can prove the origins of the variability.

A study of light curves from rapidly rotating neutron stars

Astronomy, D3, Kazutoshi Numata



around the spin axis

The hot spot near spin axis can drift !

CMB bispectrum from the second-order cosmological perturbations

• Daisuke Nitta (Astronomy D3), Eiichiro Komatsu (the university of Texas at Austin), Nicola Bartrolo (Universita di Padova), Sabino Matarrese (Universita di Padova), Antonio Riotto (CERN).



OResults



OSummary

•We obtain the general formula of the CMB angular averaged bispectrum.

ary 18-19, 2010 @Toho

•The bispectrum has maximum signal in the squeezed triangles, similar to the local-type primordial bispectrum. •However, detailed calculations show that their shapes are sufficiently different

•We conclude that the product of the first-order terms may be safely ignored in analysis of the future CMB experiments.



GENERAL PROPERTIES OF NON-RADIAL PULSATIONS (NRP): EFFECTS OF ROTATION

A P R I L I A ASTRONOMY, D3, TOHOKU UNIVERSITY

NTRODUCTION

A star's pulsation is NRP if the oscillation is in such a way as to deviate from its spherical shape. NRP is characterized by n(number of nodes of eigenfunctions between the center and surface of star), ℓ (number of surface nodal lines), and m(number of such lines that pass through the star's rotation axis), which correspond to the spherical harmonics $Y_{\ell}^{m}(\theta, \varphi)$. In the observer's co-rotating frame, waves with m < 0, propagating in the same direction to star's rotation, are called *prograde* modes; those with m > 0, traveling in the opposite direction to rotation, are called *retrograde* modes.

Here we analyze the NRP of a $4M_{\odot}$ star to see the effects of rotation to the stability of gravity modes (g-modes) at the zero-age main sequence evolution stage. The rotation effects are in the form of the Coriolis force, as the first effect of rotation, and the Centrifugal force (hence the rotational equilibrium deformation) as the second effects of rotation. The $4M_{\odot}$ itself is an SPB (Slowly Pulsating B-type star) which shows NRP of high-order (n) g-modes.

CALCULATION METHOD

By employing a series expansion in terms of the $Y_{\ell}^{m}(\theta, \varphi)$. Modes with $\ell = |m|+2(j-1)$ are defined as even modes, while $\ell = |m|+2j-1$ are odd modes with j = 1, 2, ...

$$\begin{split} \xi_{a} &= a \sum_{\ell \geq |m|} S_{\ell}(a) Y_{\ell}^{m}(\theta, \varphi) e^{i\alpha} \\ \xi_{\theta} &= a \sum_{\ell \geq |m|} \left[H_{\ell}(a) \frac{\partial}{\partial \theta} Y_{\ell}^{m}(\theta, \varphi) + T_{\ell'}(a) \frac{1}{\sin \theta} \frac{\partial}{\partial \varphi} Y_{\ell'}^{m}(\theta, \varphi) \right] e^{i\alpha} \\ \xi_{\varphi} &= a \sum_{\ell \geq |m|} \left[H_{\ell}(a) \frac{1}{\sin \theta} \frac{\partial}{\partial \varphi} Y_{\ell}^{m}(\theta, \varphi) - T_{\ell'}(a) \frac{\partial}{\partial \theta} Y_{\ell'}^{m}(\theta, \varphi) \right] e^{i\alpha} \end{split}$$

RESULTS & CONCLUSIONS

Slow Rotation

The complex eigenfrequency ω is approximated by

 $\omega = \omega_0 + mC_1\Omega + C_2\Omega^2$

where ω_0 is eigenfrequency of mode at $\Omega = 0$. Modes with $\omega_I < 0$ are unstable. For $\Omega > 0$, if $C_{1I} > 0$, the retrograde modes are stabilized and prograde modes are destabilized by rotation; if $C_{2I} > 0$, both retrograde and prograde modes are stabilized by rotation.



The first and second order effects of rotation work both for destabilizing and stabilizing the oscillation modes, but only a few modes can change their stabilities as a result of the second order effects.



If without rotation, most of the high-order g-modes are stable modes, while the unstable modes are in the certain range of ω of the low-order g-modes.

When rotation is included, C_{1l} of the high-order *g*-modes are negative, mean that the retrograde modes are destabilized and prograde modes are stabilized by rotation. Unstable modes are only found at $\ell < 10$.

Rapid Rotation



As Ω increases, many avoided crossings occur between the modes as can be seen in the behavior of ω_R . At the crossings, the mode properties are exchanged between the modes, as is shown by the behavior of ω_l . Only modes with not so different ℓ -values can cross (e.g. $\ell = 1 \& \ell = 3$) because the different of their ω_l s are not very large.



1. Preliminaries

Fuzzy logic is understood a logic with comparative notion of truth, the standard set of truth values is [0,1]. A t-norm is a binary operation * in [0,1] which is commutative, associative, non-decreasing and satisfying 0 * x = 0 and 1 * x = x for each x. Each continuous t-norm * determines uniquely its reduum \Rightarrow satisfying, $\forall x, y, z$, the condition $z \le x \Rightarrow y$ iff $x * z \le y$ (*)

Three most important continuous t-norms are

name	x + y	$x \Rightarrow y \text{ for } x > y$	$\neg x \Rightarrow 0$
Ł	$\max(0, x + y - 1)$	1-x+y	1 - x
G	$\min(x, y)$	y	$\int 1$ for $x = 0$
П	$x \cdot y$	y/x	0 for x > 0

2. Basic Fuzzy Logic and three stronger systems

The following formulas are axioms of the logic BL (A1) $(\varphi \to \psi) \to ((\psi \to \chi) \to (\varphi \to \chi))$ (A2) $(\varphi \& \psi) \rightarrow \varphi$

(A3) $(\varphi \& \psi) \rightarrow (\psi \& \varphi)$ $(A4) \left(\varphi \& (\varphi \to \psi) \right) \to \left(\psi \& (\psi \to \varphi) \right)$ $(A5a)\left(\varphi \to (\psi \to \chi)\right) \to \left((\varphi \& \psi) \to \chi\right)$ $(A5b) ((\varphi \& \psi) \to \chi) \to (\varphi \to (\psi \to \chi))$ $(A6)\left((\varphi \to \psi) \to \chi\right) \to \left(\left((\psi \to \varphi) \to \chi\right) \to \chi\right)$ $(A7)\overline{0} \rightarrow \varphi$

Łukasiewicz logic Ł is BL plus the axiom of double negation $\neg \neg \varphi \rightarrow \varphi$; Gödel logic G is BL plus the axiom of idempotence of conjunction $\varphi \rightarrow (\varphi \& \varphi)$; Product logic Π , originally introduced in [1], is defined in [2] as the extension of BL by axiom $\neg\neg\varphi \rightarrow ((\varphi \rightarrow (\varphi \& \psi)) \rightarrow (\psi \& \neg \neg \psi)).$

General algebras of truth functions for BL are called BLalgebras. A BL-algebras is a structure $\mathbf{L} = (L, \cap, \cup, *, \Rightarrow, 0, 1)$ with several properties. The axioms A1-A7 are Ltautologies for any BL-algebra L. If * is a continuous t-norm and \Rightarrow its residuum then ([0,1], min, max,*, \Rightarrow , 0,1) is a particular linearly ordered BL-algebra [0,1], given by * (called t-algebras or standard BL-algebras). The standard MV-algebra (G-algebra, Π -algebra) is just $[0,1]_*$ where * is \pounds (G, Π) t-norm. Let now C stand for BL, \pounds , G, Π , let Calgebras be BL, MV, G and II-algebras respectively.

General completeness theorem. For each formula φ , following are equivalent: (i) φ is provable in C, (ii) for each C-algebra L, φ is an L-tautology, (iii) for each linearly ordered C-algebra L, φ is an L-tautology

3. Computational complexity

Follows are survey on fuzzy propositional calculi. We consider four sets of formulas as:

 $SAT_1^{\mathcal{C}} = \{\varphi | \text{for some L-evaluation } e, e_{\mathcal{C}}(\varphi) = 1\}$ $SAT_{pos}^{\mathcal{C}} = \{\varphi | \text{for some } \mathbf{L}\text{-evaluation } e, e_{\mathcal{C}}(\varphi) > 1\}$ $TAUT_1^{\mathcal{C}} = \{\varphi | \text{for each L-evaluation } e, e_{\mathcal{C}}(\varphi) = 1\}$ $TAUT_{pos}^{c} = \{\varphi | \text{for each } \mathbf{L}\text{-evaluation } e, e_{\mathcal{C}}(\varphi) > 1\}$

Summarizing, we have the following: (1) $SAT_1^G = SAT_{pos}^G = SAT_1^{\Pi} = SAT_{pos}^{\Pi} = SAT_1^{Bool}$ is NP-

(2) $TAUT_{pos}^{G} = TAUT_{pos}^{\Pi} = TAUT^{Bool}$ is co-NP-complete.

(3) $TAUT_1^G$, $TAUT_1^\Pi$, $TAUT_1^{Bool}$ are pairwise distinct and all co-NP-complete.

(4) $SAT_{pos}^{k} \supset SAT_{1}^{k} \supset SAT^{Bool}$, all NP-complete. (5) $TAUT_{1}^{k} \subset TAUT_{pos}^{k} \subset TAUT^{Bool}$, are co-NP-complete.

3. Predicate Fuzzy Logic and three stronger systems A standard interpretation of the language is a structure $\mathbf{M} = (M, (r_P)_P, (m_c)_c)$ where each $m_c \in M$ and for each nary predicate P, $r_P: M^n \to [0,1]$. An L-interpretation over BL is a structure as above but r_P is an L-fuzzy relation, i.e. $r_P: M^n \to L$. The truth value $\|\varphi\|_{M,\nu}^L$ of a formula φ , is given by BL-algebra L, an L-interpretation M and an evaluation vassigning to each variable x an element $v(x) \in M$.

 $\|P(u_1, ..., u_n)\|_{M,v}^{L} = r_P(v(u_1), ..., v(u_n)),$
$$\begin{split} \|\varphi \to \psi\|_{\mathbf{M},v}^{\mathbf{L}} &= \|\varphi\|_{\mathbf{M},v}^{\mathbf{L}} \Rightarrow \|\psi\|_{\mathbf{M},v}^{\mathbf{L}}, \text{ similarly &,*,} \\ \|(\forall x)\varphi\|_{\mathbf{M},v}^{\mathbf{L}} &= \inf\{\|\varphi\|_{\mathbf{M},v'}^{\mathbf{L}} | v' \equiv_{x} v\}, \text{ similarly } \exists, \sup \end{split}$$

where $v' \equiv_x v$ means that v'(y) = v(y) for all variables y except possibly x. The L-interpretation M is safe if $\|\varphi\|_{M,v}^L$ is defined for all φ and v. Let $C \forall$ be BL \forall , $L \forall$, $G \forall$, $\Pi \forall$.

General completeness theorem. For each theory T and formula φ , $T \vdash_{C \forall} \varphi$ iff for each linearly ordered C-algebra L and each L-model M of T, φ is L-true in M. In particular, $\mathcal{C} \forall \vdash \varphi$ iff for each \mathcal{C} -algebra L, φ is an L-tautology.

4. Arithmetical complexity

We interested in arithmetical complexity of sets of predicate tautologies and sets of satisfiable formulas.

 $TAUT^{\mathcal{C}} = \{\varphi | \text{for all standard } \mathcal{C} \text{ algebras } \mathbf{L} \text{ and } \mathbf{L}\text{-safe } \mathbf{M}, \|\varphi\|_{\mathbf{M}}^{\mathbf{L}} = 1\}$ $genTAUT^{C} = \{\varphi | \text{for all } C \text{ algebras } L \text{ and } L\text{-safe } M, \|\varphi\|_{M}^{L} = 1\}$ $SAT^{\mathcal{C}} = \{\varphi | \text{for some standard } \mathcal{C} \text{ algebras } \mathbf{L} \text{ and } \mathbf{L}\text{-safe } \mathbf{M}, \|\varphi\|_{\mathbf{M}}^{\mathbf{L}} = 1\}$ $genSAT^{C} = \{\varphi | \text{for some } C \text{ algebras } L \text{ and } L\text{-safe } M, \|\varphi\|_{M}^{L} = 1\}$

The results are summarized in the following table.

	Provable = gen.TAUT	stand. TAUT	gen. SAT = consistent	stand. SAT
BL∀	Σ_1 -compl.	Π ₂ -hard	Π_1 -compl.	Not Arithm.
Ł∀	Σ_1 -compl.	Π_2 -compl.	Π_1 -compl.	Π_1 -compl.
$G \forall$	Σ_1 -compl.	Σ_1 -compl.	Π_1 -compl.	Π_1 -compl.
П∀	Σ_1 -compl.	Π ₂ -hard	Π_1 -compl.	Not Arithm.

Since by completeness theorem general tautologies coincide with provable formulas, $genTAUT^{c}$ is evidently Σ_1 ; Σ_1 completeness is proved in [3]. Similarly, general satisfiability coincide with consistence; thus $genSAT^{C}$ is Π_1 -complete. $TAUT^{\downarrow\forall}$ being Π_2 -complete is a classical result of Ragaz. The fact that $SAT^{\Pi \forall}$ is not arithmetical is proved in [3]. We found there are open problems concerning the arithmetical complexity on monadic logic [4].

Reference

- [1] Hájek, P.; Godo, L.; Esteva, F., A complete many-valued logic with product conjunction, Arch. Math. Logic, 35 (1996), 191-208.
- [2] Hájek, P., Metamathematics of fuzzy logic. Kluwer 1998.
- [3] Hájek, P., Fuzzy logic and arithmetical hierarchy III, Studia logica, 68 (2001), 129-142.
- [4] Hájek, P., Monadic fuzzy predicate logic, Studia logica, 71 (2002), 165-175

Some space-time integrability estimates of the solution for heat equations in two dimension

Norisuke Ioku Mathematical institute, Tohoku University, Japan sa6m02@math.tohoku.ac.jp

Heat equation in 2-dimension

 $I=(0,T),\ \Omega\subset \mathrm{R}^2:$ bounded domain

$${
m (H)} \quad egin{cases} \partial_t u - \Delta u = f, & {
m in} \quad I imes \Omega, \ u = 0, & {
m on} \ I imes \partial \Omega, \ u(0) = u_0, & {
m on} \ \{0\} imes \Omega. \end{cases}$$

Lorentz-Zygmund space

 $f\in L^p\log L^q(\Omega)\iff \int_{B_R} \left\{ \left(\lograc{R}{|x|}
ight)^q f^\sharp(x)
ight\}^p dx <\infty$

 f^{\sharp} is a symmetric-decreasing rearrangement of f. $|B_R| = |\Omega|.$

Figure of f^{\sharp}



 $\begin{array}{l} \underline{\text{Known Results}} \\ \hline \text{Maximal regularity} \\ \underline{\text{Let } 1 < p, q < \infty. \ \forall f \in L^q(I; L^p(\Omega)),} \\ \exists u \in W^{1,q}(I; L^p(\Omega)) \cap L^q(I; W^{2,p}(\Omega)) \text{ s.t.} \\ \|\partial_t u\|_{L^q(I; L^p)} + \|\Delta u\|_{L^q(I; L^p)} \leq C \|f\|_{L^q(I; L^p)} \end{array}$

PROBLEM

How about the regularity of a weak solution u for $f \in L^q(I; L\log L^p(\Omega))$?

Main Theorem

Definition (Orlicz space) –

 $\exp L^p(\Omega) := \left\{ u \in L^1_{\operatorname{loc}}; \|u\|_{\exp L^p(\Omega)} < \infty
ight\},$ where $\|u\|_{\exp L^p(\Omega)}$ is a Luxemburg norm.

 $egin{aligned} & \overline{ ext{Theorem1}} \ \overline{ ext{Let} \ f \ \in \ L^q \left(I; L^1 {\log L^p}
ight), \ u_0 \ \in \ L^q {\log L^{p-1}}} \ (0$

$$egin{aligned} \overline{ ext{Let}\; f \in L^1\left(I; L^1 ext{log}\; L^p
ight), \; u_0} \in L^1 ext{log}\; L^p \ (0$$

Remark

Theorem 1,2 can not be obtained by a simple application of the estimate of the maximal regularity.

$$egin{aligned} \overline{ ext{Let}\ p \geq 1,\ u \in \exp L^p(\Omega)}. \ \overline{ ext{Then}} \ \|u\|_{\exp L^p(\Omega)} \simeq \sup_{0 < r < R} rac{u^{\sharp}(r)}{\left(\log rac{R}{r} + 1
ight)^{1/p}}. \end{aligned}$$

Proof of Theorem 1 (Outline)

• Step 1 Use the differential inequality of u^{\sharp}

$$\partial_t u^{\sharp\sharp} - 2 rac{\partial_r u^\sharp}{r} \leq f^{\sharp\sharp},$$

where $u^{\sharp\sharp}(r) := rac{2}{r^2} \int_{B_r} u^{\sharp}(x) dx.$ • Step2 Test by $r \left(\log rac{eR}{r}
ight)^{(p-1)q} u^{\sharp\sharp^{q-1}}.$

Well-posedness for Navier-Stokes equations in modulation spaces with negative derivative indices

Tsukasa Iwabuchi (Tohoku University)

Abstract. The Cauchy problems for Navier-Stokes equations are studied in modulation spaces $M_{q,\sigma}^s(\mathbb{R}^n)$. Our aim is to reveal the conditions of s, q and σ of $M^s_{q,\sigma}(\mathbb{R}^n)$ for the existence of local and global solutions for initial data in $M^s_{q,\sigma}(\mathbb{R}^n)$.

1 Modulation spaces $M^s_{q,\sigma}(\mathbb{R}^n)$

In this section, we define Modulation spaces $M^s_{q,\sigma}(\mathbb{R}^n)$ and introduce their properties.

Definition (Modulation spaces). $\{\varphi_k\}_{k \in \mathbb{Z}^n}$: partition of unity with supp $\varphi_k \subset \{ \xi \in \mathbb{R}^n \mid |\xi - k| \le \sqrt{n} \}$. For example, in the case of \mathbb{R}^2 ,



For $s \in \mathbb{R}, 1 \leq q, \sigma \leq \infty$,

$$M_{q,\sigma}^{s}(\mathbb{R}^{n}) := \left\{ f \in \mathcal{S}'(\mathbb{R}^{n}) \mid \left(\sum_{k \in \mathbb{Z}^{n}} (1+|k|)^{s\sigma} || \mathcal{F}^{-1}\varphi_{k}\mathcal{F}f||_{L^{q}(\mathbb{R}^{n})}^{\sigma} \right)^{\frac{1}{\sigma}} < \infty \right\}$$

Proposition. [Feichtinger (1983), Toft (2004), Wang, Hudzik (2007)]

(i) $M^s_{q,\sigma}(\mathbb{R}^n) \hookrightarrow M^{s_0}_{q_0,\sigma_0}(\mathbb{R}^n)$ if $s \ge s_0, q \le q_0, \sigma \le \sigma_0$.

(ii) $M^0_{q,\min\{q,q'\}}(\mathbb{R}^n) \hookrightarrow L^q(\mathbb{R}^n) \hookrightarrow M^0_{q,\max\{q,q'\}}(\mathbb{R}^n)$ if $\frac{1}{q} + \frac{1}{q'} = 1$.

Remark. In Besov spaces $B^s_{q,\sigma}(\mathbb{R}^n)$, we use dyadic decomposition instead of square decomposition. Then, we have

 $B^s_{q,\sigma}(\mathbb{R}^n) \nleftrightarrow B^s_{q_0,\sigma}(\mathbb{R}^n)$ if $1 \le q < q_0 \le \infty$.



$\mathbf{2}$ Results

(NS)
$$\begin{cases} \partial_t u - \Delta u + u \cdot \nabla u + \nabla p = 0 & \text{for } (t, x) \in (0, \infty) \times \mathbb{R}^n, \\ \text{div } u = 0 & \text{for } (t, x) \in (0, \infty) \times \mathbb{R}^n, \\ u(0, x) = u_0(x) & \text{for } x \in \mathbb{R}^n. \end{cases}$$

	Local solutions	Global solutions
	Local solutions	Global solutions
T. Kato (1984)	$u_0 \in L^n(\mathbb{R}^n)$	small $u_0 \in L^n(\mathbb{R}^n)$
F. Planchon (1998)		small $u_0 \in \dot{B}_{q,\infty}^{-1+\frac{n}{q}}(\mathbb{R}^n)$
H. Koch, D. Tataru (2001)	$u_0 \in vmo^{-1}$	small $u_0 \in BMO^{-1}$
B. Wang, L. Zhao,B. Guo (2006)	$u_0 \in M^0_{2,1}(\mathbb{R}^n)$	
I. (2008)	$u_0 \in M^0_{\infty,\frac{n}{n-1}}(\mathbb{R}^n)$	small $u_0 \in M^0_{n,\frac{n}{n-1}}(\mathbb{R}^n)$

$$PM^s_{q,\sigma}(\mathbb{R}^n) := \{ f \in [M^s_{q,\sigma}(\mathbb{R}^n)]^n \mid \text{div } f = 0 \text{ in } \mathcal{S}'(\mathbb{R}^n) \}$$

Main Theorem. Let n, s, q and σ satisfy

$$n \ge 2$$
, $s \ge -1 + \frac{n(\sigma - 1)}{\sigma}$, $1 \le q \le \infty$.

Then, for any $u_0 \in PM^s_{q,\sigma}(\mathbb{R}^n)$ there exists T > 0 such that (NS) has a unique solution u in some subspace of $\left[C([0,T], M^s_{a,\sigma}(\mathbb{R}^n))\right]^n$.

If $q \leq n$ and u_0 is sufficiently small, the solution exists globally in time.

Remark. (i) If s = 0, then we have $\sigma \le n/(n-1)$. Therefore, this Theorem is an extension of our previous result.

(ii) Since $s \ge -1 + n(\sigma - 1)/\sigma$, we have $s \ge -1$. If s < -1, (NS) is illposedness in $M^s_{2,\sigma}(\mathbb{R}^n)$ $(1 \leq \sigma < \infty)$. Therefore, the derivative index s = -1 is optimal to obtain well-posedness. The way of the proof is applying the argument of I. Bejenaru, T. Tao (2006). They studied the nonlinear Schrödinger equation $i\partial_t u + \Delta u = u^2$ in one space dimension, and we apply their argument to the Navier-Stokes equation.

(iii) If s = -1, our initial data are included in the result by Koch and Tataru (2001). In fact, we have the following continuous embeddings

$$M_{q,1}^{-1}(\mathbb{R}^n) \hookrightarrow vmo^{-1} \quad \text{if } 1 \le q \le \infty,$$

$$M_{a,1}^{-1}(\mathbb{R}^n) \hookrightarrow BMO^{-1} \quad \text{if } 1 \le q \le n.$$

The other cases are not clear in general by $M^s_{q,\sigma}(\mathbb{R}^n) \not\hookrightarrow M^{-1}_{q,1}(\mathbb{R}^n)$ in the case of $1 \le q \le \infty$ and $s = -1 + n(\sigma - 1)/\sigma > -1$.

Proposition to prove theorems 3

We introduce the proposition to prove our theorem in the case of $-1 \leq s \leq 0$. The way of the proof is applying Banach's fixed point theorem with the following propositions.

Let $s, \tilde{s} \in \mathbb{R}, 1 \leq q, r, \sigma, \nu \leq \infty, e^{t\Delta} := \mathcal{F}^{-1} e^{-t|\xi|^2} \mathcal{F}$, we have the following estimates.

(i) $\|e^{t\Delta}u\|_{M^s_{q,\sigma}} \le C(1+t)^{-\frac{n}{2}(\frac{1}{r}-\frac{1}{q})} \|u\|_{M^s_{r,\sigma}}$ if $q \ge r$.

(ii)
$$\|e^{t\Delta}u\|_{M^s_{q,\sigma}} \leq C\left(1+t^{-\frac{s-s}{2}}\right)\|u\|_{M^s_{q,\sigma}}$$
 if $s \geq \widetilde{s}$.

- (iii) $\|e^{t\Delta}u\|_{M^{s,\sigma}_{q,\sigma}} \leq C\left(1+t^{-\frac{n}{2}(\frac{1}{\sigma}-\frac{1}{\nu})}\right)\|u\|_{M^{s,\sigma}_{q,\nu}}$ if $\sigma \leq \nu$. (iv) [Toft (2004)] If $\frac{1}{q} = \frac{1}{q_1} + \frac{1}{q_2}$ and $\frac{1}{\sigma} = \frac{1}{\sigma_1} + \frac{1}{\sigma_2} 1$, $\|uv\|_{M^{0}_{q,\sigma}} \leq C\|u\|_{M^{0}_{q_1,\sigma_1}}\|v\|_{M^{0}_{q_2,\sigma_2}}.$

(v)
$$\|e^{t\Delta}u\|_{L^{2}(0,T;M_{q,1}^{0})} \leq C\left(1+T^{\frac{1}{2}}\right)\|u\|_{M_{q,1}^{-1}}.$$

(vi) $\left\| \int_{0} \nabla e^{t\Delta} f(\tau) d\tau \right\|_{L^{2}(0,T;M^{0}_{q,1})} \leq C \left(1 + T^{\frac{1}{2}}\right) \|f\|_{L^{1}(0,T;M^{0}_{q,1})}.$ **Remark.** (1) We use the estimates (ii), (iii) and (iv) to obtain local

solutions of our theorem in the case $-1 < s \le 0$, and the case s = -1can be treated with (iv), (v) and (vi). (i) is used for the proof of the existence of global solutions.

(2) Similar estimates to (i) and (ii) are known in Sobolev spaces.

$$\left\|e^{t\Delta}u\right\|_{H^{\bar{s},q}} \le Ct^{-\frac{n}{2}(\frac{1}{r}-\frac{1}{q})} \left(1+t^{-\frac{s-\bar{s}}{2}}\right) \|u\|_{H^{\bar{s},r}} \quad \text{if } s > \tilde{s}, \ q \ge r.$$

The soul conjecture for Riemannian orbifolds

Naoki Oishi (Mathematical Institute Tohoku University)

1 Introduction

In 1994, Perelman solved the Cheeger-Gromoll soul conjecture:

<u>Theorem</u> 1. M^n : non-compact, complete Riemannian manifold with $K(M^n) \ge 0$, $\exists p \in M^n$ s.t. K(p) > 0.

 \Rightarrow The soul *S* of M^n is one point set. Consequently M^n is diffeomorphic to \mathbb{R}^n .

Here the soul S of M^n is a totally geodesic submanifold of M^n , which is constructed by some inductive way. M^n is diffeomorphic to a normal bundle of S.

One of the ideas of proving the theorem is to construct a flat rectangle on S and extend the rectangle along geodesics until it contains p.

By the way, Perelman constructed a soul on an Alexandrov space. But he couldn't generalize the theorem to the Alexandrov space. We show an analogous result to the theorem for Riemannian orbifolds with nonegative curvature, which belong to Alexandrov spaces.

2 Riemannian orbifolds

• M : complete Riemannian manifold.

 $egin{aligned} &\Gamma\subset \mathsf{Isom}(M) ext{ , }\Gamma &\curvearrowright M: ext{ discontinuously.} \ &\pi:M
i x\mapsto [x]_{\Gamma}\in M/\Gamma. \ & ilde{d}([x]_{\Gamma},[y]_{\Gamma}):=d(\pi^{-1}([x]_{\Gamma}),\pi^{-1}([y]_{\Gamma})). \ &\Rightarrow (M/\Gamma, ilde{d}): ext{ metric space.} \end{aligned}$

<u>Definition</u> . $O := (M/\Gamma, \tilde{d})$: (good) Riemannian orbifold.

<u>Definition</u> . $p \in O$ is said to be a positively curved point if $K(\tilde{p}) > 0$ holds for $\forall \tilde{p} \in \pi^{-1}(p)$. • $\tilde{a} \in M$, $\Gamma_{\tilde{a}} := \{g \in \Gamma | g\tilde{a} = \tilde{a}\}.$ $\Gamma_{\tilde{x}} \simeq \Gamma_{\tilde{y}}$ holds for $\forall x \in O, \forall \tilde{x}, \tilde{y} \in \pi^{-1}(x)$.

<u>Definition</u> . $\Gamma_x := \Gamma_{\tilde{x}}$. $x \in O$ is called a singular point if $\Gamma_x \neq \{e\}$.

• $x \in O$, γ : geodesic in O with $\gamma(0) = x$. $c(x) := \sup\{c > 0 | \text{ any } \gamma \text{ with } L(\gamma) < c \text{ can}$ extend to a geodesic $\overline{\gamma}$ with $L(\overline{\gamma}) = c$ and $\overline{\gamma}(0) = x\}$.

 $\underline{\text{Definition}} \ . \ E(A) := \inf\{c(x) \, | \, x \in A\}$ for $A \subset O.$

Remark. E(A) may not be positive. For example, let O be a cone. In this case, the origin o is a singular point, and any geodisic γ can't pass through o; see Figure.



<u>Theorem</u> 2. O : non-compact, complete Riemannian orbifold with $curv(O) \ge 0$, $\exists p \in O$: positively curved point.

Moreover, we assume the followings with respect to the soul S in O:

- (1) $\forall x, y \in S$, $\Gamma_x \simeq \Gamma_y$.
- (2) $d(p, S) \leq E(S)$.
- \Rightarrow S : one point set.

Remark. (1) implies E(S) > 0. Therefore by (2), we can apply the Perelman's method on S.

A new approach to the existence of harmonic maps

OMORI Toshiaki

(Mathematical institute, Tohoku University, Sendai, Japan)

sa6m07@math.tohoku.ac.jp

 $\operatorname{sect}^N \leq \mathbf{0} \Rightarrow \operatorname{ExpHarm}(M, N) : \operatorname{compact}$

Main Theorem
Corollary

What is exponentially harmonic maps ?

Harmonic maps

 $E(u) := \frac{1}{2} \int_{M} |du|^{2} d\mu_{g}$ $\Delta_{g} u + A(u) (\nabla u, \nabla u) = 0$

 ε -exponentially harmonic maps ($\varepsilon > 0$)

$$\mathbb{E}_{\varepsilon}(u) := \int_{M} e^{\varepsilon |\nabla u|^{2}} d\mu_{g}$$
$$\Delta_{g} u + \varepsilon \langle \nabla |\nabla u|^{2}, \nabla u \rangle + A(u) (\nabla u, \nabla u) = 0$$

Remark

 $\begin{array}{l} \ddagger u \in \operatorname{Harm}(\mathbb{T}^2, \mathbb{S}^2) \text{ with } \operatorname{deg}(u) = \pm 1 \\ \text{On the other hand} \\ \hline \mathbf{Fact} [\operatorname{Duc} \& \operatorname{Naito}] \\ \forall \ \mathscr{H} \in [M, N] \\ \exists u_{\varepsilon} \in C^{\infty}(M, N) : \mathbb{E}_{\varepsilon} \text{-minimizer in } \mathscr{H} \end{array}$

Proof of Main Theorem

Stability of the interface of a Hele-Shaw flow with two injection points

Michiaki ONODERA (Mathematics, DC3, Tohoku University)

sa5m08@math.tohoku.ac.jp

The 2nd GCOE International Symposium "Weaving Science Web beyond Particle-Matter Hierarchy"

February 18–19, 2010

1 Hele-Shaw flows with free boundaries

Hele-Shaw flows are fluid flows in an experimental device which consists of two closely-placed parallel plates. Since the gap between two plates is sufficiently narrow, one can regard them as two-dimensional flows. We consider a Hele-Shaw flow produced by the injection of fluid from two points. Let the fluid initially occupy a bounded domain $\Omega(t_0) \subset \mathbb{C}$ at the initial time $t = t_0$, and assume that $\pm i \in \Omega(t_0)$. From each point $\pm i$, additional fluid is injected at the respective rate $\alpha_1, \alpha_2 > 0$ per unit time. The fluid domain at time $t > t_0$ is denoted by $\Omega(t)$ and its boundary by $\partial\Omega(t)$. We are interested in the behavior of the interface of the Hele-Shaw flow as time tends to infinity. In particular, we consider the stability of the interface under disturbance.



One of the significant features of the Hele-Shaw flow is that the flow is characterized as a two-dimensional potential flow with the potential being its pressure p. Let p = p(z, t) be the pressure of the fluid at position $z \in \mathbb{C}$ and time t > 0. Then, p is assumed to satisfy the following equation and boundary conditions:

ſ	(1)	$-\Delta p = \alpha_1 \delta_i + \alpha_2 \delta_{-i}$	in $\Omega(t)$
ł	(2)	p = 0	on $\partial \Omega(t$
l	(3)	$-\partial_n p = v_n$	on $\partial \Omega(t)$

From (1) and (2), for each time $t > t_0$, the function p can be represented by

 $p(z,t) = \alpha_1 G_{i,\Omega(t)}(z) + \alpha_2 G_{-i,\Omega(t)}(z) \quad \text{for } z \in \Omega(t),$

where $G_{\pm i,\Omega(t)}$ are the Green's functions of $\Omega(t)$ with pole at $\pm i$, respectively. We call a family $\{\Omega(t)\}_{t>t_0}$ of domains a solution of the Hele-Shaw problem if it satisfies (1)–(3) with the function p defined above.

2 Explicit solutions

as

By means of a complex analytic method, we are able to construct an explicit solution $\{\Omega_0(t)\}_{t>t_0}$ of the equation (3). The domain $\Omega_0(t)$ is represented as the image of the unit disk D through a conformal mapping φ_t , i.e., $\Omega_0(t) = \varphi_t(D)$. The mapping φ_t is given by

$$\varphi_t(w) = \varphi(w,t) = \frac{R^4 - 1}{R} \cdot \frac{w - i\eta}{w^2 + R^2} + i\eta R,$$

where R=R(t) and $\eta=\eta(t)$ are smooth functions of t and they have the following asymptotic forms:

$$R(t) = \sqrt{\frac{\alpha_1 + \alpha_2}{\pi}} \cdot t^{1/2} + O(t^{-1/2}), \quad \eta(t) = \frac{\alpha_2 - \alpha_1}{\sqrt{\pi(\alpha_1 + \alpha_2)}} \cdot t^{1/2} + O(t^{-1/2}),$$

$$t \to \infty.$$

3 Derivation of an evolution equation

Our objective is the stability of the solution $\{\Omega_0(t)\}$ of the equation (3) under disturbance of the initial domain $\Omega_0(t_0)$. We investigate the behavior of a solution starting from a domain $\Omega(t_0)$ which is sufficiently close to $\Omega_0(t_0)$ for a large positive number t_0 . In this section we derive an evolution equation from the free boundary problem (3).

For a continuous function $r = r(\xi)$ on the unit circle \mathbb{S}^1 , we define

$$D_r := \{ w \in \mathbb{C} \setminus \{0\} \mid |w| < 1 + r(w/|w|) \} \cup \{0\}.$$

1

Let us assume that a family $\{\Omega(t)\}_{t \ge t_0}$ of domains satisfies the equation (3), and that each domain $\Omega(t)$ is represented by $\Omega(t) = \varphi_t(D_{r(\cdot,t)})$ for some time-dependent smooth function $r = r(\xi, t), \ \xi \in \mathbb{S}^1, \ t \ge t_0$. Then, after some computations we see that the function r satisfies the following equation:

$$\partial_t r(\xi,t) = -\frac{\left\langle {}^t \nabla G(w_r(\xi),t) , n_{\partial D_r}(w_r(\xi)) \right\rangle + \left\langle {}^t \nabla H_{r(\cdot,t),t}(w_r(\xi)) , n_{\partial D_r}(w_r(\xi)) \right\rangle}{\det D_w \varphi(w_r(\xi),t) \cdot \left\langle \xi , n_{\partial D_r}(w_r(\xi)) \right\rangle} - \frac{\left\langle \partial_t \varphi(w_r(\xi),t) , \left(D_w \varphi(w_r(\xi),t) \right) n_{\partial D_r}(w_r(\xi)) \right\rangle}{\det D_w \varphi(w_r(\xi),t) \cdot \left\langle \xi , n_{\partial D_r}(w_r(\xi)) \right\rangle},$$
(4)

where $w_r(\xi) := (1 + r(\xi, t))\xi$, and $\langle \cdot, \cdot \rangle$ denotes the usual inner product in \mathbb{R}^2 , and $n_{\partial D_r}(w)$ is the unit outward normal vector to ∂D_r at point w. The function G = G(w, t) is defined by

 $G(w,t) := \alpha_1 G_{\varphi_t^{-1}(i),D}(w) + \alpha_2 G_{\varphi_t^{-1}(-i),D}(w),$

and $H_{r,t} = H_{r,t}(w)$ is a unique solution to the Dirichlet problem

$$\begin{cases} \Delta H_{r,t}(w) = 0 & \text{for } w \in D_r, \\ H_{r,t}(w) = -G(w,t) & \text{for } w \in \partial D_r \end{cases}$$

Note that $r \equiv 0$ corresponds to the family $\{\Omega_0(t)\}_{t>t_0}$, and hence $r \equiv 0$ is a solution to the equation (4).



4 Main result

To state our main result, we recall some function spaces. The Hölder space $C^{k,\gamma}(\mathbb{S}^1)$ is a Banach space consisting of those functions r that are k-times continuously differentiable with their k-th order derivatives $r^{(k)}$ satisfying

$$|r^{(k)}(\xi_1) - r^{(k)}(\xi_2)| \le C|\xi_1 - \xi_2|^{2}$$

for some constant C independent of $\xi_1, \xi_2 \in \mathbb{S}^1$. Then, the little Hölder space $h^{k,\gamma}(\mathbb{S}^1)$ is defined to be the closure of the subspace $C^{\infty}(\mathbb{S}^1)$ of infinitely many differentiable functions, in the topology of $C^{k,\gamma}(\mathbb{S}^1)$.

Theorem. Suppose t_0 is sufficiently large. For $\varepsilon > 0$, there are $\delta, M > 0$ s.t. if $||r_0||_{h^{2,\gamma}(\mathbb{S}^1)} < \delta$, then there exists a unique solution $r \in C([t_0,\infty); h^{2,\gamma}(\mathbb{S}^1)) \cap C^1([t_0,\infty); h^{1,\gamma}(\mathbb{S}^1))$ satisfying

 $\|r(t)\|_{h^{2,\gamma}(\mathbb{S}^1)} + t\|r'(t)\|_{h^{1,\gamma}(\mathbb{S}^1)} \le Mt^{-1+\varepsilon}\|r_0\|_{h^{2,\gamma}(\mathbb{S}^1)}.$

Therefore, the solution $\{\Omega_0(t)\}_{t>t_0}$ is stable under disturbance, and the disturbance decays with algebraic order as time goes to infinity.

5 Outline of the proof

We treat the equation (4) as an evolution equation on $h^{1,\gamma}(\mathbb{S}^1)$ as follows:

$$\begin{cases} r' = \mathcal{F}(r,t), \quad t \ge t_0, \\ r(t_0) = r_0 \in h^{2,\gamma}(\mathbb{S}^1), \end{cases}$$

where $\mathcal{F}: h^{2,\gamma}(\mathbb{S}^1) \times [t_0,\infty) \to h^{1,\gamma}(\mathbb{S}^1)$, and $\mathcal{F}(0,t) = 0$.

 $\boxed{\text{Lemma 1.}} t \cdot D_r \mathcal{F}(r,t) \to \mathcal{A} \text{ in } \mathcal{L}(h^{2,\gamma}(\mathbb{S}^1),h^{1,\gamma}(\mathbb{S}^1)) \text{ as } t \to \infty \text{ and } r \to 0.$

Lemma 2. \mathcal{A} : sectorial in $h^{1,\gamma}(\mathbb{S}^1)$, and $\sup\{\operatorname{Re}\lambda \mid \lambda \in \sigma(\mathcal{A})\} = -1$. ('97 Escher & Simonett, '09 Vondenhoff)

Set
$$\tau := \log(t/t_0)$$
 and $\tilde{r}(\tau) := r(t)$, then
 $\tilde{r}'(\tau) = t \cdot r'(t) = t \cdot \mathcal{F}(r, t) = \mathcal{A}\tilde{r} + \mathcal{G}(\tilde{r}, \tau),$ (5)
where $\mathcal{G}(0, \tau) = 0$, $D_{\tilde{r}}\mathcal{G}(\tilde{r}, \tau) \to 0$ as $t \to \infty$ and $\tilde{r} \to 0$.

Lemma 3.] (Linearized stability for asymptotically autonomous equations) $\|r_0\|_{h^{2,\gamma}(\mathbb{S}^1)} < \delta \Rightarrow \exists \tilde{r}:$ solution of (5) satisfying

 $\|\tilde{r}(\tau)\|_{h^{2,\gamma}(\mathbb{S}^1)} + \|\tilde{r}'(\tau)\|_{h^{1,\gamma}(\mathbb{S}^1)} \le M' e^{(-1+\varepsilon)\tau} \|r_0\|_{h^{2,\gamma}(\mathbb{S}^1)}.$

• Maximal regularity property of $h^{k,\gamma}(\mathbb{S}^1)$ ('79 Da Prato & Grisvard) $\tilde{r}_j \in C([0,\infty); h^{2,\gamma}(\mathbb{S}^1)) \Rightarrow \mathcal{G}(\tilde{r}_j, \tau) \in C([0,\infty); h^{1,\gamma}(\mathbb{S}^1))$ $\Rightarrow \tilde{r}_{j+1} \in C([0,\infty); h^{2,\gamma}(\mathbb{S}^1))$ (The case where $\mathcal{G} = \mathcal{G}(\tilde{r})$ and $D_{\tilde{r}}\mathcal{G}(0) = 0$ is already known.)

Turning back to the original time variable t, we see that $\|r(t)\|_{h^{2,\gamma}(\mathbb{S}^1)} + t\|r'(t)\|_{h^{1,\gamma}(\mathbb{S}^1)} \leq Mt^{-1+\varepsilon}\|r_0\|_{h^{2,\gamma}(\mathbb{S}^1)}.$

The isomorphism between motivic cohomology and K-groups for equi-characteristic regular local rings.

Yuki Kato sa4d01@math.tohoku.ac.jp

Mathematical Institute, Tohoku University, Aoba, Sendai, JAPAN, 980-8578.

Sendai, Feb. 18 - Feb. 19, 2010

1 Notations.

- R := a regular local ring of equi-characteristic.
- $CH^r(X, n) =$ the higher Chow groups for arbitrary integral scheme X.
- $K_n(X)$ = the algebraic K-groups for arbitrally scheme X.
- $\mathcal{Z}_{eq}(X/S, r) := \mathbb{Z} [Z \subset X \mid X \to S : \text{equi-dimensional of dimension } r]$ (S := a geometrically unibranch scheme, X := a scheme of finite type over S).
- $K^{Q,T}(X)$:= the K-theory spectrum of the scheme $X \times_S T$ with family of supports consisting of all closed subschemes quasi-finite over X

Question 1.1.

 $K_n(X)_{\mathbb{Q}} \stackrel{??}{=} \bigoplus_{r \ge 0} \operatorname{CH}^r(X, n)_{\mathbb{Q}}.$

Theorem 1.2. Let R be an equi-characteristic regular local ring. Then we have an isomorphism

$$\operatorname{cl}: K_n(R)_{\mathbb{Q}} \to \bigoplus_{r \ge 0} \operatorname{CH}^r(R, n)_{\mathbb{Q}}$$

for any $n \ge 0$, where cl is the cycle-class map.

According to Voevodsky–Suslin–Friedlander [4], we can define the higher Chow group by another approach, and this coincides with Bloch's [1] for an arbitrary quasi-projective variety over a field.

2 Friedlander–Suslin–Voevodsky's motivic complex $\mathbb{Q}_{X/S}(r)[\bullet]$.

We always assume that a scheme S is noetherian, reduced and separated of finite dimensional. Write $\Box^n = \mathbb{A}^n_S$. The cubical structure of \Box^n is defined by the faces which are intersections of some of the Cartier divisor $\{t_i = 0\}$ or $\{t_i = 1\}$.

Definition 2.1.

$$\mathbb{Z}_{X/S}(r)[n] = \mathcal{Z}_{eq}(X \times_S \square_S^n \times_S \mathbb{A}^r / X \times_S \square_S^n, 0) /$$
$$\sum_{i=1}^n \pi_i^* (\mathcal{Z}_{eq}(X \times_S \square_S^{n-1} \times_S \mathbb{A}^r / X \times_S \square_S^{n-1}, 0))$$

and write $\mathbb{Q}_{X/S}(r)[n] = \mathbb{Z}_{X/S}(r)[n] \otimes \mathbb{Q}$. We obtain and a cubical complex $\mathbb{Q}_{X/S}(r)[\bullet]$ and define the motivic cohomology

 $\operatorname{CH}_{t}^{r}(X, n)_{\mathbb{Q}} = \mathbf{H}_{t}^{-n}(X, \mathbb{Q}_{X/S}(r)[\bullet]),$

where t is a topology as follows; qfh, ét, Nis or Zar.

3 The proof of the main theorem.

3.1 The construction of the cycle-class map "cl'"

Write $\mathbb{Z}^{\mathcal{Q},T}(X) = \mathcal{Z}_{eq}(X \times_S T/X, 0)$. Then the canonical map

cl:
$$K^{\mathcal{Q},T}(X) \to \pi_0(K^{\mathcal{Q},T}(X)) \to \mathbb{Z}^{\mathcal{Q},T}(X)$$

is given by taking the cycles $[\operatorname{Supp} \mathcal{F}] = \sum_{W} \operatorname{length}(\mathcal{F}_{W})[W]$ of the coherent sheaves \mathcal{F} , where $W \subset \operatorname{Supp} \mathcal{F}$ runs over all irreducible components.

Lemma 3.1. Let S be a regular noetherian scheme, T a smooth scheme of finite type over S and $f: Y \to X$ be a morphism of smooth schemes over S. Then the diagram

$$\begin{array}{c} K^{\mathcal{Q},T}(X) \xrightarrow{\mathrm{cl}} \mathbb{Z}^{\mathcal{Q},T}(Y) \\ \downarrow^{f^*} & \downarrow^{f^*} \\ K^{\mathcal{Q},T}(X) \xrightarrow{\mathrm{cl}} \mathbb{Z}^{\mathcal{Q},T}(Y) \end{array}$$

is commutative.

The main theorem is proved by the following lemmas:

Lemma 3.2. $(T_{\alpha}, f_{\alpha\beta})$ a directed inverse system of smooth schemes over S, X_0 a scheme of finite type over S. Assume $T = \varprojlim T_{\alpha}$ is also regular and noetherian. Then we have the isomorphism:

$$\varinjlim \mathcal{Z}_{eq}(X_{\alpha}/T_{\alpha},0) \otimes \mathbb{Q} \to \mathcal{Z}_{eq}(X/T,0) \otimes \mathbb{Q},$$

where we write $X_{\alpha} = X_0 \times_S T_{\alpha}$ and $X = X_0 \times_S T$.

Lemma 3.3. Let k be a field, X be a smooth scheme and F be a homotopy invariant sheaf of \mathbb{Q} -vector space on $(\mathrm{Sm}/X)_{\mathrm{Zar}}$. Then we have an isomorphisms

$$H^n_{\acute{e}t}(X,F) = H^n_{qfh}(X,F) = H^n_{Nis}(X,F) = H^n_{Zar}(X,F).$$

References

- S. Bloch, Algebraic cycles and higher K-theory, Adv. in Math. 61 (1986), no. 3, 267–304.
- [2] Eric M. Friedlander and A. Suslin, The spectral sequence relating algebraic Ktheory to motivic cohomology, Ann. Sci. École Norm. Sup. (4) 35 (2002), no. 6, 773–875 (English, with English and French summaries).
- [3] Y. Kato, The isomorphism between motivic cohomology and K-groups for equicharacteristic regular local rings., preprint.
- [4] V. Voevodsky, A. Suslin, and Eric M. Friedlander, Cycles, transfers, and motivic homology theories, Annals of Mathematics Studies, vol. 143, Princeton University Press, Princeton, NJ, 2000.

Bifurcations in semilinear elliptic equations on thin domains

Toru Kan

Mathematical Institute, Tohoku University

 $({\rm E-mail:} \verb"sa7m08@math.tohoku.ac.jp")$

1 Thin Domain and Limit Equation

Thin domain

$$\left\{ \begin{array}{ll} \Delta u + f(u,\lambda) = 0 & \mathrm{in} \ {\color{black} Q_{\varepsilon}}, \\ \frac{\partial u}{\partial \nu_{\varepsilon}} = 0 & \mathrm{on} \ \partial Q_{\varepsilon} \end{array} \right.$$

 $\lambda \in \mathbb{R}$: bifurcation parameter

 $\begin{aligned} f(0,\lambda) &\equiv 0\\ Q_{\varepsilon} &:= \{(x,y) \in \mathbb{R}^2; 0 < y < \varepsilon g(x), \ 0 < x < 1\}\\ (\varepsilon > 0 : \text{small parameter}, \ g \in C^2([0,1]), \ g > 0) \end{aligned}$



Change of variables $Q := (0,1) \times (0,1) \ni (x,y) \mapsto (x, \varepsilon g(x)y) \in Q_{\varepsilon}$

$$\longrightarrow (\mathbf{P})_{\varepsilon} \left\{ \begin{array}{l} L_{\varepsilon} u + f(u, \lambda) = 0 \quad \text{in } Q, \\ \frac{\partial u}{\partial \nu_{B_{\varepsilon}}} = B_{\varepsilon} u \cdot \nu = 0 \quad \text{on } \partial Q. \end{array} \right.$$
$$L_{\varepsilon} := \frac{1}{g} \text{div} B_{\varepsilon}, \quad B_{\varepsilon} := \left(\begin{array}{c} g \frac{\partial}{\partial x} - g_{x} y \frac{\partial}{\partial y} \\ -g_{x} y \frac{\partial}{\partial x} + \frac{1}{g} \left\{ \frac{1}{\varepsilon^{2}} + (g_{x} y)^{2} \right\} \frac{\partial}{\partial y} \end{array} \right)$$

 $\frac{\text{Limit equation at } \varepsilon = 0}{\text{The limit equation of } (P)_{\varepsilon}} \text{ at } \varepsilon = 0 \text{ is given by}$

(P)₀
$$\begin{cases} \frac{1}{g(x)}(g(x)v_x)_x + f(v,\lambda) = 0 & \text{in } (0,1), \\ v_x(0) = v_x(1) = 0. \end{cases}$$

2 Problem

How can the bifurcation diagram of $(P)_{\varepsilon}$ be approximated by that of the limit equation $(P)_0$? Do solution branches of $(P)_{\varepsilon}$ persist near those of $(P)_0$ including bifurcation points?

More precisely, we consider the following three situations:

- 1. Bifurcation from the trivial branch
- 2. Regular branch of solutions
- 3. Saddle-node bifurcation



3 <u>Definition</u>

 $\frac{\text{Norm on } H^1(Q)}{1}$

For $u \in H^1(Q)$, we define

$$\|u\| := \left(\|u\|_{H^1(Q)}^2 + \frac{1}{\varepsilon^2} \left\| \frac{\partial u}{\partial y} \right\|_{L^2(Q)}^2 \right)^{\frac{1}{2}}.$$

4 Main Results



 $\begin{aligned} (v_0, \lambda_0) &: \text{ turning point of } (\mathbf{P})_0 \\ \Rightarrow \exists \varepsilon_0 > 0, \exists C > 0 \text{ s.t. } 0 < \forall \varepsilon \leq \varepsilon_0, \exists ! (u_0^{\varepsilon}, \lambda_0^{\varepsilon}) \text{ s.t.} \\ (u_0^{\varepsilon}, \lambda_0^{\varepsilon}) \text{ is a turning point of } (\mathbf{P})_{\varepsilon}, \\ & \left\| u_0^{\varepsilon} - v_0 \right\| + \left| \lambda_0^{\varepsilon} - \lambda_0 \right| \leq C \varepsilon. \end{aligned}$



 λ

On the curvature of the pseudo-volume form defining the Carathéodory measure hyperbolicity

Shin Kikuta

Mathematical Institute, Tohoku university email : sa6m15@math.tohoku.ac.jp

Theorem [Schwarz' lemma].

 $f: \text{holomorphic on } \Delta:=\{z\in\mathbb{C}\,;\, |z|<1\},\ |f|<1,\ f(0)=0\Longrightarrow |f'(0)|\leq 1$

This theorem can be interpreted as follows.

 $|f| < 1 \implies f$: holomorphic map between Δ

$$|f'(0)| \le 1, \ f(0) = 0 \quad \Longrightarrow \quad \frac{4}{(1 - |f(0)|^2)^2} |f'(0)|^2 dx dy \le \frac{4}{(1 - |0|^2)^2} dx dy$$

But

Möbius transformation on $\Delta \Longrightarrow \begin{cases} \cdot \text{Any } z \in \Delta \text{ can be mapped to } 0 \\ \cdot \text{Poincaré volume element } v_1 := \frac{4}{(1-|z|^2)^2} dx dy \\ \text{with constant curvature } K_{v_1} = -1 \text{ is invariant} \end{cases}$

Therefore it follows that

 $f^*v_1 := \frac{4}{(1 - |f(z)|^2)^2} |f'(z)|^2 dx dy \le \frac{4}{(1 - |z|^2)^2} dx dy \quad \text{for } f : \text{holomorphic map between } \Delta$

We consider $|f'(z)|^2 dx dy$ as a change of variables by f.

Geometric interpretation of Schwarz' lemma

Any holomorphic map between Δ decreases volumes measured by v_1

With requiring this volume decreasing property of holomorphic mappings,

 v_1 on $\Delta \xrightarrow{\text{generalization}}$ what on a higher dimensional complex manifold X?

Definition-Theorem [a gerenalization of (Δ, v_1)].

Main theorem [curvature of v_X^C].

 $v_X^C := \sup\{f^*v_1^{(n)}; f \in \operatorname{Hol}(X, \mathbb{B}^n)\}$: Carathéodory pseudo-volume form ("pseudo": it may take 0 at some point $v_1^{(n)}$: Poincaré volume element which generalizes v_1 on Δ to $\mathbb{B}^n := \{z \in \mathbb{C}^n; |z| < 1\}$) \Longrightarrow All holomorphic maps decrease volumes measured with respect to the Carathéodory pseudo-volume form v_X^C

 $K_{v_X^C}(:=$ a curvature of the Carathéodory pseudo-volume form $v_X^C) \leq -1$

GÖDEL'S INCOMPLETENESS THEOREM, RECURSIVELY AXIOMATIZABLE THEORIES AND MEDVEDEV DEGREES OF UNSOLVABILITY

;

Takayuki Kihara

Mathematical Institute, Tohoku University

r

 Gödel's Incompleteness Theorem (1931) Fix a set <i>L</i> of symbols, say &, ∀, =, ∈,, and assume that an <i>L</i>-theory <i>T</i> agrees with the following 3 conditions: <i>T</i> is consistent. <i>T</i> is capable of expressing elementary arithmetic. It is knowable what axioms of <i>T</i> are. (that is, <i>T</i> is recursively axiomatizable.) Then, <i>T</i> must be incomplete, even for sentences generated from just <i>L</i>! The conclusion (of Gödel's Incompleteness Theorem) is unes 	 Goal Our purpose is to classify every possible recursively axiomatizable theories whether or not studied in the history of mankind. To simplify this problem, we focus on the spaces <i>T</i>* of ultrafilters of Lindenbaum algebras of recursively axiomatizable theories <i>T</i> with the Zariski topology. <i>T</i>* is homeomorphic to some Π⁰₁ class in the Cantor space.
nathematical propositions, <i>mathematical thinking is, and must</i> –Emil L. Post (1944), RECURSIVELY E	remain, essentially creative". NUMERABLE SETS OF POSITIVE INTEGERS AND THEIR DECISION PROBLEMS.
Definition A consistent recursively axiomatizable theory T is creative if there is a computable function F such that both affirmation and negation of $F(S)$ are unprovable in S for every consistent recursively axiomatizable extension S of T .	 Definition The Cantor space 2^N is the countably infinite topological product of the discrete space {0, 1}. A Π⁰₁ class is a computably generated closed set in the Cantor space.
 Example Peano arithmetic (PA) is creative via F(S) = Con(S), where Con(S) is a sentence expressing that "S is consistent". Zelmero-Fraenkel set theory (ZFC), a foundational axiomatic system for mathematics, is also creative. 	Theorem (Folklore) • In the Cantor space, there exists a nonempty Π_1^0 class P such that $Q \leq_M P$ for every nonempty Π_1^0 class Q . • Thus, the lattice structure $\mathcal{P}_M = ((\text{nonempty } \Pi_1^0 \text{ classes}) / \equiv_M, \leq_M)$ has the greatest element 1 and the least element 0.
 Theorem of Pour-El and Kripke (1967) There are deduction-preserving recursive isomorphisms between every creative theories. PA is isomorphic to ZFC via a deduction-preserving computable procedure! Definition Y is Medvedev reducible to X (Y ≤_M X) if there is a computable map E + X → X 	ObservationLet T be a recursively axiomatizable theory.• T is essentially incomplete $\iff 0 < \deg_M(T^*) \le 1$.• T is creative $\iff \deg_M(T^*) = 1$.• T is inconsistent $\iff \deg_M(T^*) = \infty$.• $\deg_M(T^*_0) = \deg_M(T^*_1) = 1$ $\implies T^*_0$ is computably homeomorphic to T^*_1 .
 The Medvedev degree of X (deg_M(X)) is the equivalent class of X under ≡_M. Intuition We might think of a set X as a solution set of some problem. If X has a computable element then we think X as a computably 	 Definition A Σ₁⁰ set is a computably generated open set in the discrete space N of all positive integers. y is Turing reducible to x if there is a computable map F such that F(x) = y.
 solvable problem. We might regard the empty set as a problem without solution. Y ≤_M X (via F) represents that X is more difficult problem than Y. (We have a solution F(x) to Y whenever we know a solution x to X.) Observation The empty set Ø is the most difficult problem. 	 As explained along the way, Medvedev degrees of Π⁰₁ classes are enormously beneficial to understand meta-mathematics and eventually human intelligence. Our theorem clarifies the structure of Π⁰₁ Medvedev degrees, and it suggests that some hard problems on Σ⁰₁ Turing degrees are solvable if we replace the word "Σ⁰₁ Turing" with "Π⁰₁ Medvedev".
 (Ø has the greatest Medvedev degree ∞.) A computably solvable problem is a easiest problem. (Such a problem has the least Medvedev degree 0.) 	Longstanding Open Question (1960-) Is the ∀∃-theory of the Σ ⁰ 1 Turing degrees decidable? Theorem of Cole and Kihara [2]
Essentially Complete Theories, Computably Solvable Problems M Creative Theories Easy Difficult	The $\forall \exists$ -theory of the Π_1^0 Medvedev degrees is decidable! Longstanding Open Question (1940-) Find a "concrete" intermediate Σ_1^0 Turing degree! Theorem of Cenzer-Kihara-Weber-Wu [1] The greatest <i>tree-immune-free</i> Medvedev degree exists, and it is an intermediate Π_1^0 Medvedev degree!

[1] D. Cenzer, T. Kihara, R. Weber, and G. Wu, Immunity and non-cupping for closed sets, *Tbilisi Mathematical Journal* 2 (2009), pp. 77–94.
[2] J. A. Cole, and T. Kihara, The ∀∃-theory of the effectively closed Medvedev degrees is decidable, *Archive for Mathematical Logic* 49 (2010), pp. 1–16.
[3] T. Kihara, Medvedev and Muchnik degrees of Π⁰ classes with incomplete c.e. filters, submitted.

Hypergeometric Series over a *p*-adic field Kensaku Kinjo

Mathematical Institute, Tohoku University, Sendai, Japan

p-adic fields

p: odd prime.

 $\begin{array}{l} \forall x \in Q \backslash \{0\}, \exists !n_x \in Z, \exists \ a, b \in Z \backslash pZ, \, \text{s.t.} \ x = p^{n_x}a/b.\\ \text{Let} \ |x|_p := p^{-n_x} \ (\ |0| := 0). \ \text{Then} \ | \cdot |_p \ \text{defines a distance on } Q.\\ Q_p : \text{ completion of } Q \ \text{with respect to} \ | \cdot |_p.\\ Z_p := \{x \in Q_p; |x|_p \leq 1\} \,. \end{array}$

Monsky-Washnitzer Cohomology

 $\mathrm{Frob}_p \wedge F_p[x,y]/(y^2-x(x-1)(x-\mu)).$

Then Frob_p lifts on

$$R:=Z_p\langle x,y
angle^\dagger/(y^2-x(x-1)(x- ilde{\mu})),$$

where $\tilde{\mu} \in Z_p$ is the Teichmüller lift of μ .

 $egin{aligned} ext{Overconvergent Power Series} \ Z_p \langle x,y
angle^\dagger := egin{cases} & \sum \ \sum \limits_{(i,j \geq 0} a_{i,j} x^i y^j \ | \ a_{i,j} \in Z_p, \
ho^{i+j} |a_{i,j}| o 0 \ (\exists
ho > 1) \ \end{pmatrix} \end{aligned}$

 $D^1(R)$: differential module of R over Z_p . $\mathrm{H}^1(E_\mu) := \mathrm{Coker}[d_1: R \otimes Q_p \to D^1(R) \otimes Q_p]$ (Monsky-Washnitzer cohomology for E_μ)

Property -

 $\cdot \mathrm{H}^{1}(E_{\mu})$ is a 2 dimensional Q_{p} vector space. $\cdot \mathrm{Frob}_{p}$ induces a mapping : $\mathrm{H}^{1}(E_{\mu}) \to \mathrm{H}^{1}(E_{\mu})$.

 a_1, a_2 : eigenvalues of Frob_p on $\operatorname{H}^1(E_\mu)$.

$$egin{aligned} a_1+a_2&=p-\sharp E(F_p)\ a_1\cdot a_2&=p,\ a_1 ext{ or }a_2\in Z_p^{ imes}. \end{aligned}$$

The eigenvalue in Z_p^{\times} is called the unit root of $\mathrm{H}^1(E_{\mu})$.

Family of Elliptic Curves

$$egin{aligned} m &= (p-1)/2, \, H(\lambda) := \lambda(\lambda-1) \sum\limits_{i=1}^m inom{m}{i}^2 \lambda^i, \ B &:= Z_p \langle \lambda, H(\lambda)^{-1}
angle := iggl\{ \sum\limits_{i,j \geq 0} a_{i,j} \lambda^i H(\lambda)^{-j} \ \bigg| \ a_{i,j} \in Z_p, |a_{i,j}|_p o 0 iggr\} \end{aligned}$$

 $egin{aligned} A &:= B\langle x,y
angle^{\dagger}/(y^2 - x(x-1)(x-\lambda)). \ D^1(A/B): ext{ differential module of } A ext{ over } B. \ \mathrm{H}^1 &:= \mathrm{Coker}[d: A \otimes Q_p o D^1(A/B) \otimes Q_p]. \end{aligned}$

NOTE

Let $\mu \in F_p$ is ordinary (: $\Leftrightarrow E_\mu$ is ordinary), and $\tilde{\mu}$ Teichmüller lift of $\mu \in F_p$. Then under $B \ni \lambda \mapsto \tilde{\mu} \in Q_p$, we see

$$\mathrm{H}^1 \otimes_B Q_p \simeq \mathrm{H}^1(E_\mu)$$

 $D^1(A/Z_p)$: differential module of A over Z_p .

$$\exists P,Q\in B[x] ext{ s.t. }$$

$$x(x-1)(x-\lambda)P(x)+rac{3x^2-2(1+\lambda)x+\lambda}{2}Q(x)=1.$$

$$\begin{split} & \underset{L \neq t}{\text{Property}} \underbrace{ \begin{matrix} \text{Property} \\ \text{Let } \tau &:= y P(x) dx + Q(x) dy. \\ \cdot D^1(A/Z_p) &= A d\lambda \oplus A \tau. \\ \cdot D^1(A/B) &= A(dx/y) \simeq D^1(A/Z_p) / A d\lambda. \\ \cdot D^2(A/Z_p) &:= D^1(A/Z_p) \wedge D^1(A/Z_p) = A d\lambda \wedge \tau. \end{split}$$

<u>Gauss-Manin Connectior</u>

 $\forall a \in A, \exists !L(a) \in A \text{ s.t. } d(a\tau) = L(a)d\lambda \wedge \tau.$ Therefore we obtain the mapping

$$D^1(A/B)
i a rac{dx}{y} \mapsto L(a) rac{dx}{y} \in D^1(A/B)$$

 $egin{aligned} & ext{and this induces }
abla : \mathrm{H}^1 o \mathrm{H}^1 ext{ satisfying} \ &\cdot
abla (m_1+m_2) =
abla (m_1) +
abla (m_2) \ (m_1,m_2 \in \mathrm{H}^1), \ &\cdot
abla (bm) = b' \cdot m + b
abla (m) \ (b \in B \otimes Q_p, m \in \mathrm{H}^1). \end{aligned}$

 $\boldsymbol{\nabla}$ is called Gauss-Manin connection.

Frobenius on H¹

 $\phi: B \to B: Z_p ext{ alg. hom. with } \phi(\lambda) = \lambda^p.$

 $\implies \exists \phi \text{-linear ring hom. } F : A \to A \text{ s.t. } F(z) \equiv z^p \pmod{pA}.$ This induces a ϕ -linear mapping : $\mathrm{H}^1 \to \mathrm{H}^1$, which is also denoted by F.

 $\begin{array}{l} \text{Proposition} & \\ \text{Let } \omega \text{ be the image of } dx/y \text{ in } \mathrm{H}^1. \\ \cdot \left\{ \omega, \nabla(\omega) \right\} \text{ is a basis for } \mathrm{H}^1 \text{ as } B \otimes Q_p \text{ modules.} \\ \cdot 4\lambda(1-\lambda)\nabla^2(\omega) + 4(1-2\lambda)\nabla(\omega) - \omega = 0. \\ \cdot \eta \in \ker \nabla \Leftrightarrow \eta = \lambda(1-\lambda)f'\omega - \lambda(1-\lambda)f\nabla(\omega) \text{ where } f \in B \\ \text{ satisfying } 4\lambda(1-\lambda)f'' + 4(1-2\lambda)f' - f = 0. \\ \cdot F \text{ is stable on } \ker \nabla. \\ \cdot F \text{ is also stable on } B\omega + B\nabla(\omega). \text{ Moreover, } \exists u \neq 0 \in B\omega + B\nabla(\omega) \text{ s.t. } Bu \text{ is the unique direct summand of } B\omega + B\nabla(\omega) \text{ satisfying } F(Bu) = Bu. \end{array}$

Let $\mu \in F_p$ be an ordinary and $\tilde{\mu}$ Teichmüller lift of μ . Then under $\lambda \mapsto \tilde{\mu}$, we obtain the following commutative diagram :

$$egin{array}{cccccccc} \mathrm{H}^1 & \stackrel{F}{
ightarrow} & \mathrm{H}^1 & u & \mapsto & \xi(\lambda) u \ & \downarrow & \downarrow & \downarrow & \downarrow & \downarrow & \downarrow \ \mathrm{H}^1(E_\mu) \stackrel{
ightarrow}{\mathrm{Frob}} & \mathrm{H}^1(E_\mu) \;, & v & \mapsto & \mathrm{Frob}_p(v) = \xi(ilde{\mu}) \end{array}$$

where v is the image of u in $\mathrm{H}^1(E_\mu)$. Therefore we see $\xi(\tilde{\mu})$ is the unit root of $\mathrm{H}^1(E_\mu)$.

 $egin{aligned} & ext{Theorem[Dwork]} & ext{-} & \cdot \xi(\lambda): ext{holom. on } \left\{\lambda \in Q_p \ | \ |\lambda|_p \leq 1, |H(\lambda)|_p = |\lambda(\lambda-1)|_p
ight\}. \ & \cdot \xi(\lambda) = (-1)^m \mathcal{F}(\lambda) / \mathcal{F}(\lambda^p) \quad ext{on } \left\{\lambda \in Q_p \ ; |\lambda|_p < 1
ight\}, ext{ where} \ & \mathcal{F}(x) := \sum_{n \geq 0} inom{(2n-1)/2}{n}^2 x^n \end{aligned}$

is the hypergeometric series.

On Hasse principle of purely transcendental extension field in one variable

Makoto Sakagaito

2010/2/18

Classical result

k: global field(i.e. number field or algebraic function field in one variable over finite field).

 \mathfrak{p} : prime of k including the archimedean ones if k is number field.

 $\hat{k_{\mathfrak{p}}}$: completion of k at \mathfrak{p} .

Hasse principle

local-global map

$$\mathrm{Br}(k) \to \prod_{\mathfrak{p}} \mathrm{Br}(\hat{k_{\mathfrak{p}}})$$

is injective. Here, \mathfrak{p} runs through all primes of k.

Special case

For $a, b \in \mathbb{Z}$, q: prime number There exist $x, y \in \mathbb{Q}$ s.t. $ax^2 + by^2 = q$ \Leftrightarrow There exist $x_p, y_p \in \mathbb{Q}_p$ s.t. $ax_p^2 + by_p^2 = q$ Here, p runs through all prime numbers and ∞ , $\mathbb{Q}_{\infty} = \mathbb{R}$.

Application

There exist $x, y \in \mathbb{Q}$ s.t. $q = x^2 + 26y^2$

 $\begin{cases} q \equiv 1 \mod 8 \text{ and } q \equiv \text{one of } 1, 3, 4, 9, 10, 12 \mod 13 \\ \text{or} \\ q \equiv 3 \mod 8 \text{ and } q \equiv \text{one of } 1, 3, 4, 9, 10, 12 \mod 13 \end{cases}$

Study

Related to this classical result, I studied the question of whether the Hasse principle holds in the case where the field K(X) is a purely transcendental extension in one variable of any field K.

I ascertained the result that the local-global map of K(X) restricted to

$$\operatorname{Br}(K(X))' = \operatorname{Ker}(\operatorname{Br}(K(X)) \xrightarrow{\operatorname{res}} \prod_{\mathfrak{p} \in \operatorname{Spec}(K[X])} \operatorname{Br}(K(X)_{\bar{\mathfrak{p}}})$$

is injective.

Here, $K(X)_{\bar{p}}$ is the maximal unramified extension field at $p \in \text{Spec}(K[X])$.

Remarks

In the case where m is a natural number prime to the characteristic of K,

$$\operatorname{Br}(K(X))_m = \operatorname{Ker}(\operatorname{Br}(K(X)) \xrightarrow{\times m} \operatorname{Br}(K(X))) \subset \operatorname{Br}(K(X))'$$

In the case where K is a perfect field,

$$Br(K(X))' = Br(K(X)).$$

References

[1] K.Kato, N.Kurokawa, T.Saito : Number theory.1 Fermat's dream

Large time behavior of solutions for system of nonlinear damped wave equations

Hiroshi Takeda (Tohoku University)

Feb. 18-19th, 2010

Abstract: We consider the Cauchy problem for a system of semilinear damped wave equations with small initial data. We derive the asymptotic profile of the nonlinear system corresponding to the results for the nonlinear single equation and we obtain the sufficient condition of the growth order on the nonlinear term to ensure the existence of the solution with the optimal decay. Our proof is based on the analysis for the fundamental solution of the linear damped wave equation.

1 Nonlinear damped wave system

$$\begin{split} &m\geq 2, \quad u:\mathbb{R}_+\times\mathbb{R}^n\to\mathbb{R}^m; \text{ unknown vector-valued function},\\ &F:\mathbb{R}^m\to\mathbb{R}^m; \ F_j(u)\sim\prod_{k=1}^m|u_k|^{p_{j,k}}. \end{split}$$

 $(\mathrm{DW}) \ \begin{cases} \partial_t^2 u - \Delta u + \partial_t u = F(u), \quad t > 0, \quad x \in \mathbb{R}^n, \\ u(0,x) = a(x), \ \partial_t u(0,x) = b(x), \ x \in \mathbb{R}^n. \end{cases}$

2 Problem

• When does the initial value problem (DW) have a global solution ? Find the critical condition for the nonlinear term to ensure the existence of time in global solutions!

3 Known Results

Sun-Wang (2007) $m = 2, p_{1,1} = p_{2,2} = 0, p_{1,2}, p_{2,1} > 1$

$\frac{\max\left\{p_{1,2}, p_{2,1}\right\} + 1}{p_{1,2}p_{2,1} - 1} \ge \frac{n}{2}$	$\frac{\max\left\{p_{1,2}, p_{2,1}\right\} + 1}{p_{1,2}p_{2,1} - 1} < \frac{n}{2}$
finite time blow-up	small data global existence
$(n \ge 1)$	(n = 1, 3)

4 Notation and assumptions

$$P = \begin{pmatrix} p_{1,1} & \dots & p_{1,m} \\ \vdots & & \vdots \\ p_{m,1} & \dots & p_{m,m} \end{pmatrix}, \alpha = (P - I)^{-1} \begin{pmatrix} 1 \\ \vdots \\ 1 \end{pmatrix} = \begin{pmatrix} \alpha_1 \\ \vdots \\ \alpha_m \end{pmatrix}$$

- $p_{j,k} \in [1,\infty) \cup \{0\}, \sum_{k=1}^{m} p_{j,k} > 1, j,k = 1, \cdots, m,$
- $\exists (P-I)^{-1},$
- $(a,b) \in \{W^{1,1}(\mathbb{R}^n) \cap W^{1,\infty}(\mathbb{R}^n) \times L^1(\mathbb{R}^n) \cap L^\infty(\mathbb{R}^n)\}^m$.

5 Results

Theorem. 1. (existence of the global solution) n = 1, 2, 3.

$$0 < \alpha_j < \frac{n}{2}, \quad j = 1, \cdots, m, \quad (a, b) :$$
small

 $\implies \exists! \ u(t)$: global solution of (DW) in

$$\{C([0,\infty);L^1(\mathbb{R}^n)\cap L^\infty(\mathbb{R}^n))\}^m$$

$$F_j(u) = |u_l|^{p_j} : j - l \ cyclic, \ P = \begin{pmatrix} 0 & \dots & 0 & p_1 \\ p_2 & \dots & 0 & 0 \\ \vdots & \ddots & \vdots & \vdots \\ 0 & \dots & p_m & 0 \end{pmatrix}.$$

Theorem. 2. (nonexistence of global solutions) $n \ge 1$,

$$\max_{1 \le j \le m} \alpha_j \ge \frac{n}{2},$$
$$\int_{\mathbb{R}^n} a_j(x) dx > 0, \quad \int_{\mathbb{R}^n} b_j(x) dx > 0, \qquad j = 1, \cdots, m$$

 $\implies u(t)$: blow-up i.e. $\exists T < \infty$ s.t.

 $\limsup_{t \nearrow T} \|u(t)\|_{L^1(\mathbb{R}^n) \cap L^\infty(\mathbb{R}^n)} = \infty.$

Theorem. 3. (asymptotic profile) n = 1, 2, 3,

$$\min_{1 \le j \le m} \sum_{k=1}^{m} p_{j,k} > 1 + \frac{2}{n}, \quad (a,b) : \text{small}$$

 $\implies \exists! u(t):$ global solution of (DW),

$$G_t(x) := \left(\frac{1}{4\pi t}\right)^{\frac{n}{2}} e^{-\frac{|x|^2}{4t}} \quad \text{(Heat kernel)},$$
$$\|u(t) - MG_t\|_{L^p(\mathbb{R}^n)} = o(t^{-\frac{n}{2}(1-\frac{1}{p})}), \quad \text{as} \quad t \to \infty$$

where $M = (M_1, \cdots, M_m),$

$$M_j = \int_{\mathbb{R}^n} a_j(y) + b_j(y)dy + \int_0^\infty \int_{\mathbb{R}^n} F_j(u(s,y))dy\,ds.$$

6 Remarks on the theorems

For the j-l cyclic case, we see the relationship between the Theorem 1, 2 and 3.



Kazuaki Ta ima

athematical nstitue Tohoku niversity D1

1 Classical Iwasawa Theory

Let p be a prime number. \mathbb{Z}_{p} -extension of a number field k is a Galois e tension k_{∞}/k with Galois group Γ Gal $k_{\infty}/k \simeq \mathbb{Z}_{p}$ the additive group of p-adic integers. Such a \mathbb{Z}_{p} -e tension can be regard as a tower of fields

$$k \quad k_0 \subset k_1 \subset \cdots \subset k_n \subset \cdots \subset k_\infty \quad \bigcup_{n \ge 0} k_n$$

with Gal $k_n/k \simeq \mathbb{Z}/p^n\mathbb{Z}$ where the fields k_n are uni ue sube tensions of k_{∞}/k with $k_n k p^n$.

Let $k \zeta_{p^{\infty}}$ be the e tension of k obtained by adjoining all the primitive roots of unity of p-power order. hen there e ist precisely one \mathbb{Z}_p -e tension of k inside $k \zeta_{p^{\infty}}$. his e tension is called **the cyclotomic** \mathbb{Z}_p -extension of k.

Let A_n denote the *p*-primary part of the ideal class group of k_n . y class field theory A_n is isomorphic to the Galois group of the ma imal abelian unramified *p*e tension L_n of k_n and these groups have finite order. hen X_{∞} lim A_n is isomorphic to the Galois group of the ma imal unramified abelian *p*-e tension L_{∞} of k_{∞} . ence we know that $\Lambda \quad \mathbb{Z}_p \ \Gamma$ acts naturally

on X_{∞} . In 1 Iwasawa was studied the structure of X_{∞} as Λ -module and he proved following beutiful theorem

heorem 1.1 I as a a 1959). There exist non-negative integers $\mu k_{\infty}/k$, $\lambda k_{\infty}/k$ and an integer $\nu k_{\infty}/k$ such that

 $\sharp A_n \quad p^{\mu(k_\infty/k)p^n + \lambda(k_\infty/k)n + \nu(k_\infty/k)}$

for all sufficiently large integer $n \geq .$

he integers μ and λ are the invariant of the structure of X_{∞} and sometimes also ν are called **the i asa a invariants**. specially λ rank_{\mathbb{Z}_p} X_{∞} holds. For the cyclotomic e tension the iwasawa invariants are denoted by $\mu_p k$, $\lambda_p k$ and $\nu_p k$.

Calculate the i as a invariants It is no known that general way to calculate the iwasawa invariants.

ut following result gives a simple sufficient condition such that all the iwasawa invariants vanish

heorem 1.2 I as a 1956). Suppose that p does not divide the class number $h \ k$ of k and p is non-split in k/\mathbb{Q} . Then we have $\mu_p \ k \qquad \lambda_p \ k \qquad \nu_p \ k$.

For e ample we see that $\mu_p \mathbb{Q} \quad \lambda_p \mathbb{Q} \quad \nu_p \mathbb{Q}$ for any prime p.

2 Non-abelian Iwasawa Theory

For a \mathbb{Z}_p -e tension k_{∞}/k we consider the Galois group of the ma imal unramified **not nessesary abelian** Our main purpose is to study the group structure of \widetilde{G}_{∞} by using "i as a thoretical methods". ut it is very hard to describe this structure in direct so we shall consider a filtration of \widetilde{G}_{∞}

$$\widetilde{G}_{\infty} \quad C_1 \ \widetilde{G}_{\infty} \ \supset C_2 \ \widetilde{G}_{\infty} \ \supset \cdots \supset C_i \ \widetilde{G}_{\infty} \ \supset \cdots$$

ore precisely we adopt the lower central se uence as the filtration that is $C_{i+1} G_{\infty} = C_i G_{\infty}, G_{\infty}$ $i \geq 1$ the commutator group of $C_i \widetilde{G}_{\infty}$ and \widetilde{G}_{∞} . hen we define **the** *i*-th i as a module by $X_{\infty}^{(i)}$ $C_i \widetilde{G}_{\infty} / C_{i+1} \widetilde{G}_{\infty}$. e also define $X_n^{(i)}$ as similary. Let $L_{\infty}^{(i)}$ resp. $L_{n}^{(i)}$ be the fined field of \widetilde{L}_{∞} resp. $\widetilde{\widetilde{L}}_n$ by C_{i+1} \widetilde{G}_{∞} resp. C_{i+1} \widetilde{G}_n . у definition we have $X_{\infty}^{(i)}$ Gal $L_{\infty}^{(i)}/L_{\infty}^{(i-1)}$, $X_{n}^{(i)}$ Gal $L_{n}^{(i)}/L_{n}^{(i-1)}$ especially $X_{\infty}^{(1)} \sim X_{\infty}$ and $X_{n}^{(1)} \sim A_{\infty}$ A_n are the usual iwasawa modules. Furtheremore we know that $L_{\infty}^{(i)}$ is the central *p*-class field of $L_{\infty}^{(i-1)}/k_{\infty}$ namely $L_{\infty}^{(i)}$ is the ma imal unramified pe tention of $L_{\infty}^{(i-1)}$ such that $L_{\infty}^{(i)}/k_{\infty}$ is Galois and $X_{\infty}^{(i)}$ Gal $L_{\infty}^{(i)}/L_{\infty}^{(i-1)}$ is contained in the center of Gal $L_{\infty}^{(i)}/k_{\infty}$. his is a reason why we adopt the lower central se uence as a filtration of \widetilde{G}_{∞} and this fact plays important role. In fact we know that $X_n^{(i)}$ is

finite and $X_{\infty}^{(i)}$ has the natural Λ -module structure by the virtue of the fact. Furthremore if $\mu k_{\infty}/k$ we can prove $X_{\infty}^{(i)}$ is a finitely generated torsion Λ -module so we define **the** *i*-th λ -invariant by $\lambda^{(i)} k_{\infty}/k$ rank_{\mathbb{Z}_p} $X_{\infty}^{(i)}$.

he following theorem which is an analogy of heorem 1.1 is proved by O aki

heorem 2.1 O aki 2007). Suppose $\mu k_{\infty}/k$. Then there exist an integer $\nu^{(i)} k_{\infty}/k$ for all $i \ge 1$ such that

$$\sharp X_n^{(i)} \quad p^{\lambda^{(i)}(k_\infty/k)n + \nu^{(i)}(k_\infty/k)}$$

for all sufficiently large integer $n \geq .$

From heorem 2.1 we can regard the *i*-th iwasawa module $X_{\infty}^{(i)}$ as sufficiently good objects to study the structure of \widetilde{G}_{∞} by using iwasawa theoretical methods .

References

1 .O aki Non-abelian Iwasawa theory of \mathbb{Z}_p extensions ournal fur die reine und angewandte athematik 2 2 –
The quadratic subextension of the class field of a real quadratic field

Toshihide DOI

sa6m22@math.tohoku.ac.jp Mathematical Institute, Tohoku University, Japan

The 2nd GCOE internaitonal symposium in Sendai Feb. 18 - 19, 2010.

1 INTRODUCTION

Let F be an arbitrary algebraic number field and H_F be the maximal unramified (including ∞) abelian extension of F. From the global class field theory, the extension H_F/F is finite and its galois group $G = \text{Gal}(H_F/F)$ is isomorphic to the ordinary ideal class group of F. The field H_F is called the Hilbert class field of F. (For F, Cl_F and h_F denote the ordinary ideal class group of F and its order that is called the class number of F, respectively.)

In general, it is difficult to derive H_F from F. But, if F is quadratic field, there are several way to obtain the H_F . Among these ways, there is a way for real quadratic fields with even class number. The purpose of this poster is to introduce that way with examples of $\mathbb{Q}(\sqrt{85})$.

2 h_F AND FUNDAMENTAL UNIT OF F

Let F be a real quadratic field $\mathbb{Q}(\sqrt{m})$ where m is a square-free positive rational integer, D be the discriminant of F, and ε be the fundamental unit of F.

2.1 DEFINITIONS and TERMS

Deffinition 1

If $a, b, c \in \mathbb{Q}$ satisfy $D = b^2 - 4ac$ and (a, b, c) = 1, then the roots θ, θ' of

$$ax^2 + bx + c = 0$$

are called a quadratic irrational number in D (denote by QI in D). Moreover, if θ satisfies

$$\theta > 1, \ 0 > \theta' > -1,$$

then θ is called a reduced quadratic irrational number in D (denote by RQI in D).

Set the relation "~" on \mathbb{R} as follows; For $\theta, \theta_0 \in \mathbb{R}$, if $\theta \sim \theta_0$ then there exist $a, b, c, d \in \mathbb{Z}$ such that

$$\theta = \frac{a\theta_0 + b}{c\theta_0 + d}$$
 and $\det \begin{pmatrix} a & b \\ c & d \end{pmatrix} = \pm 1.$

" \sim " satisfies equivalence condition, transitive, symmetric, reflexive. If

$$heta_n-a+rac{1}{ heta_{n+1}}-rac{a heta_{n+1}+1}{1\cdot heta_{n+1}+0},$$

then $\theta_n \sim \theta_{n+1}$ since det $\begin{pmatrix} a & 1 \\ 1 & 0 \end{pmatrix} = -1$. So each end-term θ_n of the expanding to (all numerator equal 1) continued fraction of $\theta \in \mathbb{R}$ satisfies $\theta \sim \theta_n$.

2.2 FACTS and THEOREMS

- **Fact1:** If θ is QI in D and $\theta \sim \theta_0$, then θ_0 is also QI in D and appear in the expanding to continued fraction of θ as an end-term.
- Fact2: All QIs are expanded to cyclic continued fraction. Moreover, all end-term of cyclic part of continued fraction of a QI are RQI.

Theorem 1

The number of RQI in D is finite. More precisely, the necessary and sufficient condition of that $ax^2 + bx + c = 0, a, b, c \in \mathbb{Z}$ has a root of RQI in D is below;

$$\begin{split} b < 0, \ |b| < \sqrt{D}, \\ a > 0, \ c < 0, \\ |ac| &= \frac{\sqrt{D} + |b|}{2} \cdot \frac{\sqrt{D} - |b|}{2}, \\ \frac{\sqrt{D} + |b|}{2} > a, \ |c| > \frac{\sqrt{D} - |b|}{2} \\ (a, b, c) &= 1. \end{split}$$

Theorem 2

Theorem 3

Let $F = \mathbb{Q}(\sqrt{m}), m > 0$, square-free, and D be the discriminant of F, and set h be the number of classes of {All QI in D}/~. Then

$$h = h_1$$

F,D as above. Let θ be RQI in D and ε be the fundamental unit of F and

$$heta = k_0 + rac{1}{k_1 + \cdots + rac{1}{k_{n-1} + rac{1}{a}}} = rac{p_n heta + p_{n-1}}{q_n heta + q_{n-1}}$$

be the first section of cycles of the expanding to continued fraction. Then

$$\varepsilon = q_n \theta + q_{n-1}$$

Now we can calculate h_F and ε by using above facts & theorems.

Example If $F = \mathbb{Q}(\sqrt{85})$, then D = 85. From theorem 1,

b	ac	range of $a, \left c \right $	a	c	θ
-9	1	[1,9]	1	-1	$\frac{9+\sqrt{85}}{2} = \theta_1$
$^{-7}$	9	[3, 8]	3	$^{-3}$	$\frac{7+\sqrt{85}}{6} = \theta_2$
$^{-5}$	15	[3, 7]	3	-5	$\frac{5+\sqrt{85}}{6} = \theta_3$
"	"	"	5	-3	$\frac{5+\sqrt{85}}{10} = \theta_4$

The expanding to continued fraction of θ_2 tells us $\theta_2 \sim \theta_3 \sim \theta_4$ so $h_F = 2$. And since $\theta_1 = 9 + \frac{1}{\theta_1} = \frac{9\theta_1 + 1}{\theta_1}$, we get $\varepsilon = \theta_1$.

3 HILBERT CLASS FIELD OF F WITH $h_F = 2$

To determinate the Hilbert class field H_F of a real quadratic field F with $h_F = 2$, we use following theorems.

Theorem 4

For an algebraic extention K/F, $\mathcal{D}(K/F)$, $\mathcal{D}(\alpha, K/F)$ denote the relative different of K/F and the relative different of K/F for $\alpha \in K$, respectively. Let $f_{\alpha}(x)$ be the minimal polynomial of $\alpha \in K$ over F. Then, for all $\alpha \in K$,

$$\mathcal{D}(K/F) \mid \mathcal{D}(\alpha, K/F) \text{ and } \mathcal{D}(\alpha, K/F) = (f_{\alpha}'(\alpha)).$$

Theorem 5

Let $K, F, \mathcal{D}(K/F)$ same as avobe and \mathfrak{p} be a prime of F. Then

$$\mathfrak{p}$$
 ramifies in $K/F \iff \mathfrak{p} \mid \mathcal{D}(K/F)$.

From above theorems, we only have to find $\alpha_1, \alpha_2, \ldots$ such that

$$\bigcap_{i=1} (f_{\alpha_i}'(\alpha_i)) = 1$$

to say that K/F is unramified extension.

Example

Let $F = \mathbb{Q}(\sqrt{85})$ and $K_1 = F(\sqrt{\varepsilon}), K_2 = F(\sqrt{5}), K_3 = F(\sqrt{5\varepsilon})$. Where $\varepsilon = \frac{9 + \sqrt{85}}{2} = \theta_1$. We prove that K_2/F is unaramified extension in following. Let $\alpha_1 = \sqrt{5}$ and $\alpha_2 = \frac{1 + \sqrt{17}}{2}$. Then

$$f_{\alpha_1} = x^2 - 5, \ f_{\alpha_2} = x^2 - x - 4$$

hence $\mathcal{D}(\alpha_1, K_2/F) = (2\sqrt{5}), \ \mathcal{D}(\alpha_2, K_2/F) = (\sqrt{17}).$

and hence $\mathcal{D}(\alpha$ Consequently, we obtain

$$\mathcal{D}(K_2/F) \mid \mathcal{D}(\alpha_1, K_2/F) \cap \mathcal{D}(\alpha_2, K_2/F) = (1),$$

so that K_2/F is unramified extension.

The pattern formation of head regeneration model of Hydra

Madoka Nakayama (Tohoku University)

Abstract: Hydra is a small animal living in fresh water, which is best known for its ability of regeneration. When a hydra is cut into two pieces, two hydras will regenerate. There has been several mathematical models proposed to describe this experiment. A classical model proposed by Gierer and Meinhardt in 1972, is based on the idea of diffusion- driven instability between two chemicals called activator and inhibitor. Recently, Anna Marciniak proposed new regeneration models from a more biologically refined viewpoint. Her models consist of free and bound receptors, ligands and an enzyme, and a head is formed at place of higher bound receptor concentration.

1 What is Hydra?



At first we consider the following problem for

$$\frac{\partial r_f}{\partial t} = -\mu_f r_f - br_f l + dr_b + m_1,$$
(1)
$$\frac{\partial r_b}{\partial t} = -\mu_b r_b + br_f l - dr_b,$$
(2)
$$\frac{\partial l}{\partial t} = \frac{1}{\gamma} \frac{\partial^2 l}{\partial x^2} - \mu_l l - br_f l + dr_b + p_l,$$
(3)
$$\frac{\partial pl}{\partial t} = -\delta_l \frac{p_l}{1 + p_l^2} + \frac{m_2 l r_b}{(1 + \sigma_l p_l^2 - \beta_l p_l)(1 + \alpha_l r_b)}$$
(4)

 $r_f \ge 0$:density of free receptors,

 $r_b \ge 0$:density of bound receptors,

 $l \geq 0 : \text{density of bound ligands},$

 $p_l \ge 0$:production term of ligands,

 $\mu_f, \mu_b, \mu_l, m_1, m_2, b, d, \sigma_l, \beta_l, \delta_l$: positive constants.

4 Stationary problem

To consider stationary problem, we put

 $\frac{\partial r_f}{\partial t} = \frac{\partial r_b}{\partial t} = \frac{\partial l}{\partial t} = \frac{\partial pl}{\partial t} = 0,$

and reduce (1) to the following two component system:

$$\begin{cases} \frac{1}{\gamma} \frac{d^2 u}{dx^2} + f(u, v) = 0, \\ g(u, v) = 0 \\ \frac{du}{dx} = \frac{dv}{dx} = 0, (x = 0, 1) \end{cases}$$

where $u = l, v = p_l$ and let $\mu_b \mu_f + d\mu_f = M$,

$$\begin{aligned} f(u,v) &= v - \mu_l u - \frac{m_1 \mu_b b u}{M + \mu_b b u}, \\ g(u,v) &= -\delta_l \frac{v}{1 + v^2} + \frac{m_1 m_2 b u^2}{(M + \mu_b b u + \alpha_l m_1 b u)(1 + \sigma_l v^2 - \beta_l v)}. \end{aligned}$$

To draw nullclines, we consider f(u, v) = 0, g(u, v) = 0.

$$\begin{split} f(u,v) &= 0 \Leftrightarrow v = \mu_l u - \frac{m_1 \mu_b b u}{M + \mu_b b u}, \\ g(u,v) &= 0 \Leftrightarrow \delta_l \frac{v}{1 + v^2} (1 + \sigma_l v^2 - \beta_l v) = \frac{m_1 m_2 b u^2}{M + u(\mu_b b + \sigma_l m_1 b)} =: \Psi(u) \end{split}$$

 $\begin{array}{l} \mbox{Lemma} \\ \mbox{Let } A = \frac{\psi(u)}{\delta_l}. \mbox{ Assume that} \\ (i) \ 9\sigma_l - 1 > 0 \ {\rm and} \ \sqrt{3\sigma_l} < \frac{\beta_l}{9\sigma_l - 1}, \ \sqrt{3\sigma_l} < A < \frac{\beta_l}{9\sigma_l - 1}, \ {\rm or} \\ (ii) \ 9\sigma_l - 1 > 0 \ {\rm and} \ \frac{\beta_l}{9\sigma_l - 1} < \sqrt{3\sigma_l}, \ \frac{\beta_l}{9\sigma_l - 1} < A < \sqrt{3\sigma_l}. \\ \mbox{In addition, if } A \ {\rm satisfies} \\ -(A + \beta_l) \ - \ \sqrt{(A + \beta_l)^2 - 3\sigma_l} \ < \ \frac{3\sigma_l \{\beta_l + A(1 - 9\sigma_l)\}}{2(3\sigma_l - A^2)} \ < \ -(A + \beta_l) \ + \\ \sqrt{(A + \beta_l)^2 - 3\sigma_l}, \\ \mbox{then } g(u, v) \ = \ 0 \ \ {\rm and} \ f(u, v) \ = \ 0 \ \ {\rm intersect} \ \ {\rm at three \ points} \ (u_-, v_-), (u_m, v_m), \\ (u_+, v_+). \end{array}$

From the Lemma, we obtain the relationship of parameters for which the system (1)-(4) has three stationally solutions.

Let $(u_{-} < u_{m} < u_{+})$

$$\begin{cases} v = h_0(u) (u_0 < u < u_+, v_0 < v < v_+) \\ v = h_m(u) (u_0 < u < u_+, v_+ < v < v_-) \\ v = h_1(u) (u_0 < u < u_+, v_- < v < v_1) \end{cases}$$
(5)

In particular, we focus on the curve $v = h_j(u)(j = 0, 1)$.



We substitute $v = h_j(u)$ (j = 0, 1) in f(u, v), then we have the following systems.

$$\begin{pmatrix}
\frac{1}{\gamma} \frac{d^2 u}{dx^2} + f(u, h_j(u)) = 0, \\
\frac{du}{dx} = \frac{dv}{dx} = 0 \quad (x = 0, 1).
\end{cases}$$
(6)

From the lemma, we obtain the following theorem.

Theorem

There exists a monotone increasing solution of (6) which is C^1 in [0,1].

Outline of the proof

We consider the following initial value problems

$$\frac{1}{\gamma} \frac{d^2 \hat{u}}{dx^2} + f_0(\hat{u}, h_0(\hat{u})) = 0,
\frac{d\hat{u}}{dx}(0) = 0,
\hat{u}(0) = k \quad (u_0 < k < u_+),
\frac{1}{\gamma} \frac{d^2 \bar{u}}{dx^2} + f_0(\bar{u}, h_0(\bar{u})) = 0,$$
(7)

$$\gamma \, dx^2 = 0 \, (\gamma + 0 \, (\gamma))^{-1} \, (\beta + 1)^{-1} \, (\beta + 1)^{-1}$$

We show \hat{u} is monotone increasing, and \bar{u} is monotone decreasing. Combining these facts, we can choose k and p so that the graphs of \hat{u} and \bar{u} are tangent to each other at $u = \beta$, where β is an arbitrarily fixed constant in $[u_0, u_+] \cap [u_-, 1]$.

The computational methods of the canonical height on elliptic curves

Tadahisa Nara (Mathematical Institute, Tohoku University)

Elliptic curves and the canonical height

 \mathbb{Q} : set of rational numbers

The elliptic curve E over \mathbb{Q} is defined as a curve of the form of $y^2 = x^3 + Ax + B$, $(A, B \in \mathbb{Q})$ having some conditions. And then we think of the point (∞, ∞) as is on E and denote it O. In number theory one of interests about elliptic curves is studying rational points, that is a solution in \mathbb{Q} of the equation above. Mordell-Weil theorem is one of basic guidelines for this problem, which states that all the ratoinal points on an elliptic curve form an finitely generated abelian group. This means that you can know all the rational points if you find specific points of finite number on the curve.

The height of a point is numerical size of the points for certain arithmetic. The canonical height is one of heights which is defined on elliptic curves.

Definition 0.1 The naive height on $E(\mathbb{Q})$ denoted by h is function defined by

 $h: E(\mathbb{Q}) \to \mathbb{R}, \ P = (a/c, *) \mapsto \log \max\{|a|, |c|\}.$

The canonical height on $E(\mathbb{Q})$ denoted by \hat{h} is function defined by

$$\hat{h}: E(\mathbb{Q}) \to \mathbb{R}, \ P \mapsto \lim_{n \to \infty} \frac{1}{4^n} h([2^n]P),$$

where $[2^n]P$ means 2^n times addition of P about its group operation.

The features of the canonical height are as follows.

• $\hat{h}([m]P) = m^2 \hat{h}(P)$

•
$$\hat{h}(P+Q) + \hat{h}(P-Q) = 2\hat{h}(P) + 2\hat{h}(Q)$$

Computation of the canonical height The definition above, which is usually used, is brief but not suitable for practical computaion. Neron and Tate showed that there is local height function λ_p such that $\hat{h}(P) =$ $\sum_{p:prime,\infty} \lambda_p(P)$. So the canonical height is decomposed into something local. Those factors of the decomposition are not difficult to compute. For example for all prime p that do not divide $4A^3 + 27B^2$ nor c (the denominator of the x-coordinate of P) $\lambda_p(P) = 0$, so this is a finite sum indeed. For other primes we can determine algebraically $\lambda_p(P)$ using the results of reduction type of elliptic curves. And about the case of $p = \infty$ following theorem is useful.

Theorem 0.2 (Tate) Suppose that there is an $\epsilon > 0$ so that every point Pin $E(\mathbb{Q})$ satisfies $|x(P)| > \epsilon$. Then for all $P(\neq O) \in E(\mathbb{Q}), \lambda_{\infty}(P) =$ $\log |x(P)| + \frac{1}{4} \sum_{n=0}^{\infty} 4^{-n} \log |z([2]^n P)|$, where x(P) = x-coordinate of P and $z(P) = 1 - 2A/x(P)^2 - 8B/x(P)^3 + A^2/x(P)^4$.

Using modified versions of above theorem, we see following.

Theorem 0.3 Let α and β be coprime odd integers with $\alpha \geq 1.232\beta > 3$ and E/\mathbb{Q} : $y^2 = x^3 + \alpha^6 + \beta^6$ be a elliptic curve. Suppose every order of the prime divisors of $\alpha^6 + \beta^6$ is odd and $P_1 + P_2 \notin 2E(\mathbb{Q})$, where $P_1 = (-\alpha^2, \beta^3), P_2 = (-\beta^2, \alpha^3)$. Then there is a basis of $E(\mathbb{Q})$ including P_1 and P_2 , in particular if the rank $E(\mathbb{Q}) = 2$, $\{P_1, P_2\}$ is basis of the free part of $E(\mathbb{Q})$.

Spatial branching process in random environment

NISHIMORI Yasuhito (Mathematical Institute Tohoku University)

Introduction

An outline of these notes is as follows. After defining the Galton-Watson process, we introduce the branching Brownian motion among Poisson obstacles. In the following section, we present some result on exponential growth. Our main purpose is to compare the growth order of Galton-Watson process with strictly dyadic branching Brownian motion among Poisson obstacles.

1 Galton-Watson process

We consider Galton-Watson process $\{G_n\}_{n=0}^{\infty}$ with offspring distribution $\{p_k\}_{k=0}^{\infty}$ satisfying $m = \sum_{k=1}^{\infty} kp_k > 1$ and $p_k \neq 1$ for any $k \in \mathbf{N} \cup \{0\}$.

Theorem 1 (H.Kesten, B.P.Stigum, 1966)

Let $W = \lim_{n \to \infty} m^{-n} G_n$. Then

$$\sum_{k=2}^{\infty} p_k k \log k < \infty \implies E[W] = 1 \quad . \tag{1}$$

Theorem 1 says the growth rate of $\{G_n\}$.

2 Branching Brownian motion and Poisson obstacles

Firstly, we call $(\{\hat{Z}_t\}_{t\geq 0}, P_x)$ strictly dyadic branching Brownian motion on \mathbf{R}^d with branching rate $\beta(x)$, if for any $g \in C_b^+$, $u(x,t) = E_x[e^{-\langle g, \hat{Z}_t \rangle}]$ solves

$$\begin{array}{l} \displaystyle \frac{\partial u}{\partial t} = \frac{1}{2} \triangle u + \beta (u^2 - u) \quad on \quad \mathbf{R}^d \times (0, \infty) \\ \displaystyle \lim_{t \downarrow 0} u(\cdot, t) = e^{-g(\cdot)} \quad on \quad \mathbf{R}^d \times (0, \infty) \\ \displaystyle 0 \le u \le 1, \end{array}$$

where $\langle g, \hat{Z}_t \rangle = \int_{\mathbf{R}^d} g(x) \hat{Z}_t(dx)$. $(\{\hat{Z}_t\}_{t \ge 0}, P_x)$ is following that one particle starts at $x \in \mathbf{R}^d$, performing an Brownian motion

one particle starts at $x \in \mathbf{R}$, performing an Browman motion on \mathbf{R}^d . Her life time distributions is exponential with parameter $\beta(x)$. Just as it dies, one particle splits into two one and descendants perform same as their parent.

Secondly, we set random environment by (ω, \mathbf{P}) which is a Poisson point process in \mathbf{R}^d with intensity measure $\nu(dx) = \nu dx$ where $\nu > 0$ and dx is the Lebesgue measure. For a > 0, let K_{ω} denote a random set given by *a*-neighborhood of each configuration ω :

$$K_{\omega} \stackrel{\text{def}}{=} \bigcup_{x \in supp(\omega)} \bar{B}(x, a)$$

And for fixed $0 < \beta_1 < \beta_2$, we define the strictly dyadic branching Brownian motion $(\{Z_t\}_{t\geq 0}, P^{\omega})$ with branching rate $\beta(x)$:

$$\beta(x) = \beta_1 \cdot \chi_{K_{\omega}}(x) + \beta_2 \cdot \chi_{K_{\omega}^c}(x)$$

Here χ is indicator function. Then, as long as a particle in K_{ω} , the splitting rate is less than in K_{ω}^c . Thus we may consider the random field K_{ω} as obstacles. So K_{ω} is said to Poisson obstacles and branching Brownian motion $(\{Z_t\}_{t\geq 0}, P^{\omega})$ generated by above $\beta(x)$ is strictly dyadic branching Brownian motion among Poisson obstacles(DBBP).

3 Some results

Our main purpose is to compare the growth rate Galton-Watson process with DBBP. According to the latter, the following result shows growth order of it. We denote $|Z_t| = \langle 1, Z_t \rangle$.

Theorem 2 (J. Engländer [2])

On a set of full ${\bf P}$ measure

$$\lim_{t \to \infty} \exp\left[-t\left\{\beta_2 - c(d,\nu)(\log t)^{-\frac{2}{d}}\right\}\right] E^{\omega}[|Z_t|] = 1 \qquad (2)$$

as $t \to \infty$, where ω_d is the volume of the *d*-dimensional unit ball, λ_d is the principal Dirichlet eigenvalue of $-\frac{1}{2}\Delta$ on it, and $c(d,\nu) = \lambda_d (\frac{d}{\nu\omega_d})^{-\frac{2}{d}}$.

Outline of the proof of theorem 2

We set $T_t f(x) = E_x[f(Y_t)]$ where $\{Y_t\}_{t \ge 0}$ is $\frac{1}{2} \triangle + \beta$ -diffusion, and $\{T_t\}_{t \ge 0}$ is the semigroup of it. Then

$$E[|Z_t|] = (T_t 1)(x)$$

by The first moment formula. Hence the Feynman-Kac formula implies

$$T_t 1)(x) = \mathbb{E}_x \left[\exp \left\{ \int_0^t \beta(W_s) ds \right\} \right]$$

where $(\{W_t\}_{t\geq 0}, \mathbf{P}_x)$ is d-dimensional Brownian motion. Thus,

$$E^{\omega}[|Z_t|] = \mathbb{E}_x \left[\exp\left\{ \int_0^t \beta_2 - (\beta_2 - \beta_1) \cdot \chi_{K_{\omega}}(W_s) ds \right\} \right] \\ = e^{\beta_2 t} \exp\left[-c(d,\nu)t(\log t)^{-\frac{2}{d}}(1+o(1)) \right]$$
(3)

as $t \to \infty$. To obtain the second equality of (3), we use a theorem of the large time behavior of Brownian motion among Poisson obstacles (cf.[1]).

Theorem 2 tells us that the growth order of $E^{\omega}[|Z_t|]$ is determined by the effect of each particles hitting K_{ω} .

- M. Donsker, S. R. S. Varadhan : Asymptotics for the Wiener sausage, Comm. Pure. Appl. Math., 28, 1975.
- [2] J. Engländer : Quenched law of large numbers for branching Brownian motion in random medium, Ann. Inst. H. Poincaré Probab. Statist., 2007.

Takanao Negishi (Tohoku University)

\star Introduction

For an arbitrary $c\in\mathbb{C}$, we denote by Δ_c the difference operator about c, i.e. $\Delta_c f = f(z+c) - f(z) \text{, where } f\in\mathfrak{M}(\mathbb{C}) \text{ (the set of meromorphic functions on } \mathbb{C}) \text{. Then }, \text{ we consider the difference equation } \Delta_{c_1}\Delta_{c_2}\cdots\Delta_{c_n}f = 0 \cdots (1)(c_1,c_2,\cdots,c_n\in\mathbb{C}).$ We easily see that functions which is a sum of periodic functions $P_{c_1}+P_{c_2}+\cdots+P_{c_n}$ where each P_{c_k} is a c_k -periodic meromorphic function satisfy (1). On the contrary , can all functions satisfying (1) be always represented as such a sum of periodic functions ?

We consider this periodic decomposition problem of meromorphic functions . By applying the method we used in 1-dimensional case , we try to investigate the 2-dimensional case .

\star Historical Background

```
Problem: If a function f : \mathbb{R} \longrightarrow \mathbb{R} satisfies the difference equation

\Delta_{c_1}\Delta_{c_2}\cdots\Delta_{c_n}f = 0\cdots(2) for some c_1, c_2, \cdots, c_n \in \mathbb{R}, then can f be written in the form of f = P_{c_1} + P_{c_2} + \cdots + P_{c_n}

(P_{c_k} : \mathbb{R} \longrightarrow \mathbb{R}; c_k-periodic, 1 \le k \le n)\cdots(3)?
```

Such periodic decomposition problem has already been studied in real analysis so far . However in complex analysis , it has not been considered . It started with some unpublished work of I.Z.Rusza and continued among others .

To study this problem , we have two ways . First we restrict functions to some class . A class \mathcal{F} of real functions is said to have the decomposition property (DP) , if for every $f \in \mathcal{F}$ and $c_1, \cdots, c_n \in \mathbb{R}$, (2) implies that f has a decomposition as (3) with $P_{c_k} \in \mathcal{F}$ $(1 \le k \le n)$.

• Examples of the classes having the DP

• $B(\mathbb{R}) = \{\text{bounded functions}\}$

 $\bullet BC(\mathbb{R}) = \{ \text{bounded continuous functions} \}$

 $\bullet BM(\mathbb{R}) = \{ \text{bounded measurable functions} \}$

 $\bullet B(\mathbb{Z} \to \mathbb{Z}) = \{ \text{bounded } \mathbb{Z} \to \mathbb{Z} \text{ functions} \}$

(In case of $B(\mathbb{Z} \to \mathbb{Z})$, periods c_k s are all integers)

However the class $C(\mathbb{R})=\{\text{continuous functions}\}$ and $\mathbb{R}^{\mathbb{R}}=\{\text{all real functions}\}$ do not have .

The second idea is to complement (2) with other conditions, so that these will be sufficient and necessary for the existence of periodic decompositions. I.Z.Ruzsa showed that if $\frac{c}{d} \notin \mathbb{Q}$, f(x) = x can be split into a *c*-periodic function and a *d*-periodic function. And he showed more generally, if $\frac{c_i}{c_j} \notin \mathbb{Q}$ for $i \neq j$, then (2) implies the decomposition in the form of (3).

\star The Results on the 1 – dimensional case

Before Investigating the general form (1), it is basically important to study the easier equation $\Delta_c \Delta_d f = 0 \cdots (4)$.

Concerning $\left(4\right)$, we can show the following theorem :

```
 \begin{array}{l} \textbf{Theorem 1. Let } \mathfrak{E}(\mathbb{C}) \text{ be the set of entire functions on } \mathbb{C} \text{ and } \mathfrak{M}(\mathbb{C}) \text{ be the set of meromorphic functions on } \mathbb{C} \text{ . The function } f \text{ satisfying the difference equation } \Delta_c \Delta_d f = 0 \text{ can be decomposed as follows }, \\ [A] \text{ if } c \text{ and } d \text{ are linearly independent over } \mathbb{R}, \text{ i.e. } \stackrel{c}{d} \notin \mathbb{R} \text{ ,} \\ (a) \text{ when } f \in \mathfrak{E}(\mathbb{C}) \text{ , then } f = Cz + P_c + P_d \quad (C: const, P_c, P_d \in \mathfrak{E}(\mathbb{C})) \\ (b) \text{ when } f \in \mathfrak{M}(\mathbb{C}) \text{ , then } f = P_c + P_d \quad (P_c, P_d \in \mathfrak{M}(\mathbb{C})) \\ \end{array} \\ \begin{bmatrix} B] \text{ if } \stackrel{c}{d} = \stackrel{c'}{d'} \in \mathbb{Q} \quad (\stackrel{c'}{d'} \text{ is an irreducible fraction }, c = \mu c', d = \mu d'), \\ (a) \text{ when } f \in \mathfrak{E}(\mathbb{C}) \text{ , then } f = P_c + P_d + P_{\mu}z \quad (P_c, P_d, P_{\mu} \in \mathfrak{E}(\mathbb{C}) \text{ ,} \\ P_{\mu} : \mu - \text{periodic} \\ (b) \text{ when } f \in \mathfrak{M}(\mathbb{C}) \text{ , then } f = P_c + P_d + P_{\mu}z \quad (P_c, P_d, P_{\mu} \in \mathfrak{M}(\mathbb{C}) \text{ ,} P_{\mu} : \\ \mu - \text{periodic} \\ \end{bmatrix} \\ \begin{bmatrix} C] \text{ if } \stackrel{c}{d} \in \mathbb{R} \setminus \mathbb{Q} \text{ ,} \\ (a) \text{ when } f \in \mathfrak{E}(\mathbb{C}) \text{ , then } f = Cz + P_c + P_d + \sum_{n=1}^{\infty} (Q_c^n + Q_d^n) \\ (C: const, P_c, P_d, Q_c^n, Q_d^n \in \mathfrak{E}(\mathbb{C}), Q_c^n : c - \text{periodic}, \\ Q_d^n : d - \text{periodic} \\ \end{bmatrix} \end{array}
```

The method of proof is based on the Fourier expansion of entire periodic functions and the Mittag-Leffler theorem . Applying the above theorem again and again , we can gain the following theorem :

 $\begin{array}{l} \textbf{Theorem 2. Let } c_1, c_2, \cdots, c_n \in \mathbb{C} \text{ be pairwise linearly independent over } \mathbb{R} \ . \\ [A] \ \forall f \in \mathfrak{E}(\mathbb{C}) \text{ s.t. } \Delta_{c_1} \Delta_{c_2} \cdots \Delta_{c_n} f = 0 \text{ can be decomposed in the form of} \\ f = a_1 z + a_2 z^2 + \cdots + a_{n-1} z^{n-1} + P_{c_1} + P_{c_2} + \cdots + P_{c_n} \\ (a_j \in \mathbb{C}, P_{c_j} \in \mathfrak{E}(\mathbb{C}) : c_j - periodic, 1 \leq j \leq n) \\ [B] \ \forall f \in \mathfrak{M}(\mathbb{C}) \text{ s.t. } \Delta_{c_1} \Delta_{c_2} \cdots \Delta_{c_n} f = 0 \text{ can be decomposed in the form of} \\ f = P_{c_1} + P_{c_2} + \cdots + P_{c_n} \\ (P_{c_j} \in \mathfrak{M}(\mathbb{C}) : c_j - periodic, 1 \leq j \leq n) \cdots (\mathbf{5}) \end{array}$

Corollary 1. Let $c_1, c_2, \cdots, c_n \in \mathbb{C}$ be pairwise linearly independent over \mathbb{R} . Then, an arbitrary (n-1)-degree polynomial can be decomposed in the form of (5).

Next we try to consider the 2-dimensional case , that is the decomposition problem of meromorphic functions on \mathbb{C}^2 .

Then we apply the ways that we used in 1-dimensional case . So , here we state left or right solution of $\Delta_1 f = h$ $(h \in \mathfrak{M}(\mathbb{C}))$. We denote the set of h(z)'s poles whose real parts are positive by $\{s_n\}$ and non-positive by $\{t_n\}$. By Mittag-Leffler theorem , we can decompose h into two meromorphic functions h_1, h_2 $(h = h_1 + h_2)$ where the polar set of h_1 is $\{s_n\}$ and that of h_2 is $\{t_n\}$. $\Delta_1 f_1 = h_1$ has a solution that has a form of $f_1 = \sum_{n=1}^{\infty} (h(z - n) - p_n) + E_1$ $(p_n: polynomial$, $E_1 \in \mathfrak{E}(\mathbb{C}))$. We call a solution of this type "right solution". And $\Delta_1 f_2 = h_2$ has a solution that has a form of $f_2 = -\sum_{n=0}^{\infty} (h(z + n) - q_n) + E_2$ $(q_n: polynomial$, $E_2 \in \mathfrak{E}(\mathbb{C}))$. We call a solution of this type "left solution". Then $f = f_1 + f_2$ is a solution of $\Delta_1 f = h$.

\star The 2 – dimensional case

In case of entire functions on \mathbb{C}^2 , we can show following theorem :

Theorem 3. [A] Let $e_1 = {}^t(1,0)$, $\alpha = ae_1$ ($a \notin \mathbb{R}$). Then a function $f(z,w) \in \mathfrak{C}(\mathbb{C}^2)$ satisfying $\Delta_{e_1}\Delta_{\alpha}f = 0$ has a representation as $f = Cz + P_{e_1} + P_{\alpha}$ (C: const, $P_{e_1}, P_{\alpha} \in \mathfrak{C}(\mathbb{C}^2)$) [B] Let $e_2 = {}^t(0,2)$. Then a function $f(z,w) \in \mathfrak{C}(\mathbb{C}^2)$ satisfying $\Delta_{e_1}\Delta_{e_2}f = 0$ has a representation as $f = P_{e_1} + P_{e_2}$ ($P_{e_1}, P_{e_2} \in \mathfrak{C}(\mathbb{C}^2)$)

Because the poles of meromorphic functions are not isolated, the meromorphic case is not easy. Furthermore, we have some phenomena which do not arise in 1dimensional case. For example, even the most basic type of difference equation $\Delta_c f = h$ may not have a solution in $\mathfrak{M}(\mathbb{C}^2)$. (Ex. $\Delta_{c1} f = \frac{1}{zw-1}$)

When we consider $\Delta_{e_1}f(z,w)=h(z,w)$, we regard the variable w as a parametor, that is we think the functions f(z;w) and h(z;w) in z varies as the w chenges. Then the polar distribution also depends on the value of w.

Now we restrict the value of w by removing some points from w-plane . First we remove the points w=a s.t. $h(z,a)\equiv 0$ or ∞ . Next we remove the points at which some poles of h(z;w) generate to ∞ . Finally , we remove the points w=a s.t. $I(x,a):=\inf_{y\in\mathbb{R}}|A(x+yi,a)|=0$ for $\forall x\in\mathbb{R}$ where $h(z,w)=\frac{B(z,w)}{A(z,w)}$ $(A,B\in\mathfrak{E}(\mathbb{C}^2))$.

We denote the removal set by $R = R(h; e_1)$. When $w = w_0 \in \mathbb{C} \setminus R$ has a neighborhood $U \subset \mathbb{C} \setminus R$ such that for $\forall y \in \mathbb{R}$ we can find $\exists M \in \mathbb{R}$, $\delta > 0$ s.t. h(z, w) is holomorphic in $\{z \in \mathbb{C} \mid \Re z < M, |\Im z - y| < \delta\} \times U$, then we call the set of all such points right-solvable domain $(\mathcal{S}_R(h; e_1))$. And if h is is holomorphic in $\{z \in \mathbb{C} \mid \Re z > M, |\Im z - y| < \delta\} \times U$, then we call the set of all such points left-solvable domain $(\mathcal{S}_L(h; e_1))$.

If $w_0 \in S_R(h; e_1)$, we can find a right solution of $\Delta_{e_1} f = h$ on $\mathbb{C} \times U$ where U is a neighborhood of w_0 . If $w_0 \in S_L(h; e_1)$, we can find a left solution of $\Delta_{e_1} f = h$ on $\mathbb{C} \times U.I f w_0 \in \mathbb{C} \setminus \mathbb{R}$, then we can decompose h as $h = h_1 + h_2$ where $w_0 \in S_R(h_1; e_1)$, $w_0 \in S_L(h_2; e_1)$. Therefore we can gain some local solutions of $\Delta_{e_1} f = h$. right or left solution can be extended to the right or left solvable domain.

In regard to $\Delta_{e_1}\Delta_{e_2}f=0$, I have not completed the study . However , under some assumptions , a function satisfying $\Delta_{e_1}\Delta_{e_2}f=0$ has a periodic decomposition . For example , if $\exists w_0, w_0+1 \in \mathcal{S}_R(f;e_1)$, then f has a periodic decomposition .

Difficulties of Solving Problems

Kojiro Higuchi

- Computability Theory:
- Computability theory is the study of algorithm, computability and uncomputability.
- A study of difficulties of solving problems:
- What are *problems*? How are difficulties of solving problems *defined*?
- What is the *distribution* of difficulties of solving problems like?
- My study is to clarify the distribution of difficulties of solving Π_1^0 Mass Problems:
- Difficulties are defined by the concept of *Algorithms*.

Mass Problems:=

"Sets of functions on natural numbers".

But we consider mass problems and their elements as follows:

Mass Problems =

"Sets of solutions to corresponding problems", Elements =

"Solutions to the corresponding problem".

Examples:

- 1. Make a complete table of prime numbers, 2, 3, 5, 7, 11, 13, 17, 19, 23, 29, 31, 37, ...
- 2. Calculate π by decimal expansion, 3.141592653589793238462643383279....
- 3. Construct a transcendental number except for e and π .
 - 0.153625035358421993710300427342....
- 4. Find an uncomputable function.
- 5. Find an infinite random sequence, 01101000101110110001011011101100....

6. Extend Mathematics (ZFC) to some complete consistent theory.

Degrees:=

"Difficulties of solving mass problems". P is not more difficult than Q ($P \leq Q$) =

"Some algorithm gives a solution to P by any solution to Q".

The Distribution of Degrees:

- 1. There are infinitely many degrees.
- 2. There are P and Q such that $P \not\leq Q$ and $P \not\geq Q$.
- 3. The top degree and the bottom degree exist.
- -A problem with a computable solution is of the
- -A problem without solutions is of the top degree.

bottom degree.

Algorithms:= Programs written by

"Computer Language + an Instruction Oracle".

For a solution p and a number x, the calculation of $\Phi(p; x)$ is as follows:

- 1. Calculate step by step along the algorithm Φ ,
- 2. Put p(y) into z if "z:=Oracle(y)" appears,
- 3. Let $\Phi(p; x) = Output$, if this calculation halts.

For a solution p, $\Phi(p)\uparrow =$ " Φ never halts on the input p without number x".

 $\Pi_1^0 \text{ Mass Problems:=}$ "Sets of p's satisfying $\Phi(p)\uparrow$ for an algorithm Φ ".

Examples:

- 1. Make a complete table of prime numbers, 2, 3, 5, 7, 11, 13, 17, 19, 23, 29, 31, 37, ...
- 2. Calculate π by decimal expansion, 3.141592653589793238462643383279....
- 5. Find an infinite random sequence, 01101000101110110001011011101100....

6. Extend Mathematics (ZFC) to some complete consistent theory.

The Distribution of Degrees:

- 1. There are infinitely many degrees.
- 2. There are P and Q such that $P \not\leq Q$ and $P \not\geq Q$.
- 3. The top degree and the bottom degree exist.
- -The example 6 is of the top degree.

-A problem with a computable solution is of the bottom degree.

A Problem on Π_1^0 Mass Problems:

Does there exist P and Q such that P < Q and no problem R such that P<R<Q?

Undecidability and weak theory of concatenation

Yoshihiro HORIHATA

Mathematical Institute, D2 E-mail: sa6m31@math.tohoku.ac.jp

• Background of the theory of concatenation

- * In 2005's, A. Grzegorczyk defined the new theory of concatenation denoted by TC, and he proved that TC is undecidable theory([1]).
- * In 2008's, A. Grzegorczyk and K. Zdanowski proved that the theory TC is essentially undecidable. And they left the open question whether the weak arithmetic Q is interpretable in TC.
- * For the above question, in 2009's, V. Švejdar, A. Visser and M. Ganea independently gave a positive answer.

• Main results

- \ast We defined the new theory of concatenation WTC and proved that this system is properly weaker than TC.
- * We proved the theory WTC is Σ_1 -complete, that is, for any Σ_1 sentence φ , if φ is true in the standard model of WTC, then is provable in WTC.
- * We proved that WTC interprets R. This implies that WTC is essentially undecidable. About the converse of this, we conjecture that R interprets WTC.

1. Robinson's weak arithmetic Q 3. Our new theory WTC Oer theory has the following axioms: for each The arithmetic Q is PA-(Induction) whose axioms $x, y, z, u, v \in \{a, b, c\}^*,$ are followings: (W1) $x^{-}\varepsilon = \varepsilon^{-}x = x$. (Q1) $\forall x \forall y (S(x) = S(y) \rightarrow x = y).$ (W2) $x^{(y)}(z) = (x^{y})^{(z)}z$. (Q2) $\forall x(S(x) \neq 0)$. (W3) $x^y = u^v \to \exists w ((x^w = u \land y = w^v) \lor (x = u)$ $u^{\frown}w \wedge w^{\frown}y = v)).$ (Q3) $\forall x (x \neq 0 \rightarrow \exists y (x = S(y))).$ (W4) $\alpha \neq \varepsilon \land x \land y = \alpha \rightarrow x = \varepsilon \lor y = \varepsilon$. (Q4) $\forall x(x+0=x).$ (W5) $\beta \neq \varepsilon \land x^{\frown} y = \beta \rightarrow x = \varepsilon \lor y = \varepsilon.$ (Q5) $\forall x \forall y (x + S(y) = S(x + y)).$ **(W6)** $\gamma \neq \varepsilon \land x \land y = \gamma \rightarrow x = \varepsilon \lor y = \varepsilon$. (Q6) $\forall x(x \cdot 0 = 0)$. (W7) $\alpha \neq \beta \land \beta \neq \gamma \land \gamma \neq \alpha$. (W8) $\forall x(x \sqsubseteq u \rightarrow \bigvee x = v).$ (Q7) $\forall x \forall y (x \cdot S(y) = x \cdot y + x).$ 2. Robinson's very weak arithmetic R 4. Our conjecture Robinson's arithmetic R has the following axioms: for each standard number m, n, Q ТC \bowtie ∇ ∇ (R1) $\bar{m} + \bar{n} = \overline{m+n}$. ?? ⊵ ⊴ R WTC (R2) $\overline{m} \cdot \overline{n} = \overline{m \cdot n}$. (R3) $\bar{m} \neq \bar{n}$ (if $m \neq n$). References [1] A. Grzegorczyk. Undecidability without arithmeti-(R4) $\forall x (x \leq \bar{n} \rightarrow x = \bar{0} \lor \cdots \lor x = \bar{n}).$ zation. Studia Logica, Vol. 79, No. 1, pp. 163-230, 2005. (R5) $\forall x (x < \overline{n} \lor \overline{n} < x).$

Davies' Conjecture for Pseudo-Schrödinger Operators and Its Applications to Penalization Problem

Masakuni MATSUURA*[†] Mathematical Institute, Tohoku University

February 18, 2010[‡]

1 Abstract

We call $(-\Delta)^{\alpha/2} + V(x)$ $(0 < \alpha \le 2, V \in C_0^{\infty}(\mathbb{R}^d))$ a "pseudo-Schrödinger operator" and we consider asymptotic behavior of heat kernels of stochastic processes associated with pseudo-Schrödinger operators. It is well known that Pinchover has solved Davies' conjecture for second-order elliptic operators. (See, Theorem 1.2, [Pin1] or Theorem 1.1, [Pin2]). We derive the result of Davies' conjecture for pseudo-Schrödinger operators from Pinchover's results as follows.

Conjecture 1. Let $k_t(x,y)$ and $\varphi(x)$ be the heat kernel and the ground state associated with pseudo-Schrödinger operator. Then, if H_{α} is subcritical or null-critical, then for every $x, y \in \mathbb{R}^d$,

$$\lim_{t \to \infty} k_t(x, y) = 0.$$

If H_{α} is positive-critical, then for every $x, y \in \mathbb{R}^d$,

$$\lim_{\to\infty}\frac{k_t(x,y)}{k_t(0,0)}=\varphi(x)\varphi(y).$$

Though we can easily prove subcritical cases, we have not completed the proof of critical cases yet.

2 Preliminaries

The following stories are well known.

Let $H_{\alpha} = (-\Delta)^{\frac{\alpha}{2}} + V(x)$ be a pseudo-Schrödinger operator. Then, there is a (C_0) -semigroup $(T_l)_{l \geq 0}$ on $C_{\infty}(\mathbb{R}^d)$, which is generated by H_{α} . Riesz's representation theorem implies the existence of the heat kernel $k_l(x, dy)$ such that

$$T_t v(x) = \int_{\mathbb{R}^d} k_t(x, dy) v(y) \tag{1}$$

for all $v \in C_0^{\infty}(\mathbb{R}^d)$, every $x, y \in \mathbb{R}^d$, and every $t \ge 0$. Then, the heat kernel $k_t(x, dy)$ is the minimal fundamental solution of the initial problem with pseudo-Schödinger operator.

For all $v \in C_0^{\infty}(\mathbb{R}^d)$ and every $x \in \mathbb{R}^d$, $T_t v$ can be represented as

$$T_l v(x) = \mathbb{E}_x \left[e^{\int_0^l V(X_s) ds} v(X_l) \right].$$
⁽²⁾

Here, X_t is the symmetric α -stable process. Further, the range of principal eigenvalues λ_0 of H_{α} is non-negative.

3 Motivation

We would like to solve Davies' conjecture for pseudo-Schrödinger operators in order to solve penalization. Let us suppose the following assumption.

Assumption 2. If H_{α} is positive critical, then

$$\lim_{t \to \infty} \frac{k_t(x, y)}{k_t(0, 0)} = \varphi(x)\varphi(y).$$
(3)

for every $x, y \in \mathbb{R}^d$.

[†]Web:http://www.math.tohoku.ac.jp/~sa9d10/

 $\ensuremath{^{\ddagger}\text{Poster}}$ session, Global COE symposium "Weaving Science Web beyond Particle-Matter Hierarchy".

Then, there is a limit distribution \mathbb{P}_0^{φ} such that

$$\lim_{t \to \infty} \frac{\mathbb{E}_0 \left[e^{\int_0^t V(X_u) du} \Sigma \right]}{\mathbb{E}_0 \left[e^{\int_0^t V(X_u) du} \right]} = \mathbb{E}_0^{\varphi} \left[\Sigma \right]$$
(4)

for all $\Sigma \in \mathscr{F}_s$. Indeed, if we define

$$\begin{split} \mathbb{P}_{0}^{\varphi} &:= \frac{\varphi(X_{s})}{\varphi(X_{0})} e^{\lambda_{0}s + \int_{0}^{s} V(X_{u}) du} \mathbb{P}_{0}, \\ \text{hen as } t \to \infty, \\ &= \frac{\frac{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} \Sigma \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} \right]} \\ &= \frac{\mathbb{E}_{0} \left[\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} \Sigma \right] \right] \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} 2\mathbb{E}_{0} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \circ \theta_{s} \right] \right]} \\ &= \frac{\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{0} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \circ \theta_{s} \right] \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right] \right]} \\ &= \frac{\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right] \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right] \right]} \\ &= \frac{\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right] \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right] \right]} \\ &= \frac{\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right] \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right]} \right]} \\ &= \frac{\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right] \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right]} \right]} \\ &= \frac{\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right]} \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} V(X_{u}) du} 2\mathbb{E}_{X_{s}} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right]} \\ &= \frac{\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{0} \left[e^{\int_{0}^{s} V(X_{u}) du} 2\mathbb{E}_{0} \left[e^{\int_{0}^{t} - s V(X_{u}) du} \right]} \right]}{\mathbb{E}_{0} \left[e^{\int_{0}^{t} - s V(X_{u}) du} 2\mathbb{E}_{0} \left$$

 $k_{t-s}(X_s, y)$

 $\frac{k_{l-s}(X_s,y)}{k_l(0,0)} \to \varphi(X_s)\varphi(y) \mathbb{P}_0\text{-a.s.}$

as $t \rightarrow \infty$ in the last computation.

We call such a problem "Feynman-Kac penalization".

- [ABJ] Anker, J.P., Bougerol, P., and Jeulin, T., *The infinite Brownian loop* on a symmetric space, Rev. Mat. Iberoamericana 18, pp. 41 – 97, 2002.
- [Dav1] Davies, E.B., Non-Gaussian aspects of heat kernel behavior, J. London Math. Soc. (2) 55, pp. 105 – 125, 1997.
- [Dav2] Davies, E.B., *Heat Kernel and Spectral Theory*, Cambridge Univ. Press, Cambridge, 1989.
- [Jac1] Jacob, N., *Pseudo-differential operators and Markov processes, Volume I*, Imperial College Press, 2001.
- [Jac2] Jacob, N., *Pseudo-differential operators and Markov processes, Volume III*, Imperial College Press, 2005.
- [Pin1] Pinchover, Y., Large Time Behavior of the Heat Kernel and the Behavior of the Green Function near Criticality for Nonsymmetric Elliptic Operators, Journal of Functional Analysis 104, pp. 54 – 70, 1992.
- [Pin2] Pinchover, Y., On Davies' conjecture and strong ratio limit properties for the heat kernel, Advanced Studies in Pure Mathematics XX, pp. 1 – 13, 2005.

^{*}Mail: sa9d10@math.tohoku.ac.jp

Torsion points of abelian varieties with values in infinite extension fields

Yuken Miyasaka

Mathematical Institute, Tohoku University

(E-mail : sa7m27@math.tohoku.ac.jp)

We consider a problem "the torsion part $A(L)_{tor}$ of an abelian variety A over a number field or p-adic field K wih values in some infinite extension L of K is finite or infinite ?". On the other hand, the periodic points of an arithmetic dynamical system have some relationship with the torsion points of abelian varieties, and there also is a problem related to the finiteness. We compare the side of the variety with that of the dynamics in this problem.

SIDE OF VARIETY

K: number field, A/K: abelian variety, L/K: extension A(L): set of K-rational points of A with values in L. $A(L)_{tor} := \{P \in A(L) : nP = 0 \text{ for } \exists n \in Z\}$

Well-Kown Fact L/K: finite extension $\Rightarrow \#A(L)_{\text{tor}} < \infty$

• It is true even if we replace K with a p-adic field.

UNIFORMLY BOUNDED

Thm. (Mazur-Morel) -

A = E: elliptic curve over number field K L/K: finite extension $\exists C = C(d)$: constant s.t. $\#E(L)_{\text{tor}} < C$.

Mazur proved it when K = Q, and Morel proved in the general case. Their proof of this theorem is very difficult !
It is an open problem in the case of abelian varieties.

Problem

- A: abelian variety over number field or p-adic field KIf L is an infinite extension field of K,
- $A(L)_{tor}$ is finite or infinite?

GLOBAL AND LOCAL RESULTS

In the global case, there are many results for this problem. We pick up the following Ribet's result.

Thm. (Ribet) *K* : number field, $L := K(\mu_{\infty}) \Rightarrow \#A(L)_{\text{tor}} < \infty$

Ribet proved this theorem by reducing it to the following local result which is proved by Imai and Serre.

Thm. (Imai-Serre) K: p-adic field, $L := K(\mu_{p^{\infty}})$ A: abelian variety with ordinary reduction $\Rightarrow \#A(L)_{tor} < \infty$

Recently, This local result is refined to more large field L by Ozeki.

MANIN- MUMFORD CONJECTURE

Finally we note that for curves with a higher genus, it is known the following Manin- Mumford conjecture, proved by Raynaud.

Thm. (Raynaud) K: alg. closed field of char. 0, C/K: curve of genus > 1 J: Jacobi variety of $C \Rightarrow C \cap J(K)_{tor} < \infty$

SIDE OF DYNAMICS

$$\begin{split} &K: \text{number field, } f(X) \in K(X), \ L/K: \text{extension} \\ &P \in \mathrm{P}^1(L) \text{ is periodic, if } f^{(n)}(P) = f \circ \cdots \circ f(P) = P. \\ &\Pi_f(L) := \{P \in \mathrm{P}^1(L): f^{(m)}(P) \text{ is periodic } {}^\exists m \in \mathbf{Z} \} \end{split}$$

Thm. (Nothcott) L/K: finite extension $\Rightarrow \Pi_f(L)$: finite

• There are not results for a *p*-adic field.

UNIFORMLY BOUNDED (DYNAMICAL VERSION)

Conj. (Morton-Silverman)

L : finite extension of number field K, $f(X) \in K(X)$ Then $\exists C$: constant depending only on L and $\deg(f)$ s.t. $\#\Pi_f(L) < C$.

• This conjecture includes the Mazur-Morel's theorem stated the side of variety, thus it is considered that it is difficult to prove it. it is not known even when $f(X) = X^2 + c$.

Problem (dynamical version) -

f(X): rational function over number field or *p*-adic field If *L* is an infinite extension field of *K*, $\Pi_f(L)_{tor}$ is finite or infinite?

LATTÉS MAP

The Lattés map f_E is one of the most important rational maps in the study of this problem. It is obtained by the projection of the multiplication by m map on an elliptic curve Esuch as the following:

The Lattes map has the property $\#\Pi_{f_E}(L) \leq \#E(L)_{tor}$ for an extension L of K. Therefore the results in the side of the variety give some examples of the finiteness problem in the side of the dynamics.

DYNAMICAL MANIN- MUMFORD CONJECTURE

Finally, we state the dynamical system version of the Munin-Mumford conjecture. This conjecture is also an open problem.

Conj. (Zhang) -

 $\begin{array}{l} K: \text{alg. closed field of char. 0, } \phi: \text{morphism on } \mathbf{P}^N(K) \\ X: \text{variety in } \mathbf{P}^N \\ \Pi_\phi(K) \cap X \subset X: \text{dense} \Leftrightarrow X: \text{pre-periodic variety} \end{array}$

Maximum principle for a biological model related to the motion of amoebae

Harunori Monobe (Tohoku University)

February , 2010

1 Model of amoebae motion

Tamiki Umeda (Kobe University) proposed a biological model of amoebae considering the motion and the chemical reaction in the body. We modify Umeda's model from a mathematical viewpoint, and analyze the following model:

$$(\mathbf{P}) \begin{cases} u_t = \Delta u + k_1 w(t) - k_2 u & x \in \Omega(t), \ t > 0, \\ u = 1 + A \kappa + B V & x \in \partial \Omega(t), \ t > 0, \\ V = -\nabla u \cdot \mathbf{n} + g(u) w(t) & x \in \partial \Omega(t), \ t > 0, \\ u = \phi & x \in \Omega(0), \end{cases}$$

where

$$\begin{split} &\Omega(t): \text{domain of } \mathbf{R}^2 \ , \ \partial\Omega(t): \text{boundary of } \Omega(t), \\ &\kappa = \kappa(x,t): \text{curvature of } \partial\Omega(t), \\ &V = V(x,t): \text{normal velocity of } \partial\Omega(t), \\ &\mathbf{n} = \mathbf{n}(x,t): \text{outward normal vector to } \partial\Omega(t), \\ &g(\cdot): \text{smooth function}, \ w(t) = C_0 - \int_{\Omega(t)} u \, dx \\ &k_1, k_2, A, B, C_0: \ \text{positive constant}, \ \phi(x): \text{initial data}. \end{split}$$

In (P), u(x, t) is an unknown function depending on space and time, and the function represents a F-actin which assumes the role of bones and muscles in the body of amoebae. w(t) represents a G-actin which is produced by F-actin. This model is mainly constructed by the law of conservation of mass and the relationship between G-actin, F-actin and the normal velocity of the boundary $\partial \Omega(t)$:

 $u_t + \operatorname{div}(uv) = k_1 w(t) - k_2 u, \quad V = v \cdot \mathbf{n} + g(u)w(t).$



Figure 1: Chemical reaction and Motion of amoebae

2 Motivation

In Umeda's model, chemotactic substance gives a moebae a positive effect [Figure 2]. Therefore, the function g which represents the activity of a chemical reaction is positive. However, there are some chemotactic substance which also give a moebae a negative effect [Figure 3]. Now we consider the situation that there are such a substance around a moebae, and we try taking g as a negative function. Then we expect a moebae to become smaller. However, it is natural for a moebae not to become smaller than a certain size. We predict the sign of the function g on the basis of a symptotic behavior in some special conditions.



Figure 2: Case of positive effect



Figure 3: Case of negative effect

3 Main Result

Let |x| = r, v(r,t) = u(x,t) in (P), then (P) becomes

$$(\text{RP}) \begin{cases} \frac{\partial v}{\partial t} = v_{rr} + v_r/r + k_1 w(t) - k_2 v & r \in (0, s(t)), \ t > 0, \\ v = 1 + A/s(t) + B\dot{s}(t) & r = s(t), \ t > 0, \\ \dot{s}(t) = -v_r + g(v)w(t) & r = s(t), \ t > 0, \\ v_r = 0 & r = 0, \ t > 0, \\ v(r, 0) = \phi & r \in (0, s(t)), \end{cases}$$

where s(t) is an unknown function depending only on time, and the function represents the radius of the circular domain which is the body of amoebae.

Definition 1. A pair (v, s(t)) is said to be a solution of (RP) if there exists T > 0 such that

$$(v, \ s(t)) \in C^{2+\alpha, \ (2+\alpha)/2} \left(\bigcup_{0 \le t \le T} [0, s(t)] \times \{t\} \right) \times C^{(3+\alpha)/2}([0, T])$$

satisfies (RP) for some $\alpha \in (0, 1)$.

Assumption 1. We assume that the initial data satisfy the following condition:

$$s(0) \in \left(0, \frac{-A + \sqrt{A^2 + 4C_0/\pi}}{2}\right), \quad \max_{r \in [0, s(0)]} \phi = \phi(s(0), 0),$$
$$\phi(r, 0) > 0, \quad C_0 - 2\pi \int_0^{s(0)} r\phi \, dr > 0.$$

Then we have the following results:

Lemma 1. Let (v, s(t)) be a solution of (RP) and the initial data satisfies Assumption 1. Moreover, g(v) is depend only on time and satisfies

$$0 > g > -\frac{1}{BC_0}.$$

Then v > 0, w(t) > 0, $\dot{s}(t) < 0$ for all $t \in [0, T]$.

Theorem 1. Suppose that the initial data satisfy the same assumption as in Lemma 1 and

$$\frac{k_1 C_0}{k_2} < \phi(s(0), 0).$$

Then

$$\max_{Q_T} v = \max_{0 \le t \le T} v(s(t), t).$$

Moreover, if $T = \infty$, then $s(t) \to 0$, $\max_{r \in [0,s(t)]} v \to \infty$ as $t \to \infty$.

Remark 1. From Lemma 1 and Theorem 1, if g is negative at any time, we may construct an extinction solution. Therefore, if the model (P) have validity from a biological viewpoint, we expect the sign of g to be non-negative when amoebae is smaller than a certain size.

References

H. Monobe, Existence of solutions for a mathematical model related to the motion of an amoeba (preprint).

Asymptotic Behavior of Solutions to the Drift-Diffusion Equation in the Whole Space

Masakazu YAMAMOTO (sa5m270math.tohoku.ac.jp) Mathematical Institute, Graduate School of Science, D3

February 18, 2010

Abstract

We consider the asymptotic behavior of the solution to the Cauchy problem for the Nernst–Planck type drift-diffusion equation arising from the plasma dynamics model. For our problem, it is already proved that the time global existence and decay of the solution. We also show an asymptotic expansion of the solution as $t \to \infty$.

1 Introduction

We study the following Cauchy problem for the drift-diffusion equation.

(1)
$$\begin{cases} \partial_t u - \Delta u + \nabla \cdot (u \nabla \psi) = 0, \quad t > 0, \quad x \in \mathbb{R}^3, \\ -\Delta \psi = -u, & t > 0, \quad x \in \mathbb{R}^3, \\ u(0, x) = u_0(x) \ge 0, & x \in \mathbb{R}^3. \end{cases}$$

The drift-diffusion equation is the model of a plasma dynamics. The unknown functions u = u(t, x) and $\psi = \psi(t, x)$ denote the density of charges and the potential of statistic electric field, respectively.

The well-posedness and the global existence of solution are already proved([5]). The mass conservation and L^p -decay estimate for the solution to (1) with large initial data were derived. Moreover, the effect from the non-linear part decays faster than the top term from the linear part([4]). Namely, the estimate $||u(t) - MG(1+t)||_p = o(t^{-\gamma})$ $(t \to \infty)$ holds for $1 \leq p \leq \infty, \ \gamma = \frac{3}{2}(1 - \frac{1}{p})$, where $G(t, x) := (4\pi t)^{-3/2}e^{-|x|^2/(4t)}$ and $M := \int_{\mathbb{R}^3} u_0(y) dy$.

2 Main Result

Now, our main concern here is the second asymptotic expansion of the solution. For $u_0 \in L_2^1(\mathbb{R}^3) \cap L^{\infty}(\mathbb{R}^3)$, where $L_2^1 := \{f \in L^1 \mid |x|^2 f \in L^1\}$, we introduce the following functions:

$$\begin{split} V_0(t,x) &:= G(t,x) \int_{\mathbb{R}^3} u_0(y) dy, \quad V_1(t,x) := \nabla G(t,x) \cdot \int_{\mathbb{R}^3} y u_0(y) dy \\ J(t,x) &:= \int_0^t \nabla e^{(t-s)\Delta} \cdot (V_0 \nabla (-\Delta)^{-1} V_0)(s) ds, \\ K(t,x) &:= -\frac{1}{3} \log(1+t) \Delta G(t,x) \int_{\mathbb{R}^3} y \cdot (V_0 \nabla (-\Delta)^{-1} J \\ &+ J \nabla (-\Delta)^{-1} V_0)(1,y) dy. \end{split}$$

Then the following estimate holds.

Theorem 1 ([9]) Let
$$u_0 \in L_2^1(\mathbb{R}^3) \cap L^{\infty}(\mathbb{R}^3)$$
 and $u_0 \ge 0$. Then for $1 \le p \le \infty, \ \gamma = \frac{3}{2}(1 - \frac{1}{p})$, the following estimate holds
(2) $\|u(t) - V_0(t) - V_1(t) - J(t) - K(t)\|_p = O(t^{-\gamma - 1})$ as $t \to \infty$.

Moreover, the functions J and K satisfy $J, K \neq 0$.

We should denote that the functions V_0, V_1, J and K satisfy the following equalities for any $\lambda > 0$:

(3)
$$\lambda^{3}V_{0}(\lambda^{2}t,\lambda x) = V_{0}(t,x),$$
$$\lambda^{4}V_{1}(\lambda^{2}t,\lambda x) = V_{1}(t,x), \quad \lambda^{4}J(\lambda^{2}t,\lambda x) = J(t,x).$$

Theorem 1 states that the asymptotic expansion of the solution to (1) contains logarithmic term at the rate $t^{-\frac{3}{2}(1-\frac{1}{p})-1}$.

There are very much related model called the Navier–Stokes equation and the Keller–Segel equations. They appears in a model for the incompressible fluid flow and the chemotaxis respectively. For the Navier–Stokes equations, an asymptotic behavior of the solution was considered ([1, 2]). In the case of the Keller–Segel equation, it was shown that there exists a special term likes J(t) in an asymptotic expansion of the solution ([3, 6]). Moreover, in the even-dimensional cases, the asymptotic expansion of the solutions contains a logarithmic term at the rate $t^{-\frac{n}{2}(1-\frac{1}{p})-\frac{n+1}{2}}$ ([8]).

3 Outline of the Proof

In order to prove Theorem 1, we see the first asymptotic expansion estimate for the solution.

Proposition 2 (cf.[7]) Under the same assumption as in Theorem 1, the following estimate holds for $1 \le p \le \infty$ and $\gamma := \frac{3}{2}(1 - \frac{1}{n})$;

 $||u(t) - V_0(1+t) - V_1(1+t) - J(1+t)||_p \le C(1+t)^{-\gamma - 1}\log(2+t).$

To have an estimate for $\nabla \psi$, we prepare the following Sobolev type estimate. Lemma 3 (Hardy–Littlewood–Sobolev's inequality) Let $f \in L^p(\mathbb{R}^3)$ for $1 . Then <math>\nabla (-\Delta)^{-1} f \in L^{p_*}(\mathbb{R}^3)$ and $\|\nabla (-\Delta)^{-1} f\|_{p_*} \leq C \|f\|_p$ for $3/2 < p_* < \infty$ with $\frac{1}{p_*} = \frac{1}{p} - \frac{1}{3}$.

The Cauchy problem (1) is equivalent to the following integral equation:

(4)
$$u(t) = e^{t\Delta}u_0 + \int_0^t \nabla e^{(t-s)\Delta} \cdot \left(u\nabla(-\Delta)^{-1}u\right)(s)ds, \quad t > 0, \ x \in \mathbb{R}^3,$$

where $\{e^{t\Delta}\}_{t\geq 0}$ is the heat semigroup. A combination of Proposition 2 and Lemma 3 immediately gives

$$i\nabla(-\Delta)^{-1}u \sim V_0\nabla(-\Delta)^{-1}V_0 + V_0\nabla(V_1 + J) + (V_1 + J)\nabla(-\Delta)^{-1}V_0.$$

Then the nonlinear part on the right hand side of (4) converges to

$$\int_{0}^{5} \nabla e^{(t-s)\Delta} \cdot \left(u\nabla(-\Delta)^{-1}u\right)(s)ds \\ \sim \int_{0}^{t} \nabla e^{(t-s)\Delta} \cdot \left(V_{0}\nabla(-\Delta)^{-1}V_{0}\right) \\ + V_{0}\nabla(V_{1}+J) + (V_{1}+J)\nabla(-\Delta)^{-1}V_{0})(1+s)ds \\ \sim \sum_{|\beta|=1} \nabla^{\beta}\nabla G(t,x) \cdot \int_{0}^{t} \int_{\mathbb{R}^{3}} (-y)^{\beta}(V_{0}\nabla(-\Delta)^{-1}V_{0}) \\ + V_{0}\nabla(V_{1}+J) + (V_{1}+J)\nabla(-\Delta)^{-1}V_{0})(1+s)ds,$$

where we use the Taylor expansion. The right hand side of (5) contains a logarighmic term K. Indeed, by the scaling arguments (3), we have

$$\begin{split} &\sum_{|\beta|=1} \nabla^{\beta} \nabla G(t,x) \int_{0}^{t} \int_{\mathbb{R}^{3}} (-y)^{\beta} (V_{0} \nabla (-\Delta)^{-1}J + J \nabla (-\Delta)^{-1}V_{0})(1+s,y) dy ds \\ &= \sum_{|\beta|=1} \nabla^{\beta} \nabla G(t,x) \int_{0}^{t} (1+s)^{-1} \int_{\mathbb{R}^{3}} (-(1+s)^{-1/2}y)^{\beta} (V_{0} \nabla (-\Delta)^{-1}J \\ &\quad + J \nabla (-\Delta)^{-1}V_{0})(1,(1+s)^{-1/2}y)(1+s)^{-3/2} dy ds \\ &= -\log(1+t) \sum_{j=1}^{3} \partial_{j}^{2} G(t,x) \int_{\mathbb{R}^{3}} \eta_{j} (V_{0} \partial_{j} (-\Delta)^{-1}J + J \partial_{j} (-\Delta)^{-1}V_{0})(1,\eta) d\eta \end{split}$$

where we put $(1+s)^{-1/2}y = \eta$ and use the relation $\int_{\mathbb{R}^3} \eta_j (V_0 \partial_k (-\Delta)^{-1}J + J\partial_k (-\Delta)^{-1}V_0)(1,\eta)d\eta = 0$ $(j \neq k)$ in the second equality. This term gives the desired function K since the following relation holds:

$$\int_{\mathbb{R}^3} \eta_j (V_0 \partial_j (-\Delta)^{-1} J + J \partial_j (-\Delta)^{-1} V_0) (1, \eta) d\eta$$

= $\frac{1}{3} \int_{\mathbb{R}^3} \eta \cdot (V_0 \nabla (-\Delta)^{-1} J + J \nabla (-\Delta)^{-1} V_0) (1, \eta) d\eta \quad (j = 1, 2, 3).$

Applying Perseval's equality, we can confirm that $\int_{\mathbb{R}^3} \eta \cdot (V_0 \nabla (-\Delta)^{-1} J + J \nabla (-\Delta)^{-1} V_0)(1, \eta) d\eta \neq 0.$

- [1] Carpio, A., SIAM J., Math. Anal. 27 (1996), 449–475.
- [2] Fujigaki, Y., Miyakawa, T., SIAM J. Math. Anal. 33 (2001), 523–544.
- [3] Kato, M., Differential Integral Equations 22 (2009), 35-51.
- [4] Kawashima, S., Kobayashi, R., Funkcial. Ekvac. 51 (2008), 371–394.
- [5] Kurokiba, M., Ogawa, T., J. Math. Anal. Appl. 342 (2008), 1052-1067.
- [6] Nagai, T., Yamada, T., J. Math. Anal. Appl. 336 (2007), 704-726.
- [7] Ogawa, T., Yamamoto, M., Math. Models Methods Appl. Sci. 19 (2009), 939–967.
- [8] Yamada, T., Hiroshima Math. J. **39** (2009), 363–420.
- [9] Yamamoto, M., RIMS Kôkyûroku Bessatsu B15 (2009), 189–208.

A study of the idea of systematic knowledge: On the relation between nature and spirit in the organizational view of nature

Abstract

In the studies of the humanities and science which have become increasingly technical, sophisticated and interdisciplinary, the idea of the assimilation of literature and science attracted growing academic interest in recent years. This idea naturally requires not aggregate but systematic knowledge accomparied by a philosophy that addresses the question of what it is that we conceive as knowledge. It was the German philosophy as it evolved during the 18-79th century, which, based on works on systematic or speculative knowledge, attempted to establish an organizational view of nature becoming known under the notion of natural philosophy. Here, nature and spirit, which were traditionally conceived as divided notions within the mechanical view on nature, were now viewed in a new light as something unified through life. Although modern natural philosophy thus the relation between nature and spirit in this particular theoretical notion as well as the systematic knowledge itself which attempts to apprehend such an organizational whole.

1. Transition of the philosophical view on nature

(1) The teleological view on nature (e.g. Aristotle)

- Nature
- = The essence of things which have a source of movement in themselves
- Unidentifiable materials are identified through their form.
 Formal ("eidos") ⇔ "Telos" of the self-making contained in nature.
- Everything from the individual to the world as a whole has a necessary
- purpose.
 An ordered cosmos held by "unmoved movers (kinoun akineton) = God" "Not only is immaterial being or absolute mind logically prior to nature, but the differentiation of mind into minds is prior to nature also."

(Collingwood[1960],p.90)

Order based on the

Christian world view

God

Nature

Nature

natura naturans

(2) The mechanistic view on nature (e.g. Renaissance, Modern science)

Descartes

Substance = God, spirit, matter Spirit = res cogitans Matter = res extensa Nature as matter ⇔ causality as mechanism of transmission of movements (mechanical causality)

 ■ Principle of causality in the mechanistic view on nature
 ⇒ The concept of "causa finalis" and "causa formalis" is not valid. Only "causa efficiens" formuus is valid.

(3) The organizational view on nature (e.g. Romanticism, German Idealism)

 Opposition to the mechanistic view on nature Nature isn't mechanical but living.

 Nature as subject in the work of Schelling While nature produces itself, it is produced from natura naturata

2. Epistemology as the basis of each view on nature

(1) The teleological view on nature

natural knowledge on the schema "hyl #" - "eidos" "eideische = qualitative Naturerkenntnis"(cf. Hirschberger [1953], p.33)

(2) The mechanistic view on nature

natural knowledge on the dualism of subject and object "mechanistische = quantitative Naturerkenntnis"(cf. ibid.)

(3) The organizational view on nature

systematic and holistic knowledge of nature and spirit as life To grasp the whole not as an aggregation but a system.

3. The relation between nature and spirit in the organizational view on nature (part. on Schelling and Hegel)

Schelling's natural Philosophy

The concept of "living nature"

I he concept of "inving nature" "So long as I myself an identical with Nature, I understand what <u>a living nature</u> is as well as I understand my own life; I apprehend how this <u>universal life of</u> <u>Nature</u> reveals itself in manifold forms, in progressive developments, in gradual approximations to freedom. As soon, however, as I separate myself, and with me everything ideal, from Nature, nothing remains to me but <u>a dead</u> <u>edject</u>, and I case to comprehend how a *life utsikal* me can be possible." (Ideas for a philosophy of nature as introduction to the study of this science (1797), p.36)

Nature and spirit ("I") are ideally identical as something living.

The concept of Organization

The concept of Organization "Every organic product carries the reason of its existence in *itself*, for it is cause and effect of itself. No single part could arise except in this whole, and this whole itself consists only in the *interaction* of the parts [...]. The organism, however, is not mere appearance, but is *itself* object, and indeed an object subsisting through itself, in itself whole and indivisible, and because in it the form is inseparable from the matter, the *origin* of an organism, as such, can no more be explained mechanically than the origin of matter itself." (ibid, p. 31)

Hegel

@ "Spirit is higher than nature."

a"Spirit is higher than nature."
"if [...] that infinite expansion and this infinite withdrawal into itself, are completely one [...], spirit is higher than nature. For if nature is absolute self-intuition and the actuality of the infinited differentiated meditation and development [Entfaltung], then spirit, which is the intuition of itself as itself - or absolute cognition - is, in the withdrawal of the universe into itself. John the scattered totality of this multiplicity which it [Le spirit] encompasses, and the absolute ideality of this same multiplicity, in which it nullifies this separateness and reflects it into itself as the unmediated point of unity of the infinited work of the infinite dense point of unity of the infinite of the infinite dense point of unity of the infinite dense point of unity

- Hegel also understands nature as something living and subjective, especially the living process ("Lebensprozeß"). Unlike on Schelling, spirit is not identical with Nature on Hegel.
 - Hegel thinks that nature itself and conceptualization of nature "completely are

@ "Impotence of nature (Ohnmacht der Natur)"

"It is not only that in nature the play of forms has unbounded and unbridled It is not only that in nature the play or forms has <u>unbounded and unbrated</u> contingency, but that each shape by <u>itself</u> is devided if its holfung. Life is the highest to which nature drives in its determinate being, but as merely natural loca, life is submerged in the irrationality of externality, and the living individual is bound with another individuality in every moment of its existence, while spiritual manifestation contains the moment of a free and universal relation of spirit to itself."

(The Philosophy of Nature, § 248 Remark)

"<u>The impotence of nature</u> is to be attributed to its only being able to maintain the determinations of the Notion in an abstract manner, and to its exposing the foundation of the particular to determination from without." (ibid., § 250)

⇒ Hegel's theory of "the impotence of nature" is a radical form of the idealistic

Indeed Hegel also considers nature as an organic life, but it is classified or opposed to nature and spirit according to the degree of the realization of the idea of freedom. Hegel's view on nature does not represent a return to the mechanistic view which implies a human supremacy over nature. But he also attempts to prove the possibility that the spirit forms itself freely and voluntarily in the necessity and autonomy accide of nature is emphasized.

4. Spirit and nature on Whitehead

Whitehead's "Process Philosophy" or "Philosophy of organism" @ Whitehead: An exponent of the organizational view on nature i . he 20[±] ntury

- We are merely endeavoring to exhibit the type of relations which hold between the entities which we in fact perceive as in nature." (*The Concept of Nature*) \Leftrightarrow Getting ride of the relation between subject and object, and metaphysical arguments
- @ The first category of Existence "Actual Entities (also termed Actual Occasions), or Final Realities, or Res Verae." (Process and Reality, p.22)
 The reformed subjectivist principle
- Whitehead does not suppose the" immovable mover", an absolute spirit, or subjective substance on the basis of the organism, but he considers the world in its solidarity through actual entities and occasions.

Conclusion

- Because the modern organizational view of nature (part. Schelling) takes its origin from the humanism of romanticism (elevation from the finite to the infinite), identity between nature and spirit is described from the element of intellectual intuition, and it is not finally problematic how is a correspondence between the contructed (experienced) and true nature. ~ However it is also said that this point of view is a source of today's environmental ethics which criticizes anthropocentrism.
- It is the most important point on the organizational view on nature that nature and spirit only statically unite, but that both reciprocally and organically relate and creatively evolve. Therefore this view doesn't fall into anthropocentrism like in the mechanistic view even if spirit oft might have the initiative in observations, speculation and so on. The idea of Systematic knowledge is also a cosmology of nature and spirit.

- Georg Wilhelm Friedrich Hegel, Gesammelte Werke, in Verbindung mit der Deutschen Forschungsgemeinschaft, Hamburg 1966 ff.
 Schellings Werke, Nach der Originalausgabe in neuer Anordnung herausgegeben von Manfred Schröter, München 1927 ff.
 G.W.F. Hegel, Political Writings, edited by Laurence Dickey and H. B. Nisbet, traniated by H. B. Nisbet, Cambridge University Tress, 1999.
 Hegel S Philosophy of nature, edited and translated by M. J. Petry, Allen & Unwin, E. W. I. Schlene University Tress, 1999.

- London 1970. ¹⁹ J. ¹⁰ J. ¹⁰ J. ¹⁰ L. ¹⁰ L.

Medical technology and surrogate decision-making

1. INTRODUCTION

With regard to medical technology it is bioethics which emphasizes the principle of respecting the autonomy of the individual. Yet, in the case of doubtfully autonomous or non-autonomous patients the principle of bioethics may not be granted or even be absent. In this presentation, I wish to discuss the framework of standards for surrogate decision-making in order to make decisions on behalf of individuals incompetent to seek decisions on their own.

2. INCOMPETENT PATIENTS AND STANDARDS FOR SURROGATE DECISION-MAKING

In the case of some patients only insufficient autonomy is granted or may even be absent.

Ex. persistent vegetative state patients, dementia patients

 \rightarrow formerly competent patients

newborns suffering a serious illness

 \rightarrow never competent patients

Two types of standards for surrogate decision-making in Bioethics

[a] The Pure Autonomy Standard

(-The Substituted Judgment Standard)

[b] The Best Interests Standard

3. The pure autonomy standard and its problems

This standard requires the surrogate decision-maker to formulate a decision for formerly autonomous, now incompetent patients based on the patient's prior autonomous preferences (precedent autonomy).

- → The decision might be made based on values held by the patient that are little relevant to it
 - The precedent autonomous preferences might conflict with the interests of the patient
 - The preferences or choices of the reasonable (competent) people could be different from the preferences or values of the incompetent people

4. The best interests standard and its limits

This standard requires the surrogate decision-maker to determine the option that result in the highest net benefit for incompetent patient by evaluating burdens and benefits of the

Haruka HIKASA (Philosophy D3, Tohoku University)

available options.

- → The decisions might fail to reflect the patient's subjective benefits, because surrogate decision-makers evaluate benefits and burdens, more or less, dependently on objective judgment (QOL chosen by a reasonable person)
 - Whether the burdens should be limited to physical pain and suffering

5. Towards a new method of surrogate decision-making

 In evaluating the result of medical interventions for incompetent patients, surrogate decision-makers must attempt to ascertain the patient's *present* point of view

 \rightarrow cognitive science and neuroscience might be useful for interpretations of this present point of view

- 2) Surrogate decision-makers take the patient's present benefits (preferences, experiential interests) and the precedent autonomy (formerly preferences, values, perspectives) into consideration in order to make decisions on behalf of formerly competent patients.
 → the need to distinguish the patient from the person who he/she used to be, while at the same time regarding the patient as a person who continues to live a human life
- 3) With regard to 1) and 2), not only the physical but also the mental (spiritual) pain and suffering of the incompetent patient should be included in the assessments.
- 4) Surrogate decision-makers assume that some medical interventions will be extreme terror on patients who are unable to understand the reasons for these burdens (ex. invasive or immobilizing treatments)
- 5) Protecting and promoting the abilities which still held by incompetent patients. These abilities must be regarded as important.

REFERENCES

Beauchamp, Tom L., and Childress, James F. 2001: Principles of Biomedical Ethics. 5th ed., New York,
Oxford University Press. Dresser, Rebecca and Whitehouse, Peter J. 1994: "The incompetent patient on the slippery slope", Hastings Center Report, Vol.24, No.4, pp.6-12. Dresser, Rebecca. 1995: "Dworkin on Dementia. Elegant Theory, Questionable Policy", in:Hastings Center Report, Vol.25, No.6, pp.32-38.
Dworkin, Ronald. 1994: Life's Dominion, New York. Jonsen, Albert R. and Siegler, Mark and Winslade, William J. 2002: CLINICAL ETHICS: A Practical Approach to Ethical Decisions in Clinical Medicine. 5th ed., MeGraw-Hill Medical Publishing.

Hume's empiricism and the experimental method of reasoning

SUGAWARA Hiromichi (Philosophy, D1) Graduate School of Arts and Letters, Tohoku University

Introduction

Hume intended to introduce the experimental* method of reasoning into moral subject, a method of which Isaac Newton successfully provides a secure foundation in the natural philosophy. As is commonly known, Hume confines, in his works, our speculations to our perceptions, namely in empiricistical terms of impressions, and the copy of them, that is ideas. Therefore his claims place emphasis on the importance of phenomenal resources given to us by experiences. According to him, if we engage in the study of real nature and the operations of the external world beyond experiences, we would encounter serious difficulties. Yet, it seems that hume has not necessarily the faithfulness in his own statements with regard to the empiricistic doctrine. It is here, that I argue, that we can reconsider Hume's empiricism.

I Two empiricisms in moral sciences

Hume's empiricism is interpreted in various way. D.Garrett finely divides the characters of Hume's empiricism into five kinds: such as methodological, conceptual, nomological, explanatory, and reductive empiricism (Garrett: 29-38). I take out two empiricisms from them that correlate with the following arguments.

1. Methodological empiricism

Hume writes in his Treatise that 'the only solid foundation we can give to this science itself must be laid on experience and observation' (THN: xvi). Garrett construes Hume's attitude that 'observation should be the main determinant of theory, and that, in case of apparent conflict, theory should generally be revised to accommodate the interpreted observations, ...' (Garrett: 30). That is to say, the theory of science of man has acceptability of revision in response to observations.

2. Conceptual empiricism

According to Garrett, it is 'the view that the semantic content of thought is always fully derived from things or features of things as they have been encountered in sensory or reflective experience' (Garrett: 33). Hume considers his Copy Principle as that 'all our simple ideas in their first appearance are derive'd from simple impressions, which are correspondent to them, and which they exactly represent' (THN: 9).



I Newton's four rules of scientific reasoning

Newton originally thoughts that 'the basic problem of philosophy seems to be to discover the forces of nature from the phenomena of motions and then to demonstrate the other phenomena from these <u>forces</u>' (MP: 382). What he seeks out is not the ultimate nature of forces of nature, but the forces of nature insofar as nature revealed by experiment and observations. After demonstrating the mathematical principles of the laws and conditions of certain motions, and powers or force in Philosophiae Naturalis Principia Mathematica, he first gives rules for the study of natural philosophy to demonstrate the frame of 'the System of the World' (MP: 794-796).

Rule 1:

No more causes of natural things should be admitted than are both true and sufficient to explain their phenomena.

⇒ We could call it a principle of simplicity. Newton's commentary on the first rule above is that '...more causes are in vain when fewer suffice. For nature is simple and does not indulge in the luxury of superfluous causes'. Hume says in his Treatise that 'we must endeavour to render all our principles as universal as possible, by tracing up our experiments to the utmost, and explaining all effects from the simplest and fewest causes, ...'(THN: 5). (We may associate this understanding with the principle of Ockham's razor.)

Rule 2:

Therefore, the causes assigned to natural effects of the same kind must be, so far as possible, the same.

This rule clearly influences Hume's fourth rule in 'Rules by which to judge of causes and effects': 'The same cause always produces the same effect...' (THN: 116). In An Enquiry concerning the Principles of Morals, Hume refers to this rule as 'Newton's chief rule of philosophizing' (EPM: 98).

Rule 3:

Those gualities of body that cannot be intended and remitted and that belong to all bodies on which experiments can be made should be taken as qualities of all bodies universally.

⇒ This rule was added in the second edition of MP. We could call it a principle of universality. The content of this rule seem to be conveyed as Hume's positive attitude toward inductive reasonings supported by his association of ideas, rather than negative one that falls into skepticism.

Rule 4:

In experimental philosophy, propositions gathered from phenomena by induction should be considered either exactly or very nearly true notwithstanding any contrary hypotheses, until yet other phenomena make such propositions either more exact or liable exceptions.

⇒ This rule was added in the third edition of MP. Newton's commentary on the last rule above is that 'This rule should be followed so that arguments based on induction may not be nullified by hypotheses'. Newton's insight here is that we should give a certain authority to the method of induction, which makes provisional propositions possible. Hume argues the same many times. Newton appears to accept a priori principles in rules 1, 2, and 3. But, as E.A.Burtt said, Newton restricts them by fourth rule within narrower limits (Burtt: 194).

References

THN: Hume, D. [1739-40], A Treatise of Human Nature, edited by D. F. Norton and M. J. Norton: Oxford University Press, 2000.

EHU: Hume, D. [1748], An Enquiry concerning Human Understanding, edited by T. L. Beauchamp: Oxford University Press, 1999. EPM: Hume, D. [1751], An Enquiry concerning the Principles of Morals, edited by T. L. Beauchamp: Oxford University Press, 1998

Burtt, E.A. [1925], The Metaphysical Foundations of Modern Physical Science; A Historical and Critical Essay: BiblioLife, 2009; 邦訳『近代科学の形而上学的基礎』:平凡社、1988年. Garrett, D. [1997], Cognition and Commitment in Hume's Philosophy, Oxford University Press. MP: Newton, I. [1687-1726], The Principia, Mathematical Principles of Natural Philosophy, A New Translation by I. Bernard Cohen and Anne Whitman: University of California Press, 1999. Strawson, G. [2000], "Objects and Power" in The New Hume Debate, edited by R.Read and K.A.Richman: Routledge.

Williams, M. [1996], Unnatural Doubts,; Epistemological Realism and the Basis of Scepticism: Princeton University Press.

II Assumptions of Hume's Arguments

There are several notable arguments in the work of Hume which seem not to be genuinely based on his empiricism, and which are concerned with the existence of body, and the existence of power or force. These arguments, I think, underlie the empiricistical arguments as an assumption or belief. That is to say, they are tacit understandings for Hume. Therefore they need to be received in his studies without verification. And it seems that Hume uses the term 'necessity connection' as determination of the mind, and the term 'power' or 'force' as unknown quality of objects.

1. Arguments concerning the existence of body

(1) I need not observe, that a full knowledge of the object is not requisite, but only of those qualities of it, which we believe to exist (THN: 116).

(2) ...'tis in vain to ask, whether there be body or not? That is a point, which we must take for granted in all our reasonings (THN: 125).

2. Arguments concerning the existence of powers or forces

(3) The operations of nature are independent of our thought and reasoning, I allow it; ... (THN: 113).

(4) The scene of the universe are continually shifting, and one object follows another in an uninterrupted succession; but the power or force, which actuates the whole machine, is entirely concealed from us, and never discovers itself in any of the sensible qualities of body (EHU: 136).

These descriptions are neither grounded on a priori knowledge or principle of the external world, nor are they grounded on his empiricism that would have been influenced by Newton, in that Hume does not properly employs the experimental method of reasoning. They are rather our ordinary belief or assumption, on which we speak and behave in daily life, and with which scientific investigations begin. These thoughts maintain the view of naïve causal realism proposed by G.Strawson that though we cannot perceive the necessary connection between objects, it nevertheless really exists (Strawson: 31-48). And this interpretation is partly due to the view of contextualism suggested by M.Williams that claims relativity between ordinary or scientific statement and skeptical statement. (Williams: 22-31).



Conclusion

· Hume's empiricism is based on ordinary assumptions or beliefs, but they needs not to be verified, because they make all the reasoning possible.

·Hume's experimental method of reasoning is influenced by Newton's fourth rule, in that they acknowledge the authority of the provisional character of induction, in the same manner, as a conductive principle in moral sciences.

•Though Hume believes in the existence of the external objects or real powers, he has to adopt agnosticism concerning them in his empiricism. Because his aim of science of man is to provide a secure and solid foundation for moral sciences.

Footnote

At that time, the term 'experiment' had two meanings: As an 'experiment' in natural science, and as an 'empirical fact'. Hume mainly employed the concept based on its second meaning.

Epistemic Deference and Transmission of Knowledge



NIHEI Mariko (D3, Philosophy Faculty of Arts and Letters, Tohoku University)

Introduction



Does B Know that a quark is a fundamental particle?

We believe in countless expert claims without being able to confirm their truth and in many cases we lack the necessary competence to do so. It is difficult to acquire all relevant evidence concerning particular situations by ourselves because we do not go through extensive specialized training.

Strange as it may seem to the non-philosopher, mainstream epistemologists have paid little attention to the phenomenon of epistemic dependence or intellectual division of labor. The aim of my research is to advance a new epistemological theory to explain this phenomenon and to provide a vivid description of the social-public character of knowledge. For this aim, in this presentation, I focus on the following three works in order to show how the traditional epistemology does not keep in step with the situation of 'epistemic dependence' which is to be a commonplace affair in our time. 1.) I present the conception of 'epistemic dependence' suggested by John Hardwig.

- 2.) According to Hardwig, if we receive 'epistemic dependence' then we must face the dilemmatic choice between two epistemological models . I will point out that both options are incomplete.
- 3.) I clarify that Hardwig falls into this dilemma because he has not perfectly freed himself from the old epistemological obsession and that if we would abandon this obsession the dilemma could be eliminated.

1.) Hardwig's insistence on epistemic dependence In Hardwig [1985]:

- According to traditional epistemology, one can only have rational reasons for believing p, if he has evidence for p; and evidence is anything that establishes the truth of *p*. •But, in our culture, the more is known that is relevant to the truth of one's beliefs the nor anyone is able to know by himself.
- •Therefore, we can never avoid some epistemic dependence on experts.
- Non-expert B can acquire *knowledge* of some **proposition** *p* **from expert A**, even though
- (a) B has not performed the inquiry capable of providing the evidence for p,
- (b) B is not competent to perform that inquiry,
- (c) B is not competent to assess the merits of the evidence provided by A,
- (d) B cannot understand what p means.

•We must say that B's belief is rationally justified and B knows p if we do not want to receive that a very large percentage of beliefs in this complex culture are irrational. According to Hardwig, to accept this conception of epistemic dependence is to deviate from the traditional epistemological view. Instead of the old view, he suggests two options: the epistemic deference model and the epistemic community model.

*This naming is not by Hardwig but by my own.

3.) Tentative conclusion: Drawing the moral from Hardwig's story





• In this model the primary knower is not individual A or B but community of A and B. • p is known not by any person but by the community. Thus, members of group cannot say that 'I know that p', but only

we know that p

Knowledge and its evidence are common property of the community

•By belonging to the epistemic community, each member's belief is justified believing

This model retains the intuition that the knower must possess evidence within himself, but permits a 'community' or 'group' to count as knowers.

* This model represents a kind of naïve communitarian epistemology. See Welbourne[1981] and Kush[2002]

- Hardwig does not define what the criterion for an epistemic community is.
- •Granting that p is common knowledge of communities, this is not to mean expert A and non-expertB have the same information about p.

'We know that p'.

• When something is transmitted as knowledge, A and B are already standing in a normative context. (Ex. teacher-student, doctor-patient, that is, the context of 'epistemic deference'?)

Should we make a decision in favor of either one? Perhaps not. Surely, Hardwig points out the narrow conception of evidence in old individualistic epistemology and presents the new direction towards a socialistic epistemology. Yet he remains too conservative because he endures the old notion of 'knowledge'. He seems to treat knowledge as something non-temporal and to suppose that knowledge itself is invariable through the transmission from one person to another. In other words, supposing that p could be known, that is that p is knowledge, then for all positive knowers (whether individuals or not) who know p posses p while those who do not know p do not posses p. However when we abandon this assumption we need not to faces the dilemma between an all-out deference to experts and a posit community as one whole epistemic subject. (Like Kusch[2002], we can define knowledge as a kind of entitlement and commitment among persons concerning testimony.)Obviously, there is some epistemic deference in our culture. But I think that 'deference' is constructed in dynamic processes of transmission of testimony among members of a community. Therefore, I argue that 'knowledge' represents not something fixed-timelessness but something historical-variable, just like deference.

Bibliography

Brewer [2006]: "Scientific expert testimony and Intellectual due process", in Selinger & Crease(eds.) The philosophy of Expertise, Columbia U.P.

Hardwig [1985]: "Epistemic dependence", in The Journal of Philosophy, Vol.82. No.7. pp.335-349.

Kusch [2002]: Knowledge by Agreement, Oxford U.P. Welborne [1981]: "The community of knowledge", in The Philosophical Quarterly, Vol.31. No.125. pp.302-314.

1. INTRODUCTION

This study explores the meanings of technics for human beings with regards to "reliability" or "durability" that technology creates.

2. HUMAN BEINGS AS "LIFE"

In the first place, organisms fight and eat each other. "Life" (in German: *Leben*) produces the tissues of itself and survives after these die. Human beings live in the same fashion.

 \rightarrow We fight for survival, kill other lives, and produce children who live over the death of their parents. In this process, we cannot help feeling the finiteness of ourselves.

Thus T. Hobbs and G. W. F. Hegel argued that human beings want to overcome this natural process(ex: excessive emotions like fear). For this purpose, we need to find a way for continuous control of nature. Such techniques include politics and science technology (politics and science are fundamentally related to each other in this point).

3. MAGIC AS PRIMITIVE TECHNOLOGY

Parents maintain themselves by producing their children, and the family exists longer than their individual members. According to Hegel, this durability gives human beings a **most primitive and elemental intuition of "eternity"**. (Therefore we tend to favor whatever is stable and continue to exist longer than we do.)

 \rightarrow Thus, magic as the primitive technology took a role of protecting people against the fierceness of nature—and this force was God for them. He was a king who could handle the magic. He provided his people with "durability" of their family.

4. Modern state of affairs

In modern times, in "the twilight of the faith in God" (M. Horkheimer), the power of magic ceases to exist. But is it only the magic that extinguished?

The kingship extinguished as did the conventional family. (Hegel foresaw this fact in the early 18th century.)—"Falling down of all the value" (F. Nietzsche). As J. Habermas argues, that "securalization" and

SUZUKI Ryozo (Tohoku University : Philosophy D3)

"demagicalization (in German: Entzauberung)" are the back of rationalization and only the abstract political power remains without old-fashioned kings.

5. WHAT KIND OF TECHNOLOGY IS THE GREATEST CONCERN TODAY?-AND ITS REASON

From the point of practical significance, technology encompasses a deep relation with the hope for *durability* (speaking metaphysicaly, "eternity"!) — in short, *the continuance of family, the stability of everyday life* against life's changeability or ambiguity. ("*Life against death*" — Norman. O. Brown)

 \rightarrow For that reason, one of the deepest interests in science technology at present concentrates on reproductive technologies. And in this sense, as A. Gehlen has pointed out, the reason for why human beings come to depend on petrochemistry becomes obvious. Because this technology is more *stable* than other natural things.

<u>6. "The eternal irony"</u>

On the other hand, economical changes in modern times demand from us "the quickness to adapt, ability to correctly react to stimuli and specialized skill" (Horkheimer). But we know that the family represents something that we cannot substitute artificially. In this sense, Hegel calls the characteristics of the family "the eternal irony", as long as it is opposed to human work, because **the community and society where the rationality has a meaning depend on members of the family**.

If we cannot go back to 'good old times', then we should reconsider the meaning of technology into which human beings lay their hopes, over and over again.

Hegel. *Phänomenologie des Geistes* : in Werke in zwanzig Bänden, Bd. 3, Frankfurt am Main, 1973.

- M. Horkheimer, Traditionelle und kritische Theorie und andere Aufsätze (Japanische Edition), 1974.
- A. Gehlen, Anthropologische Forschung, Reinbek bei Hambrug, 1961.
- N. O. Brown, Life against death, Wesleyan University Press, 1959.
- J. Harbermas, *Technik und Wissenschaft als>Ideologie* <, Frankfurt am Main, 1968.

^{*} References

What is ethically problematic in Biogenetics?

Takuma OBARA (Philosophy, D3, Tohoku Univ.)

1.序

本発表では、今日の遺伝子工学系科学の発展にともなって起こりうる倫理的な問題 を明確化することを試みている。初めに、その種の科学が今後もたらすであろう事態に ついて言及し、次に、それに対してなされる一般的な倫理的反論について言及する。 最後に、そうした反論によって実は覆い隠されてしまっている問題、我々が真に考察せ ねばならない問題を提示する。

2. 遺伝子工学的発展がもたらす事態と一般的倫理的反論

遺伝子工学にみられるような科学的進展の主な帰結は、一言で言って、「人間性」の 終焉である。ひとたびその構造の規則性が遺伝子レベルで知られてしまえば、人間的 有機体は操作可能な対象と化す。もはや神秘的なところは何もなく、すべては物理法 則に従って計算可能、予測可能な物的存在とみなされる。もちろんこうした考えはそれ ほど新しいものではないが、しかし遺伝子という(これまでのところ)我々のもっとも基礎 的な構成単位のレベルでそうであるという点で、深刻である。それは遺伝子工学が 我々の尊厳の感覚や自律といった「人間本性」に影響を及ぼすことを意味している。 我々がどれほど自由にどれほど自律的に行動した(と思っている)としても、それはあ らかじめ遺伝子にプログラムされていた行動にすぎない、というわけだ。

この場合、人格としてのその人のアイデンティティーの核心が破壊される、という反 論が考えられる。それは同時に教育概念の破壊でもある。というのも、我々は、人は教 養形成を通じて道徳的アイデンティティーを発達させる、と考えているからだ。しかし、 自然科学系の究極の帰結にしたがえば、どれだけ教育し訓練して道徳的人格を形成 しようと、人は遺伝子に即して行動するとみなされる。教育の観念は無意味化する。そ の結果、たとえば、いわゆる更生的訓練(懲役)ではなく、生化学的遺伝子工学的な手 段で犯罪を捜滅するような社会になるだろう。犯罪者には、過剰な攻撃衝動を抑える 薬物の服用が強制され、あるいは、彼の人格から攻撃性そのものを取り除くような遺 伝子操作が施されるだろう。彼は文字通り人格を矯正されるわけである。

3. 旧来の倫理への固執は倫理的ではない

しかしだからといって、根本的な人間性を守るために遺伝子操作のような技術を制 限しなければならないという態度をとるのは誤っている。それは、知っているにもかか わらずあたかも知らないふりをするという、物神的な分裂態度である。すなわち、「科学 が主張していることを私は十分によく知っているが、にもかかわらず、人間性(という見 かけ)を保持するために、私はそれを無視し、あたかも自分がそれを知らないかのよう に振る舞う」という態度である。こうした態度は真の問題、つまり、自由や自律、尊厳と いった人間本性が新たな科学的条件によってどう変化するのか、という真の問題に直 面することを妨げてしまう。

かつてデカルトはこうした旧来の倫理観(に固執する態度)を「一時的倫理」と呼んだ。 なぜ「一時的」かと言えば、今や新たな人間観の地平が開かれているからである。新た な地平へ乗り出すことはこれまでの土台を失うことであり、支えなしに自立することであ る。それは一つの危機である。それを回避するためにも、これまでの慣習やならわしに 今一度そしてこれを最期に身を浸す必要がある。しかし、それはあくまで新たな地盤を 築くためであって、旧来の土台を守るためではない。一度知ってしまった以上、もはや 純真無垢な無知へ帰ることはできない。今やこの新たな現実のもとにとどまり、回避し たり隠蔽したりすることなく、それを引き受けなければならない。これまでの倫理に固執 すること、それは実際には倫理的態度ではない。新たな地平で新たな倫理を構築しよ うとすること、これこそが真に倫理的な態度なのである。

4. 遺伝子工学系科学の発展において、何が真に倫理的な問題であるのか?

先に挙げた「矯正」の例を考えてみよう。仮に罪を犯したある人間が、懲役としての 数年の更生訓練の代わりに、衝動を抑制する薬品を常用することを条件にすぐさま社 会復帰したとしよう。おそらく人々は容易には彼を受け入れられないだろう。しかもそれ は更生訓練を経て社会復帰した人より以上にである。なぜか。それは、彼が自律的で はなく他律的に衝動を抑えていると人は考えるからである。言い換えれば、薬物は「外 的」操作とみなされ、その人固有の特性とは認められえない。人は人間の内面性に価 値を置き、それゆえ、人間の内面を形成するものとしての教育や訓練を評価してきた のである。したがって、彼は本質的には(薬なしには)未だ攻撃的人間とみなされる。

この話をもう一歩進めてみよう。今述べたような衝動を抑える薬ならばまだ容易に外 的と判断することができる。では、そうではなくてたしかに更生訓練によって自律的に 衝動を抑えられるようになったが、実はその訓練において苦闘に耐えることを可能する 薬を使っていたとしたらどうか。外的操作は確かに認められるが、当人が苦しい経験を したこともまた確かである。さらには、上でも述べたように、彼の人格そのものを矯正し たとしたらどうか。要するに、どこまでが個人に属するもので、どこからが外的操作の 範囲になるのか。人間における内的なものと外的なものの境界のこうした曖昧化が、 遺伝子工学のような科学の発展によってもたらされる事態である。こうした条件に立っ て、そもそも人間性とはいったい何なのかを考察すること、それが真に哲学的な問いで あろう。こうした問題こそが真に倫理的な問題として考察され、解答されなければなら ない問題なのである。

1. Introduction

In this presentation, I will try to precise the ethical problems that may be posed with particular scientific breakthroughs such as success in decoding the human genome. First, I will elaborate on what such scientific advances can bring about. Second, I am going to mention general, ethical oppositions against it. Finally, I will propose problems that this opposition unconsciously hides, and that we should really consider.

2. Situations that progress in biogenetics can bring about

and general, ethical oppositions against them

The main consequence of the scientific breakthroughs in biogenetics is the end of "humanity". Once we know the rules of its construction at the level of the genome, human organisms are transformed into objects amenable to manipulation. Human beings are not mysterious and they are regarded as material objects which act by the laws of physics. Of course, this idea is not new, but it is serious in the sense that it is true at the level of genomes which are our most basic components. This means that our behavior is only a result of the program of genomes even if (we believe) we act freely or autonomically.

Against this, opposition has been expressing arguing that the very heart of one's identity as a person could be destructed in such cases. Moreover, the concept of education could also be destructed given the fact that we think that we develop our moral identity through education. However, according to ultimate consequences which natural sciences has, human beings act by the rules of genomes no matter how they educate themselves and form their moral identity. The notion of education may be rendered meaningless. As a result, for example, there would be a society that has intention to fight crime through direct biochemical or biogenetic intervention instead of imprisonment at hard labor. The criminals would be compelled to take medications against excessive aggression or be administered to biogenetic manipulations to remove their aggression itself from their personalities. Literally, they would be reformed.

3. Persisting in old ethics is not ethical

However, it is false to assert that technology such as biogenetic manipulations should be restricted for preserving the fundamental humanity. It is a fetishistic splitting attitude: you ignore it while you know it. That is to say, "I know very well what science claims, but, nonetheless, in order to retain [the appearance of] humanity, I choose to ignore it and act as if I don't know it". This attitude prevents us from confronting the true question: 'How do these new scientific conditions transform and reinvent the very human nature such as freedom, autonomy, and dignity?'

Descartes called (the persistence of) the old ethics the "provisional ethic". The reason why it is called "provisional" is that a new horizon is opening now. To take a step toward a new field is to lose our previous foundation and to stand on our own legs without any support. It represents a crisis. In order to evade it, it is necessary to put ourselves into our previous custom or convention again and for the last time. However, it is not to protect the old foundation but to construct the new one. Once we know it, there is no return to innocent ignorance. We have to tarry in the new reality and accept it rather than to evade or conceal it. Persisting old ethics is in fact not ethical attitude. It is attempting to constitute new ethics on the new field that is truly ethical attitude.

4. What is ethically problematic in biogenetic breakthroughs?

Let us consider the above mentioned example of "reformation". What would happen if a criminal immediately returned to his society on condition that he takes medications against his own aggression, instead of receiving a sentence of imprisonment? Certainly, people would consider it to be difficult to accept him much more than a person that receives a sentence of imprisonment. Why? That is because people would think he can control his own aggression not autonomically but heteronomously. In other words, medications are thought to be "external" manipulation, not his proper personality. We have valued the human's interior or inherence and thus estimated education or training as forming it. In short, he would still be regarded as a dangerous person essentially (without medicine).

Finally, let us go one more step further in this example. It is easy to judge drugs that we have just mentioned as external. Now, what would happen if he came to be able to autonomically control his own compulsion through intensified training, whereas he had taken medicine to enable him to endure struggle or agony during the training? Moreover, as I have mentioned above, what would happen if his personality itself was reformed? In a word, can we draw a line of division between the internal and the external within an individual? It is this blurring of the borderline that scientific advances such as biogenetics may bring about. It seems to be an authentic question of philosophy to consider what humanity is under that condition at all. This is the very problem that should be reflected and responded as being the truly ethical one.

The Ontological Genesis of the Theoretical Attitude

Tetsurou YAMASHITA (Tohoku University)

If we conceive scientific activities radically, that is in terms of the ontological structure of the human being, we may understand that science is not the primary way of our being but a modification of practical activities in everydayness. In this presentation I will sketch the process of the ontological genesis of the theoretical attitude. In doing so, I am dependent heavily on *Martin Heidegger*'s ontological insights regarding the human being.

I. The Structure of the Understanding of Being



Human beings exist, understanding being of beings (entities). Understanding has the structure in which we project beings toward their being (horizon).

\underline{II} . The Way in which We are Concerned with Beings



III. Practical Attitude

Availableness and Discovery of the Available Beings



Those properties can, however, be discovered only in a definite practical context; that is, the available beings don't beforehand have properties without practical context, such as physical weight or spatial position etc. The properties of the available beings are regulated and structured according to their context. They don't have the same degree of clarity; which aspect should be centered or put on the periphery, or which property in the centered aspect should come to the fore, and so on, are all determined according to practical attitude. Thus, in practices "subject" is extremely involved in

practical context; subject should behave and take the



foming the connection of purposes. Subject is involved in it, behaving along it on the one hand and functioning as the end of it on the other.

2. Formation of Theoretical Attitude

viewpoint as the context demands.

"Moving-boundaries-away" invalidates practical context; thus, all the viewpoint (aspects) which have already been built in this context are released, without being regulated according to practical context. From these viewpoints each, we can divide the wholeness of objects into the various fields of sciences (space, matter, life, history and so on). Theoretical attitude is formed through this operation.



Each field of science is understood from a single viewpoint (monistically). Because subject is free from involvement in practical context, subject can take an optional viewpoint and discover optional properties which have already been understood, but not clearly. By fixing a viewpoint, we can articulate the field more clearly and minutely than we can expect in everyday practices.

For example, mathematical physics finds its field by discovering in advance the consistently existing beings(matter) and then noticing previously their constituent moments specifiable quantitatively (ex. motion, force, position, time). This operation is called *the mathematical projection of nature*. Discovery of facts, organization of concepts, regulation of methods including experiments, decision of the form of arguments, and so on, are all possible under the light of this projection.

Tools occupy their own proper places, related with other specific tools in each practical context (a connection of purposes). I call this way of being in which tools exist (extremely) context-dependently, "availableness" (readiness-to-hand).

In practices, we can discover unthematically various properties of the available beings in terms of our purpose in each context. The discovered properties are pluralistic according to their context.

IV. Change into Theoretical Attitude

1. Change of Understanding of Being

Theoretical attitude is formed through the change from availableness into *"forwardness-to-hand"*, the way of being in which beings exist constantly, independent of subject in practical context. Forwardness-to-hand is formed through the release of subject from the involvement in practical context.



This change into forwardnessto-hand contains "movingboundaries-away" ("Entschränkung" in German). This operation removes the boundaries between the proper places of available beings, and changes various places(environmental) into uniform positions (spatial). Through this change, beings lose their practical context

(or availableness) and appear as the wholeness of independent objects. Subject is independent of objects; thus he/she can take an optional viewpoint voluntarily.

Conclusion and Suggestion

Because of its projection toward a specific viewpoint, theoretical attitude can discribe (a part of) nature in such clarity and preciseness that practical attitude cannot. However, we cannot constitute the genuine wholeness of nature by combining these right but partial pictures of nature. Theoretical pictures of nature are only results of the division of a specific viewpoint. Nature as it is, the source from which we can derive various viewpoints, always exceed theoretical discovery in abundance. The substance of nature is the possibility that we can derive various theory from it. So science cannot access being of nature in this sense. This is by no means a fault of science, but the price for its clarity and preciseness. From this standpoint, it is necessary to consider the validity of the conception of scientific comprehensive picture of nature, and to clarify ontological significance of scientific trurh, that is, significance of non-ontological truth of science.



Tohoku University Scienceweb GCOE "Weaving Science Web beyond Particle-matter Hierarchy"





Risk, uncertainty and the precautionary principle: How to deal

with scientific uncertainty?

FUJIO Yasuhiko Graduate School of Arts and Letters, Tohoku University

Introduction

The precautionary principle(PP) requires that we should adopt approaches such as regulating or Risk deriving from the development of science and technology

We are in need of risk analysis and management in order to coexist with risks, whose approaches are based on the theory of probability or statistics.

Scientific uncertainty and risk management

We must assess the occurrence of undesirable events probabilistically or statistically. There are some realms of science in which we cannot use such approaches because of a lack of sufficient knowledge with regard to complex phenomena (ex. global warming).

Uncertainty in risk management Such procedures as RCBA(risk-cost-benefit analysis)embody

Emergence and development of the concept of risk

We cannot obtain complete knowledge about future events; we can only anticipate future events probabilistically



The magnitude of risk is calculated as the product of probability of the occurrence of undesirable events and severity of damage which they might pose.

The magnitude of risk

or statistically.



Accordingly, we must make decisions concerning issues of risk under uncertainty of the occurrence of undesirable damage in the future.

Precautionary principle

The origin of the conception of PP

The cholera epidemic in London in 1854: John Snow, the founder of epidemiology, recommended removing the handle of the water pump in order to stop the cholera epidemic. (At that time, it was not proved that cholera was transmitted by polluted water.)

This case may suggest that the conception of PP was known and applied earlier before PP has emerged in the 1970s.





The emergence of PP

PP has been adopted in the legislation for protection of human health and the environment. The most prominent and frequently cited version of PP is probably the 1992 Rio Declaration:

> "In order to protect the environment, the precautionary approach shall be widely applied by States according to their capabilities. Where there are threats of serious or irreversible damage, lack of full scientific certainty shall not be used as a reason for postponing cost-effective measures to prevent environmental degradation."

Conclusions

- We are forced to deal with risk under uncertainty, for we can only recognize the occurrence of undesirable events in the future probabilistically or statistically.
- · We must perform risk management on the ground that there might be scientific uncertainty.
- PP does not disregard the role of science, but acknowledges the limitations of science and the confidence in science.
- We should elucidate the ethical implication of PP and make PP more available in processes of policy-making so that moral theories can address issues of risk.

Problematic of risk management

uncertainties; models used in RCBA are not always perfectly sophisticated models of real events.

In principle, we never entirely know what happened, happens and will happen in the real world, which makes us more modest about our actions that might pose harm to the human health or the environment.

uncertainty.

Example of failed risk management: BSE in UK

The government did not acknowledge the

infection with BSE to humans at first. The Southwood Committee's Report recognized

the slim possibility of humane infection with BSE.

However, the government paid no attention to it, but overestimated the safety. Afterwards the government acknowledged the infection with BSE and the failure of risk management.

Negligence or underestimation of uncertainty of scientific knowledge or evaluation of risk might lead to the failure of risk management.

One of the most remarkable characteristics of PP is that even if there exist no full scientific evidences it is possible and legitimate to regulate or prohibit products or procedures which might pose serious and irreversible damage to human health or the environment.

Criticism against PP

Opponents of PP characterize the defect of PP as:

- 1) The definition of PP is too vague to serve as a regulatory standard,
- 2) PP forces decision-makers to pay unreasonable attention to extremely unlikely scenarios,
- 3) PP would lead to more risk-taking or another risk,
- 4) PP is a value judgment, not a scientific one,
- 5) PP does not take science seriously and marginalizes the role of science in decision-making.

Objections against criticism of PP

Proponents of PP argue against each criticism:

- 1') Vagueness of PP does not necessarily imply uselessness in practice; PP may be given more precise formulations through elaboration and practice,
- 2') PP is not willing to prohibit all actions which may pose harm; all actions may have unforeseen, more or less, harmful consequences, so it is impossible and unreasonable to ban all actions.

3') The way of framing of the decision problem to which PP is applied might cause more risks: PP itself does not pose more risk.

Reference

- 大竹千代子、東賢一『予防原則 人と環境の保護のための基本理念』合同出版、2005年 Commission of the European Communities, Communication from the Commission on the
- Precautionary Principle, http://ec.europa.eu/environment/docum/20001 en.htm, 2000.
- Hansson, S.O., "Risk and ethics: three approaches", In Risk: Philosophical Perspectives, ed. Lewens, T., Routledge.2007. Jensen, K. K., "The Moral Foundation of the Precautionary Principle", Journal of Agricultural and
- Environmental Ethics, 15(1), 2002, p.39-55.
- Sandin, P., "Common-sense precaution and varieties of precautionary principle", In Risk: Philosophical Perspectives, ed. Lewens, T., Routledge, 2007.
- Sandin, P. et al., "Five charges against the precautionary principle", Journal of Risk Research,5(4),2002,p.287-99.
- Sunstein, C., Laws of Fear: Beyond the Precautionary Principle, Cambridge University Press, 2005.

The limits of our knowledge and science

As science is the process for exploring veiled realms, in cutting-edge research the scientific knowledge is incessantly renewed by new scientific discoveries.

There are, however, some cases where we cannot elucidate the causal mechanism of complex phenomena at the present time due to the complexity of phenomena(ex. global warming)

Thus, scientific uncertainty is due to the plasticity or limits of our knowledge.

banning the use of chemical substances or new technology in order to protect the human health and

uncertainty, concerning the severe and irresistible harm to the human health and the environment.

phenomena, and PP as a reasonable principle for decision making confronted with such

PP is thought of as the principle which is to be applied in the issues of risk, under scientific

I wish to examine the notion of scientific uncertainty with regard to complex



Attitudes towards uncertainty in science

In regulating or prohibiting a certain product or procedure, it seems to be necessary to prove the causal relationship between products and the possible harmful effect.

As we know, since scientific knowledge might be renewed, we must not have too much confidence in science.

Therefore, we, especially policy-makers, must keep in mind that science has limitations.

4') All decision rules including PP are value-based, 5') PP is not based on science, but does not contradict science.

Outlook for PP

Development of PP for better policy-making

In the Communication from the Commission on the Precautionary Principle(2000) the Commission of the European Communities noted that measures based on the PP:

·based on an examination of potential costs and benefits of action or lack of action,

·subject to review, in the light of new scientific data, ·capable of assigning responsibility for producing the scientific evidence necessary for a more comprehensive risk assessment.

The foundation of PP in philosophy and ethics

Philosophy might be able to argue 'uncertainty' in risk management in the light of philosophy of probability or statistics.

"Our current moral theories are not suitable to deal with issues of risk." (Hanson, 2007)

We should, therefore, explicate the ethical implication of PP.



Precautionary principle: an approach to risks

the environment despite the lack of sufficient scientific certainty.



The mechanism of suppressed dynamical friction in a constant density core of dwarf galaxies

Inoue Shigeki (Tohoku Univ.)

Abstract : Dynamical friction problem is a long-standing dilemma. In dwarf galaxies, dynamical friction on their globular clusters (GCs) is too strong to keep their orbital motions. Nonetheless, GCs do exist even in current dwarfs. However, a solution have been proposed. If dwarf galaxies have a cored dark matter halo which has a constant density region in its center, the dynamical friction is significantly weakened. But, the mechanism of the suppressed dynamical friction has not been clarified yet.

By means of N-body simulation, I find that the mechanism of the suppressed dynamical friction is the effect of orbital resonance between the GC and the halo.







A favorable answer to "the cusp-core problem" : In Fig.1, I show a result of my simulation. As mentioned above, the GC can survive from the strong dynamical friction in the cored structure. This

0.1

result impllies that the cored halo is a favorable dark matter structure for dwarf galaxies on the point of the existence of the GCs in dwarfs.

Fig.2 $\Delta E/\Phi_0$ 0.05 Energy change -0.05 -0.1 -0.1 0.3 0.1

The mechanism of the suppressed dynamical friction : I examine orbital energies and changes of the energies of halo particles $(E-\Delta E \text{ diagram}, Fig. 2)$. I find that the energies of some particles largely increase or decrease (the green & red squares on Fig.2) and these particles have prograde rotations with the GC. These particles resonate with the GC and have very complex orbits. In

Fig.3, I show a example of the orbits in the rotating coodinate of the GC.

I extract such resonant particles and calculate energy transfers from the resnant particles to the GC. I find that the resonant particles inject energy to the GC (Fig.4).

This result means that the dynamical friction on the GC is canceled out by the resonant particles in the core structure.





t [Gyr]

10